
Transportation Network Analysis

Volume I: Static and Dynamic Traffic Assignment

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Preface

This book is the product of ten years of teaching transportation network analysis, at the The University of Texas at Austin, the University of Wyoming, the University of Connecticut, and Portland State University. This version is being released as a *public beta*, with a final first edition hopefully to be released within a year. In particular, we aim to improve the quality of the figures, add some additional examples and exercises, and add historical and bibliographic notes to each chapter. A second volume, covering transit, freight, and logistics, is also under preparation.

Any help you can offer to improve this text would be greatly appreciated, whether spotting typos, math or logic errors, inconsistent terminology, or any other suggestions about how the content can be better explained, better organized, or better presented.

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Difficulty Scale for Exercises

Inspired by Donald Knuth's *The Art of Computer Programming*, the exercises are marked with estimates of their difficulty. The key reference points are:

- 0** : A nearly trivial problem that you should be able to answer without any pencil-and-paper work.
- 20** : A straightforward problem that may require a few minutes of effort, but nothing too difficult if you have given the chapter a good read.
- 40** : A typical problem requiring some thought, or a few attempts at solution in different ways, but the answer should yield after some dedicated effort.
- 60** : A problem of above-average difficulty, where the correct approach is not obvious. You may require a bit of scratch paper or computer work as you try out different approaches before settling on a solution.
- 80** : A highly challenging or involved problem, which may be appropriate as a course project or other long-term study.
- 100** : An open problem in the research literature, whose solution would be a substantial advance in the understanding of transportation networks.

The tens digit indicates the intellectual difficulty of the exercise, while the ones digit indicates the amount of calculation required. An exercise rated at 50 may require more cleverness and insight than one ranked at 49, but the ultimate solution is shorter. Of course, each student will find different problems more challenging than others.

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Part I

Preliminaries

Chapter 1

Introduction to Transportation Networks

This introductory chapter lays the groundwork for traffic assignment, providing some overall context for transportation planning in Section 1.1. Some examples of networks in transportation are given in Section 1.2. The key idea in traffic assignment is the notion of equilibrium, which is presented in Section 1.3. The goals of traffic assignment are described in Section 1.4. Traffic assignment models can be broadly classified as static or dynamic. Both types of models are described in this book, and Sections 1.5 and 1.6 provide general perspective on these types of models. Prerequisites for understanding this material are briefly discussed in Section 1.7.

1.1 Transportation Networks

Planning helps ensure that transportation spending and policies are as effective as possible. As transportation engineers and researchers, we support this process by developing and running models which predict the impact of potential projects or policies — for instance, what would be the impact on city traffic and emissions if an extra lane was added on a major freeway? If the toll on a bridge was reduced? If streetcar lines are installed downtown? In this way, the benefits of projects can be compared with their costs, and funding and implementation priorities established. Depending on the models used, a variety of measures of effectiveness can be considered, and one may want to know the impacts of a project on mobility, congestion, emissions, equity, toll revenue, transit ridership, infrastructure maintenance needs, or countless other metrics.

This is rather difficult. To predict ridership on a new transit line with complete accuracy would require knowing how many trips every single possible rider makes, and the decision process each one of these potential riders uses when deciding whether or not to use transit. Unlike trusses or beams, human beings can behave in ways that are impossible to predict, maddeningly inconsistent,

and motivated by a variety of factors difficult to observe (many of which occur at a subconscious level). Further, most transportation infrastructure lasts for decades, meaning that effective planning must also predict the impact of projects and policies decades into the future. And so far, we've only considered the pure engineering dimension: introduce local, state, and federal politics into the mix, other stakeholders such as neighborhood associations and transit agencies, and a public with a variety of priorities (is it more important to reduce congestion, increase safety, or have livable communities?), and the picture only grows more complicated. What to do?

Enter the mathematical model. The purpose of a mathematical model is to translate a complicated, but important, real-world problem into precise, quantitative language that can be clearly and unambiguously analyzed. By their nature, models cannot account for all of the possible factors influencing planning. As the famous statistician George Box once quipped, "All models are wrong, but some models are useful." A useful model is one which provides enough insight that good decisions can be made. To do this, a model must capture the most important characteristics of the underlying system; must not require more input data than what is available for calibration; and must not require more time and memory than what available hardware permits.

Further, just because a model is useful does not mean it cannot be improved. Indeed, this is the goal of transportation researchers around the world. The usual pattern is to start with a model which is simple, transparent, insightful... and also wrong. This simple model can then be improved in ways to make it more correct and useful, and this is the general pattern which will be seen in this book. The first network models you will see are such gross simplifications of reality that you may question whether they can truly be of value in practice. Perhaps they can, perhaps they can't; but in any case, they form a foundation for more advanced and realistic models which relax the assumptions made earlier on.

For this reason, as a student of transportation planning, you should always be looking for the assumptions involved in everything you see. All models make assumptions which are not entirely correct. The relevant questions are, how much does this assumption limit the applicability of the model, and how easy would it be to relax this assumption? If you're looking for a research topic, finding an existing model and relaxing an assumption is often a good approach. With this book, if you clearly understand all of the assumptions underlying each model, and how they differ from those made in other models introduced, you're 90% of the way there.

Networks are a type of mathematical model which are very frequently used in the study of transportation planning. This introductory chapter gives a very brief overview of transportation networks, and provides a sketch for the remainder of the book.

This book covers both *static* and *dynamic* network models. Static models assume that network conditions are at steady-state, while dynamic models represent changes in congestion and demand patterns over the course of several hours or a day. Static models were the first to be developed historically, and

remain the most commonly-used in current transportation planning practice. Dynamic models are more realistic in portraying congestion, but require more data for calibration and validation, and more computational resources to run. Solving and interpreting the output of dynamic models is also more difficult. As research progresses, however, more planners are using dynamic models, particularly for applications when travel conditions are changing rapidly during the analysis period. This chapter will present a balanced perspective of the advantages and disadvantages of dynamic traffic assignment *vis-à-vis* static assignment, but one of them is worth mentioning now: dynamic traffic assignment models are inherently mode-specific.

That is, unlike in static assignment (where it is relatively easy to build “multimodal” networks mixing roadway, transit, air, and waterway infrastructure), the vast majority of dynamic traffic assignment models have been specifically tailored to modeling vehicle congestion on roadways. In recent years, researchers have started integrating other modes into dynamic traffic assignment, and this area is likely to receive more attention in years to come. However, the congestion model for each mode must be custom-built. This is at once an advantage (in that congestion in different modes arises from fundamentally different sources, and perhaps ought to be modeled quite differently) and a disadvantage (a “generic” dynamic traffic assignment model can only be specified at a very high level). For this reason, this book will focus specifically on automobile traffic on roadway networks. This is not meant to suggest that dynamic traffic assignment cannot or should not be applied to other modes, but simply an admission that covering other modes would essentially require re-learning a new theory for each mode. Developing such theories would make excellent research topics.

1.2 Examples of Networks

Networks are fundamental to the study of large-scale transportation models representing an entire metropolitan area, a state, or multistate regions. They can be applied in many contexts, including alternatives analysis, developing congestion pricing plans, identifying bottlenecks and critical infrastructure, shipping and freight logistics, multimodal planning, and disaster evacuation planning, to name only a few. The reason network models are so useful, and so broadly applicable, is because a mathematical network is a simple, compact, and flexible way to represent a large, complicated system.

A network consists of *links* and *nodes*. In transportation applications, a link usually represents a means of travel from one point to another: a road segment between two intersections, a bus route between two stops, and so on, as seen in Figure 1.1. The nodes, in turn, are the endpoints of the links. Quite often, nodes are adjacent to multiple links, so a node representing an intersection may adjoin multiple links representing road segments. Nodes and links may also be more abstract; for instance, links in a multimodal network might represent a transfer from one transport mode to another. The level of detail in a network varies from application to application. For multistate freight models, major

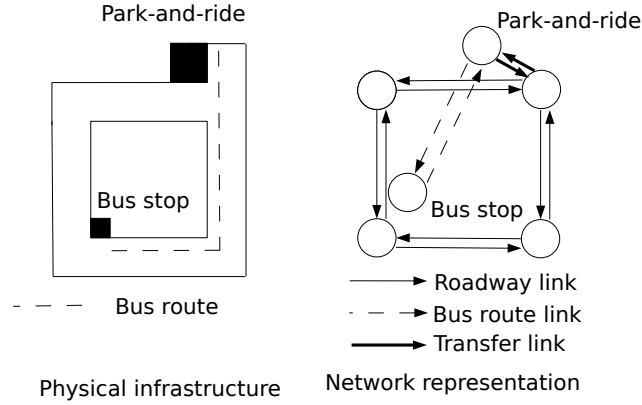


Figure 1.1: Nodes and links in transportation networks.

Table 1.1: Nodes and links in different kinds of transportation networks.

Network type	Nodes	Links
Roadway	Intersections	Street segments
Public transit	Bus or train stops	Route segments
Freight	Factories, warehouses, retailers	Shipping options
Air	Airports	Flights
Maritime	Ports	Shipping channels

highways may be the only links, and major cities the only nodes. For a city's planning model, all major and minor arterials may be included as well. For a more detailed model, individual intersections may be "exploded" so that different links represent each turning movement (Figure 1.2). Other examples of transportation networks are shown in Table 1.1.

1.3 The Notion of Equilibrium

The nature of transportation systems is that of multiple interacting systems. Congestion is determined by the choices travelers make: where, when, how often to travel, and by what mode. At the same time, these choices depend on congestion: travelers may choose routes or departure times to avoid congestion. These two "systems" (travel choices and system congestion) are thus interdependent and interrelated, with a circular or chicken-and-egg quality to their relationship. This interdependency lies at the root of transportation analysis. It is at once interesting, because of the complexity of transportation systems involving both humans and physical systems; challenging, because we must find a way to resolve this circular dependency; and frustrating, because obvious-

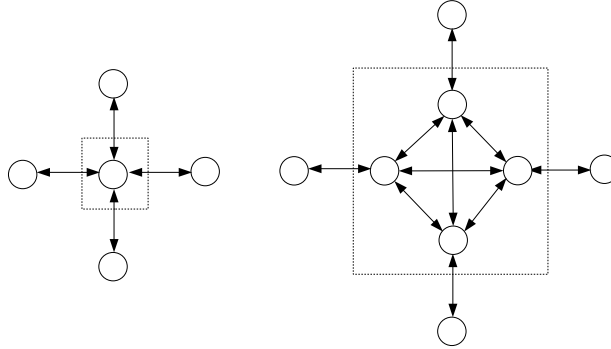


Figure 1.2: Two representations of the same intersection.

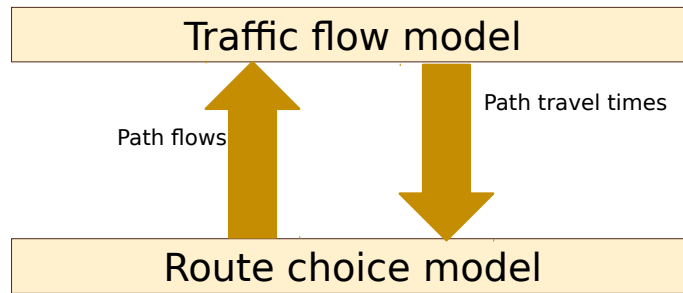


Figure 1.3: Traffic assignment as interacting systems.

looking policy interventions can actually be counterproductive. Some examples of “paradoxical” effects will be seen in Chapters 5 and 12.

The schematic in Figure 1.3 shows the dependency between the choices made by travelers (sometimes called the *demand side*), and the congestion and delay in the system (sometimes called the *supply side*) in the basic traffic assignment problem. Each of these systems requires a distinct set of models. Demand-side models should be behavioral in nature, identifying what factors influence travel choices, and how. Supply-side models are often based in traffic flow theory, queueing theory, computer simulation, or empirical formulas describing how congestion will form.

It is not difficult to imagine other types of mutually-dependent transportation systems. Figure 1.4 shows how one might model traffic assignment in a region with a privately-operated toll road. Now, there are three systems. In addition to the demand side and supply side from before, the private toll operator can also influence the state of the network by choosing the toll in some way, such as maximizing toll revenue. But this choice is not made in isolation: as the toll

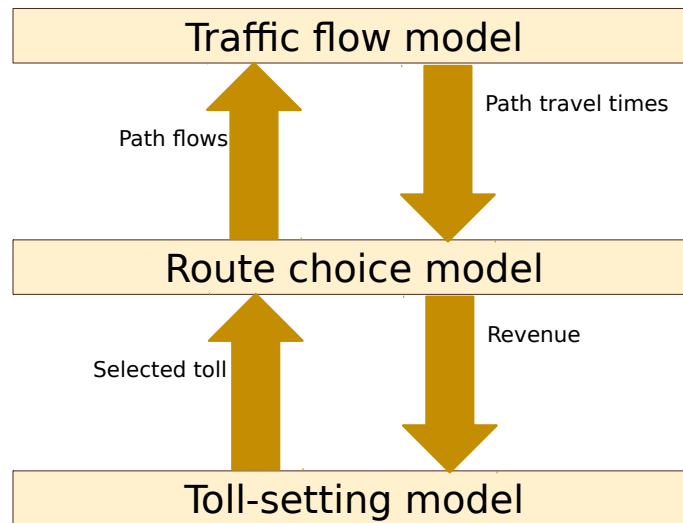


Figure 1.4: A more complicated traffic assignment problem, with tolls.

is increased, drivers will choose alternate routes, suggesting that driver choices are affected by tolls just as the toll revenue is determined by driver choices. It is fruitful to think of other ways this type of system can be expanded. For instance, a government agency might set regulations on the maximum and minimum toll values, but travelers can influence these policy decisions through the voting process.

The task of transportation planners is to somehow make useful predictions to assist with policy decision and alternatives analysis, despite the complexities which arise when mutually-dependent systems interact. The key idea is that a good prediction is *mutually consistent* in the sense that all of the systems should “agree” with the prediction. As an example, in the basic traffic assignment problem (Figure 1.3), a planning model will provide both a forecast of travel choices, and a forecast of system congestion. These should be consistent in the sense that inputting the forecasted travel choices into the supply-side model should give the forecast of system congestion, and inputting the forecasted system congestion into the demand-side model should give the forecast of travel choices. Such a consistent solution is termed an *equilibrium*.

The word equilibrium is meant to allude to the concept of economic equilibrium, as it is used in game theory. In game theory, several agents each choose a particular action, and depending on the choices of all of the agents, each receives a payoff (perhaps monetary, or simply in terms of happiness or satisfaction). Each agent wants to maximize their payoff. The objective is to find a “consistent” or equilibrium solution, in which all of the agents are choosing actions which maximize their payoff (keeping in mind that an agent cannot control

Table 1.2: Alice and Bob's game; Alice chooses the row and Bob the column.

		Bob	
		Cactus Café	Desert Drafthouse
Alice	Cactus Café	$(-1, -1)$	$(1, 1)$
	Desert Drafthouse	$(1, 1)$	$(-1, -1)$

another agent's decision). A few examples are in order.

Consider first a game with two players (call them Alice and Bob), who happen to live in a small town with only two bars (the Cactus Café and the Desert Drafthouse). Alice and Bob have recently broken off their relationship, so they each want to go out to a bar. If they attend different bars, both of them will be happy (signified by a payoff of $+1$), but if they attend the same bar an awkward situation will arise and they will regret having gone out at all (signified by a payoff of -1). Table 1.2 shows the four possible situations which can arise — each cell in the table lists Alice's payoff first, followed by Bob's. Two of these are boldfaced, indicating that they are equilibrium solutions: if Alice is at the Cactus Café and Bob at the Desert Drafthouse (or vice versa), they each receive a payoff of $+1$, which is the best they could hope to receive given what the other is doing. The states where they attend the same bar are not equilibria; either of them would be better off switching to the other bar. This is a game with two equilibria in which Alice and Bob always attend the same bar each week.¹

A second game involves the tale of Erica and Fred, two criminals who have engaged in a decade-long spree of major art thefts. They are finally apprehended by the police, but for a minor crime of shoplifting a candy bar from the grocery store. The police suspect the pair of the more serious crimes, but have no hard evidence. So, they place Erica and Fred in separate jail cells. They approach Erica, offering her a reduced sentence in exchange for testifying against Fred for the art thefts, and separately approach Fred, offering him a reduced sentence if he would testify against Erica. If they remain loyal to each other, they will be convicted only of shoplifting and will each spend a year in jail. If Erica testifies against Fred, but Fred stays silent, then Fred goes to jail for 15 years while Erica gets off free. (The same is true in reverse if Fred testifies against Erica.) If they both testify against each other, they will both be convicted of the major art theft, but will have a slightly reduced jail term of 14 years for being cooperative. This game is diagrammed in Table 1.3, where the "payoff" is the negative of the number of years spent in jail, negative because more years in jail represents a worse outcome. Surprisingly, the only equilibrium solution is for both of them to testify against each other. From Erica's perspective,

¹There is also a third equilibrium in which they each randomly choose a bar each weekend, but equilibria involving randomization are outside the scope of this book.

Table 1.3: Erica and Fred's game.

		Fred	
		Testify	Remain silent
Erica	Testify	$(-14, -14)$	$(0, -15)$
	Remain silent	$(-15, 0)$	$(-1, -1)$

Table 1.4: Ginger and Harold's game.

		Harold	
		Heads	Tails
Ginger	Heads	$(+1, -1)$	$(-1, +1)$
	Tails	$(-1, +1)$	$(+1, -1)$

she is better off testifying against Fred *no matter what Fred will do*. If he is going to testify against her, she can reduce her sentence from 15 years to 14 by testifying against Fred. If he is going to stay silent, she can reduce her sentence from one year to zero by testifying against him. Fred's logic is exactly the same. This seemingly-paradoxical result, known as the *prisoner's dilemma*, shows that agents maximizing their own payoff can actually end up in a very bad situation when you look at their combined payoffs!

A third game, far less dramatic than the first two, involves Ginger and Harold, who are retirees passing the time by playing a simple game. Each of them has a penny, and on a count of three each of them chooses to reveal either the head or the tail of their penny. If the pennies show the same (both heads or both tails), Ginger keeps them both. If one penny shows heads and the other shows tails, Harold keeps them both. (Table 1.4). In this case, there is *no* equilibrium solution: if Ginger always shows heads, Harold will learn and always show tails; once Ginger realizes this, she will start showing tails, and so on *ad infinitum*.

You may be wondering how these games are relevant to transportation problems. In fact, the route choice decision can be seen as a game with a very large number of players. Some drivers may choose to avoid the freeway, anticipating a certain level of congestion and trying to second-guess what others are doing — but surely other drivers are engaging in the exact same process.² Each of these three games has bearing on the traffic assignment problem. The game with Alice and Bob shows that some games have more than one equilibrium solution

²To borrow from Yogi Berra, nobody takes the freeway during rush hour anymore — it's too congested.

(an issue of equilibrium *uniqueness*), although the two equilibrium solutions are the same if you just look at the total number of people at each bar and not which individuals go at each. What does it mean for transportation planning if a model can give several seemingly valid predictions? The game with Erica and Fred shows that agents individually doing what is best for themselves may lead to an outcome which is quite bad overall, an issue of equilibrium *efficiency*. As we will see later on, in transportation systems this opens the door for seemingly helpful projects (like capacity expansion on congested roads) to actually make things worse. The game with Ginger and Harold is a case where there is no equilibrium at all (an issue of equilibrium *existence*). If this could happen in a transportation planning model, then perhaps equilibrium is the wrong concept to use altogether. These questions of uniqueness, efficiency, and existence are important, and will appear throughout the book.

The three example games described above can be analyzed directly, by enumerating all the possible outcomes. However, transportation systems involve thousands or even millions of different “players” and an analysis by enumeration is hopeless. The good news is that the number of players is so great that little is lost in assuming that the players can be treated as a continuum.³ This allows us to work with smooth functions, greatly simplifying the process of finding equilibria. The remainder of this chapter introduces the basic traffic assignment problem in terms of the equilibrium concept and with a few motivating examples, but still in generally qualitative terms and restricted to small networks. The following two chapters provide us with the mathematical vocabulary and network tools needed to formulate and solve equilibrium on realistic, large-scale systems.

1.4 Traffic Assignment

There are many possible measures of effectiveness for evaluating the impacts of a roadway transportation project or policy. However, many of these can be calculated if one can predict the number of drivers on each roadway segment. These are called *link flows*. Predicting link flows allows a city or state government to evaluate different options.

If link flows are the output of a planning model, the main input is demographic data. That is, *given certain information about a population (number of people, income, amount of employment, etc.), we want to predict how many trips they will make, and how they will choose to travel*. Census records form an invaluable resource for this, often supplemented with travel surveys. Commonly, a medium-to-large random sample of the population is offered some money in exchange for keeping detailed diaries indicating all of the trips made within the

³This is analogous to solving structural design problems by assuming the usual stress-strain relationships, which assume a continuous material. In reality, a beam or column is composed of many distinct atoms, not a homogeneous material — but surely it is unnecessary to model each atom separately. The continuum assumption works almost as well and is much, much easier to work with.

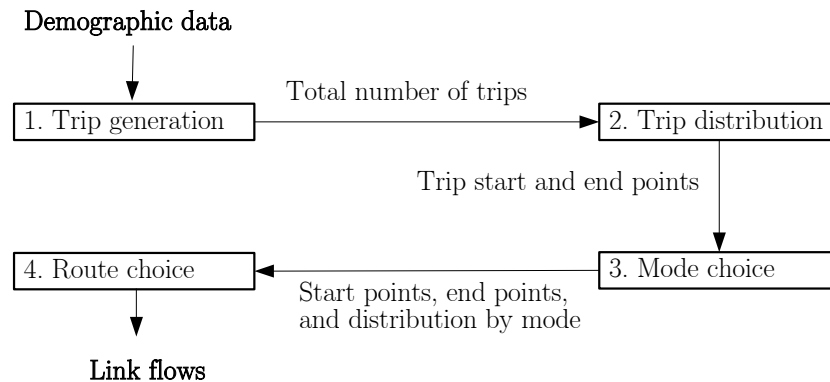


Figure 1.5: Schematic of the four-step process.

next several weeks, including the time of day, reason for traveling, and other details.

To get link flows from demographic data, most regions use the so-called *four-step model* (Figure 1.5). The first step is *trip generation*: based on demographic data, how many trips will people make? The second is *trip distribution*: once we know the total number of trips people make, what are the specific locations people will travel to? The third is *mode choice*: once we know the trip locations, will people choose to drive, take the bus, or use another mode? The fourth and final step is *route choice*, also known as *traffic assignment*: once we know the modes people will take to their trip destinations, what routes will they choose? Thus, at the end of the four steps, the transition from demographic data to link flows has been accomplished.⁴

Demographics are not uniform in a city; some areas are wealthier than others, some areas are residential while others are commercial, some parts are more crowded while other parts have a lower population density. For this reason, planners divide a city into multiple *zones*, and assume certain parameters within each zone are uniform. Clearly this is only an approximation to reality, and the larger the number of zones, the more accurate the approximation. (At the extreme, each household would be its own zone and the uniformity assumption becomes irrelevant.) On the other hand, the more zones, the longer it takes to run each model, and at some point computational resources become limiting. Typical networks used for large metropolitan areas have a few thousand zones. Zones are often related to census tracts, to make it easy to get demographic information from census results.

The focus of this book is the last of the four steps, traffic assignment. In the beginning, we assume that the first three steps have been completed, and we know the number of drivers traveling between each origin zone and desti-

⁴In more sophisticated models, the four steps may be repeated again, to ensure that the end results are consistent with the input data. There are also newer and arguably better alternatives to the four-step model.

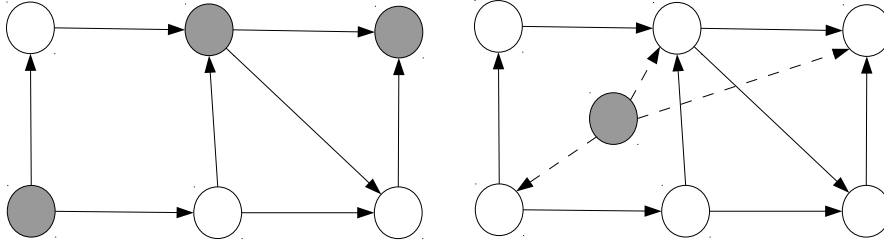


Figure 1.6: Centroids (shaded) coinciding with existing infrastructure, and artificial centroids. Dashed links on the right represent centroid connectors.

nation zone. From this, we want to know how many drivers are going to use each roadway segment, from which we can estimate congestion, emissions, toll revenue, or other measures of interest. We’ve already discussed several of the pieces of information we need in order to describe traffic assignment models precisely, including zones and travel demand. The final piece of the puzzle is a representation of the transportation infrastructure itself: the transportation network, which is described in the next section.

It is usually convenient to use a node to represent each zone; such nodes are called *centroids*, and all trips are assumed to begin and end at centroids. The set of centroids (for “zones”) is thus a subset of the set of nodes, defined in the next chapter. Centroids may coincide with physical nodes in the network. Centroids may also represent artificial nodes which do not correspond to any one physical point, and are connected to the physical infrastructure with links called *centroid connectors* (dashed lines in Figure 1.6).

1.5 Static Traffic Assignment

Figure 1.3 is the template for all traffic assignment models, be they static or dynamic: the choices of travelers lead to congestion patterns in the network (as predicted by a traffic flow model), and these patterns in turn influence the choices travelers make. The difference between static and dynamic traffic assignment lies in the traffic flow models used. Historically, static assignment models were the first to be developed, and research into dynamic models arose from the need to improve earlier, static models. Dynamic traffic assignment thus has a large number of parallels with static assignment; but where they differ, this difference is often intentional and important. Understanding these distinctions is key to knowing when dynamic models are appropriate to use.

1.5.1 Overview

In static assignment, the traffic flow model is based on *link performance functions*, which map the flow on each link to the travel time on that link. Mathe-

matically, if the notation (i, j) is used to refer to a roadway link connecting two nodes i and j , then x_{ij} is the flow on link (i, j) and the function $t_{ij}(x_{ij})$ gives the travel time on link (i, j) as a function of the flow on (i, j) . These functions $t_{ij}(\cdot)$ are typically assumed to be nonnegative, nondecreasing, and convex, reflecting the idea that as more vehicles attempt to drive on a link, the greater the congestion and the higher the travel times will be. A variety of link performance functions exist, but the most popular is the Bureau of Public Roads (BPR) function, which takes the form

$$t_{ij}(x_{ij}) = t_{ij}^0 \left(1 + \alpha \left(\frac{x_{ij}}{u_{ij}} \right)^\beta \right) \quad (1.1)$$

where t_{ij}^0 and u_{ij} are the free-flow time and capacity of link (i, j) , respectively, and α and β are shape parameters which can be calibrated to data. It is common to use $\alpha = 0.15$ and $\beta = 4$.

With such functions, the more travelers choose a path, the higher its travel time will be. Since travelers seek to minimize their travel time, travelers will not choose a path with high travel time unless there is no other option available. Indeed, if travelers only choose paths to minimize travel time, and if they have perfect knowledge of network conditions, then the network state can be described by the *principle of user equilibrium*: all used paths between the same origin and destination have equal and minimal travel times, for if this were not the case travelers would switch from slower routes to faster ones, which would tend to equalize their travel times.

It is not difficult to show that this *user equilibrium* state is not socially optimal, and that other configurations of traffic flow can reduce the average travel time (or even the travel time for all drivers) compared to the user equilibrium state. In other words, individual drivers seeking to minimize their own travel times will not always minimize travel times throughout the network, and this latter *system optimal* state can be contrasted with the user equilibrium one.

The prime advantage of using link performance functions like that in equation (1.1) is that the user equilibrium and system optimum states can be found with relative ease, even in realistic networks involving tens of thousands of links. Part II of the book discusses this in detail, showing how the static assignment problem can be formulated using the mathematical tools of optimization, fixed point, and variational inequality problems. These three representations of the equilibrium problem can be linked to powerful mathematical results which assure the existence and uniqueness of user equilibrium solutions under mild conditions. Efficient algorithms allow these states to be identified in a matter of minutes on large-scale networks.

For these reasons, static traffic assignment has been widely used in transportation planning practice for decades, and remains a powerful tool that can be used for performing alternatives analysis and generating forecasts of network conditions.

1.5.2 Critique

There are also a number of serious critiques of static assignment models, focused primarily on the link performance functions. By definition, static models do not monitor how network conditions (either demand or congestion) change over time, and implicitly assume a steady-state condition. This is clearly not the case. There are additional, subtler and more fundamental problems with link performance functions. This section describes a few of these problems.

First, not all vehicles on a link experience the same travel time. Even if the demand for travel on a link exceeds the capacity, the first vehicles to travel on that link will not experience much delay at all, while vehicles which arrive later may experience a very high delay. Together with the principle of user equilibrium, this means that the paths chosen by travelers will also depend on when they are departing. Route choices during periods of high congestion will be different from route choices made while this congestion is forming or dissipating. Furthermore, the travel time faced by a driver on a link depends crucially on the vehicles in front of them, and very little on the vehicles behind them. (A driver must adjust their speed to avoid colliding with vehicles downstream; a driver has no obligation to change their speed based on vehicles behind them.) This asymmetry is known as the *anisotropic* property of traffic flow, and it is violated by link performance functions — an increase in flow on the link is assumed to increase the travel time of all vehicles, and directly using link performance functions in a dynamic model would lead to undesirable phenomena, such as vehicles entering a link immediately raising the travel time for all other vehicles on the link, even those at the downstream end.

Second, the use of the word “flow” in static assignment is problematic. In traffic engineering, *flow* is defined as the (time) rate at which vehicles pass a fixed point on the roadway, and *capacity* is the greatest possible value for flow. By definition, flow cannot exceed capacity. However, the BPR function (1.1) imposes no such restriction, and it is common to see “flow-to-capacity” ratios much greater than one in static assignment.⁵ Instead, the x_{ij} values in static assignment are better thought of as *demand* rather than actual flow, since there is no harm in assuming that the demand for service exceeds the capacity, but it is impossible for the flow itself to exceed capacity. And for purposes of calibration, demand is much harder to observe than actual flow. These issues do not have clean resolutions.

Third, and related to the previous issue, link performance functions suggest that lower-capacity links have higher travel times under the same demand. But consider what happens at a freeway bottleneck, such as the lane drop shown in Figure 1.7. Congestion actually forms *upstream* of a bottleneck, and downstream of the lane drop there is no reason for vehicles to flow at a reduced speed. In

⁵There are several ways to add such a restriction, but these are less than satisfactory. A link performance function which tends to $+\infty$ as capacity is reached introduces numerical issues in solving for equilibrium. Explicitly adding link capacity constraints to the traffic assignment problem may make the problem infeasible, during peak periods there may be no way to assign all vehicles to the network without (temporarily) exceeding capacity.

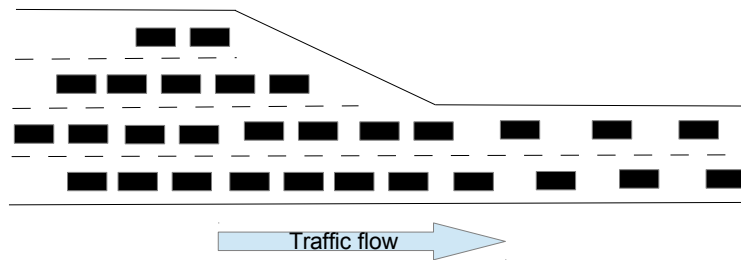


Figure 1.7: Congestion arises upstream of a bottleneck, not on the link with reduced capacity.

reality, it is upstream links that suffer when the demand for traveling on a link exceeds the capacity, not the bottleneck link itself.

Fourth, in congested urban systems it is very common for queues to fill the entire length of a link, causing congestion to spread to upstream links. This is observed on freeways (congested offramp queues) and in arterials (gridlock in central business districts) and is a major contributor to severe delay. In addition to the capacity, which is a maximum flow rate, real roadways also have a *jam density*, a maximum spatial concentration of vehicles. If a link is at jam density, no more vehicles can enter, which will create queues on upstream links. If these queues continue to grow, they will spread even further upstream to other links. Capturing this phenomenon can greatly enhance the realism of traffic models.

For all of these reasons, dynamic traffic assignment models shift to an entirely different traffic flow model. Rather than assuming simple, well-behaved link performance functions for each link we turn to traffic flow theory, to find more realistic ways to link traffic flow to congestion. Some early research in dynamic traffic assignment attempted to retain the use of link performance functions — for instance, modeling changes in demand by running a sequence of static assignment models over shorter time periods, with methods for linking together trips spanning multiple time periods. While this addresses the obvious shortcoming of static models, that they cannot accommodate changes in demand or model changes in congestion over time, it does nothing to address the more serious and fundamental problems with link performance functions described above. For this reason, this approach has largely been abandoned in favor of entirely different traffic flow models.

1.6 Dynamic Traffic Assignment

Dynamic traffic assignment arose from efforts to resolve the problems with static assignment noted in the previous section. While it is being used more and more in practice, it has not completely supplanted static traffic assignment. This is partially due to institutional inertia, but is also due to a few drawbacks associated with more realistic traffic flow models. Both these advantages and

drawbacks are discussed in this section.

1.6.1 Overview

The biggest difference between static and dynamic traffic assignment is in the traffic flow models used. A number of different traffic flow models are available, and a number of them are discussed in the following chapter. At a minimum, a traffic flow model for dynamic traffic assignment must be able to track changes in congestion at a fairly fine resolution, on the order of a few seconds. To do this, the locations of vehicles must be known at the same resolution. Most of them also address some or all of the shortcomings of link performance functions identified above.

The behavior of drivers is similar to that in static assignment in that drivers choose paths with minimum travel time. However, since congestion (and therefore travel time) changes over time, the minimum-time paths also change with time. Therefore, the principle of user equilibrium is replaced with a principle of *dynamic user equilibrium*: All paths used by travelers departing the same origin, for the same destination, *at the same time*, have equal and minimal travel times. Travelers departing at different times may experience different travel times; all the travelers departing at the same time will experience the same travel time at equilibrium, regardless of the path they choose. By virtue of representing demand changes over time, dynamic traffic assignment can also incorporate departure time choice of drivers, as well as route choice. Some dynamic traffic assignment models include both of these choices, while others focus only on route or departure time choice. Which choices are appropriate to model depends on which is more significant for a particular application, as well as the data and computational resources available. Most chapters of this book focus only on route choice, and as a general rule we will assume that departure times are fixed. In a few places we show how departure time choice can be added in.

It is important to specify that this equilibrium is based on the travel times the drivers actually experience, not the “instantaneous” travel times at the moment they depart. That is, we do not simply assume that drivers will pick the fastest route based on current conditions (as would be provided by most advanced traveler information services), but that they will anticipate changes in travel times which will occur during their journey. This suggests that drivers are familiar enough with the network that they know how congestion changes with time. This distinction is important — dynamic traffic assignment equilibrates on experienced travel times, not instantaneous travel times.

This means that it is impossible to find the dynamic user equilibrium in one step. Experienced travel times cannot be calculated at the moment of departure, but only after the vehicle has arrived at the destination. Therefore, dynamic traffic assignment is an iterative process, where route choices are updated at each iteration until an (approximate) dynamic user equilibrium solution has been found. This iterative process virtually always involves three steps, shown in Figure 1.8:

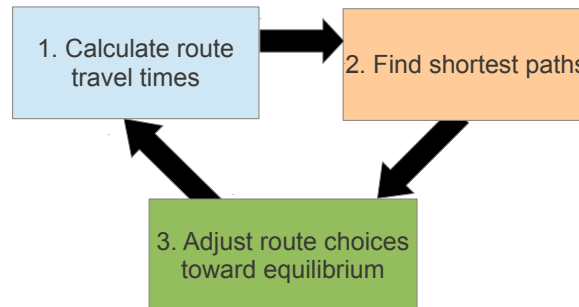


Figure 1.8: Three-step iterative process for dynamic traffic assignment.

Network loading: This is the process of using a traffic flow model to calculate the (time-dependent) travel times on each link, taking the routes and departure times of all drivers as inputs. In static assignment, this step was quite simple, involving nothing more than evaluating the link performance functions for each network link. In dynamic traffic assignment, this involves the use of a more sophisticated traffic flow model, or even the use of a traffic simulator. Network loading is discussed in Chapter 10.

Path finding: Once network loading is complete, the travel time of each link, at each point in time, is obtained. From this, we find the shortest path from each origin to each destination, *at each departure time*. Since we need experienced travel times, our shortest path finding must take into account that the travel time on each link depends upon the time a vehicle enters. This requires *time-dependent shortest path* algorithms, which are discussed in Chapter 11.

Route updating: Once time-dependent shortest paths have been found for all origins, destinations, and departure times, vehicles can be shifted from their current paths onto these new, shortest paths. As in static assignment, this step requires care, because shifting vehicles will change the path travel times as well. A few options for this are discussed in Chapter 12, along with other issues characterizing dynamic equilibrium. Unfortunately, and in contrast with static assignment, dynamic user equilibrium need not always exist, and when it exists it need not be unique.

1.6.2 Critique

The primary advantage of dynamic traffic assignment, and a significant one, is that the underlying traffic flow models are much more realistic. Link performance functions used in static assignment are fundamentally incapable of representing many important traffic phenomena. However, dynamic traffic assignment is not without its drawbacks as well.

Dynamic assignment requires considerably more computational time and memory than static assignment. As computers advance, this drawback is becoming less severe; but regardless of the computational resources available, one can run a large number of static traffic assignments in the time required for a single dynamic traffic assignment run. It may be advantageous to examine a large number of scenarios with static assignment, rather than a single run with dynamic traffic assignment, particularly if there is a lot of uncertainty in the input data.

Dynamic assignment models also require more input data for calibration. In addition to the usual link parameters such as capacity and free-flow time, dynamic traffic assignment models require a time-dependent origin-destination matrix, often known as a *demand profile*. Estimating even a static origin-destination matrix is difficult; estimating a dynamic demand profile can be even harder.

Furthermore, in addition to simply requiring more input data, dynamic traffic assignment also requires more accurate input data. Dynamic traffic assignment tends to be much more sensitive to the input data. Unfortunately, these models are more sensitive precisely because they are more realistic — features such as queue spillback mean that even a single erroneous capacity value can have ramifications throughout the entire network, not just on the link with the wrong value. If there is great uncertainty in the inputs (for instance, when attempting to predict demand decades into the future), then using a dynamic traffic assignment model may actually be further away from the truth than a static model, despite its more “realistic” congestion model.

Separately, dynamic traffic assignment is a relatively young field relative to static assignment, and there is no consensus on a single formulation. There are a large number of software packages (and an even larger number of academic models) which make differing assumptions and can lead to distinct results. By contrast, the optimization, variational inequality, and fixed point formulations in static assignment are completely standard, and therefore these models are more transparent. This book attempts to provide a high-level perspective, along with detailed discussions of a few of the most common modeling choices.

Finally, dynamic traffic assignment generally lacks neat, exact mathematical properties. In many dynamic traffic assignment models, one can create examples where no dynamic user equilibrium exists, or where this equilibrium is not unique. Proving convergence of iterative schemes is also more difficult.

All of these drawbacks must be traded off against the increased realism of dynamic traffic assignment models. Both static and dynamic traffic assignment models have their place as distinct tools for transportation engineering, and you should learn to identify circumstances where one or the other is more appropriate. As a general rule of thumb, dynamic models are most appropriate when the input data are known with high certainty (as in present-day traffic studies), and when queue lengths or other detailed congestion measures are required. Static models, by contrast, perform best when there is considerable uncertainty in the input data or when running a large number of scenarios is more important than detailed congestion information.

1.7 Target Audience

This book is primarily intended for first-year graduate students, but is also written with other potential audiences in mind. The content should be fully accessible to highly-prepared undergraduate students, and certain specialized topics would be appropriate for advanced graduate students as well. The text covers a large number of topics, likely more than would be covered in one or two semesters, and would also be useful for self-paced learners, or practitioners who may want in-depth learning on specific topics. We have included some supplementary material in optional sections, marked with an asterisk, which we believe are interesting, but which can be skipped without loss of continuity.

The most important prerequisites for this book are an understanding of multivariate calculus, and the intellectual maturity to understand the tradeoffs involved in mathematical modeling. Modeling does not involve one-size-fits-all approaches, and dogma about the absolute superiority of one model or algorithm over another is scarce. Instead, the primary intent of this book is to present a survey of important approaches to modeling transportation network problems, as well as the context to determine when particular models or algorithms are appropriate for real-world problems. Readers who can adopt this perspective will gain the most from the book.

Chapter 3 covers the mathematics needed for the network models in this book. Readers of this book will have significantly different mathematical backgrounds, and this chapter is meant to collect the necessary results and concepts in one place. Depending on your background, some parts of this chapter may need only a brief review, while other parts may be completely new.

While this book does not explicitly cover how to program these models and algorithms into a computer, if you have some facility in a programming language, it is highly instructive to try to implement them as you read. Many network algorithms are tedious to apply by hand or with other manual tools (calculator, spreadsheet). Computer programming will open the door to applying these models to interesting problems of large, realistic scale.

1.8 Exercises

1. [10] If all models are wrong, how can some of them be useful?
2. [10] All of the links in Figure 1.1 have a “mirror” connecting the same nodes, but in the opposite direction. When might mirror links not exist?
3. [23] A network is called *planar* if it can be drawn in two dimensions without links crossing each other. (The left network in Figure 1.2 is planar, but not the network on the right.) Do we expect to see planar graphs in transportation network modeling? Does it depend on the mode of transportation? What other factors might influence whether a transportation network is planar?

4. [35] Table 1.1 shows how networks can represent five types of transportation infrastructure. List at least five more systems (not necessarily in transportation) that can be modeled by networks, along with what nodes and links would represent.
5. [45] What factors might determine how large of a geographic area is modeled in a transportation network (e.g., corridor, city, region, state, national)? Illustrate your answer by showing how they would apply to the hypothetical projects or policies at the start of Section 1.1.
6. [45] What factors might determine the level of detail in a transportation network (e.g., freeways, major arterials, minor arterials, neighborhood streets)? Illustrate your answer by showing how they would apply to the hypothetical projects or policies at the start of Section 1.1.
7. [21] Provide an intuitive explanation of the prisoner's dilemma, as described in the Erica-Fred game of Table 1.3. Why does it happen? Name other real-world situations which exhibit a similar phenomenon.
8. [1] Explain the difference between the demand for travel on a link, and the flow on a link.
9. [42] Explain why a less realistic model less sensitive to correct input data may be preferred to a more realistic model more sensitive to correct inputs, and in what circumstances. Give specific examples.

Chapter 2

Network Representations and Algorithms

This chapter introduces networks as they are used in the transportation field. Section 2.1 introduces the mathematical terminology and notation used to describe network elements. Section 2.2 discusses two special kinds of networks which are used frequently in network analysis, acyclic networks and trees. Section 2.3 then describes several ways of representing a network in a way computers can use when solving network problems. This third section can be skipped if you do not plan to do any computer programming. Section 2.4 describes the shortest path problem, a classic network algorithm which plays a central role in traffic assignment.

2.1 Terminology

Because of its relative youth, there are a variety of notational conventions in the transportation networks community. A common notation is adopted in the book to present the methods and topics in a consistent manner, but you should be aware that other authors may use slightly different conventions and definitions. Table 2.1 gives an example of some terms which are often used synonymously (or nearly synonymously) with ours. These differences are mainly a matter of style, not substance, but when reading other books or papers you should pay careful attention to the exact wording of their definitions.

The fundamental construct we will utilize is the *network*. A network is a collection of *nodes*, and a collection of *links* which connect the nodes. A network is denoted $G = (N, A)$, where N is the set of nodes and A is the set of links. Figure 2.1(a) shows a simple network, with four nodes in the set $N = \{1, 2, 3, 4\}$ and five links in the set $A = \{(1, 2), (1, 3), (2, 3), (2, 4), (3, 4)\}$. Notice that the notation for each link contains the two nodes connected by the link: the upstream node is called the *tail* of the link, and the downstream node the *head*. We will often refer to the total number of nodes in a network as n

Table 2.1: A thesaurus of network terminology

Terminology in this book	Alternative terms
Network	Graph
Node	Vertex
Link	Arc, edge, line
Tree	Arborescence

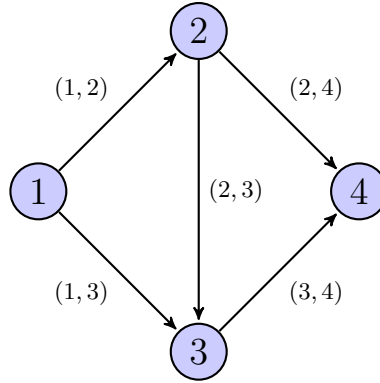


Figure 2.1: Example network notation

and the total number of links as m . This book is solely concerned with the case where n and m are finite — networks with infinitely many nodes and links are sometimes studied in theoretical mathematics, but rarely in transportation applications.

All of the networks in this book are *directed*. In a directed network, a link can only be traversed in one direction, specified by the ordering of the nodes. Therefore, $(1,2)$ and $(2,1)$ represent different links: they connect the same nodes, but represent travel in opposite directions. Unless stated otherwise, we assume that there are no “self-loops” (i,i) from a node to itself, and no parallel links, so the notation (i,j) is unambiguous. This is not usually a restrictive assumption, since we can introduce artificial nodes to split up self-loops or parallel links, as shown in Figure 2.2. The upper left panel in this figure shows a network with a self-loop $(2,2)$. By introducing a fourth node in the bottom left, we have divided the self-loop into two links $(2,4)$ and $(4,2)$. In the upper right panel of the figure, there are two networks connecting nodes 2 and 3, so the notation $(2,3)$ does not tell us which of these two links we are referring to. By introducing a fourth node in the bottom right, we now have three links: $(2,3)$, $(2,4)$, and $(4,3)$, which collectively represent both of the ways to travel between nodes 2 and 3 in the original network, but without parallel links. These new nodes are *artificial* in the sense that they do not represent physical transportation infrastructure, but are inserted for modeling convenience. Artificial nodes (and

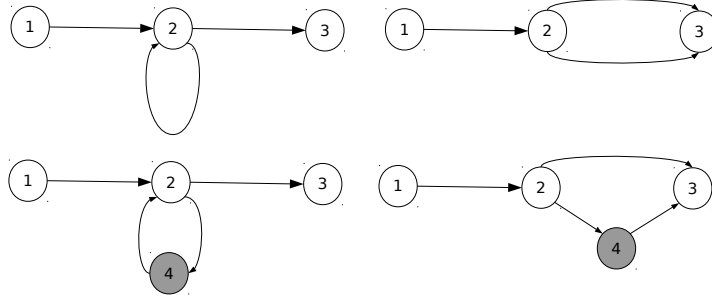


Figure 2.2: Artificial links can be introduced to simulate self-loops and parallel links.

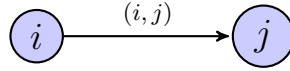


Figure 2.3: Generic two-node network

artificial links) play an important role in simplifying certain network problems, as discussed throughout the book.

To introduce some more terminology, Figure 2.3 shows a two-node network with nodes i and j connected by the link (i, j) . In this figure, we say that link (i, j) is *incident* to both i and j because it is connected to both. We would further say that (i, j) is an *outgoing* link of i and an *incoming* link of j . If $(i, j) \in A$, then we say that node j is *adjacent* to node i . In this example, j is adjacent to i , but i is not adjacent to j . Adjacency can also be applied to links; two links are adjacent if the head of the first link is the tail of the second link. In Figure 2.1, both links $(2, 3)$ and $(2, 4)$ are adjacent to $(1, 2)$.

The *forward star* of a node i is the set of all outgoing links, denoted $\Gamma(i)$; the *reverse star* is the set of all incoming links, denoted $\Gamma^{-1}(i)$. In Figure 2.1, $\Gamma(2) = \{(2, 3), (2, 4)\}$ and $\Gamma(3) = \{(3, 4)\}$, while $\Gamma^{-1}(2) = \{(1, 2)\}$ and $\Gamma^{-1}(3) = \{(1, 3), (2, 3)\}$.

The *degree* of a node is the total number of links incident to that node. The node degree can be separated into the *indegree* and *outdegree* of a node. The indegree is the total number of incoming links at a node, while the outdegree is the number of outgoing links at a node. The degree is then the sum of the indegree and outdegree. Referring back to Figure 2.1, node 2 has an indegree of 1, an outdegree of 2 and a degree of 3.

A *path* π is a sequence of adjacent links connecting two nodes i_0 and i_k . We can either write π as an ordered set of links

$$\{(i_0, i_1), (i_1, i_2), (i_2, i_3), \dots, (i_{k-1}, i_k)\},$$

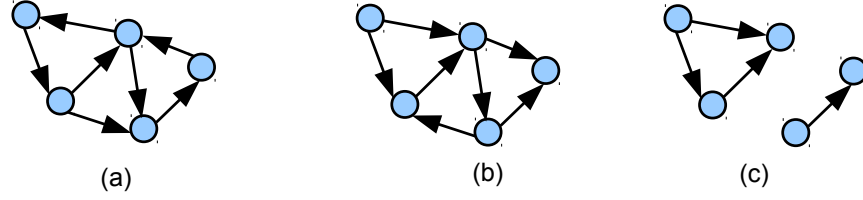


Figure 2.4: Networks which are (a) strongly connected; (b) connected, but not strongly connected; (c) disconnected.

or more compactly, by the nodes passed on the way with the notation

$$[i_0, i_1, i_2, i_3, \dots, i_{k-1}, i_k].$$

A path is a *cycle* if the starting and ending nodes are the same ($i_0 = i_k$). Paths which contain a cycle are called *cyclic*; paths which do not have a cycle as a component are called *acyclic*. Cyclic paths are uncommon in transportation problems, so unless stated otherwise, we only consider acyclic paths. Let Π^{rs} denote the set of all acyclic paths connecting nodes r and s , and let Π denote the set of all acyclic paths in the entire network, that is, $\Pi = \cup_{(r,s) \in N^2} \Pi^{rs}$. A network is *connected* if there is at least one path connecting any two nodes in the network, assuming that the links can be traversed in either direction (that is, ignoring the direction of the link); otherwise it is *disconnected*. A network is *strongly connected* if there is at least one path connecting any two nodes in the network, obeying the given directions of the links. Figure 2.4 shows networks which are strongly connected; connected, but not strongly connected; and disconnected.

2.2 Acyclic Networks and Trees

Networks which do not have any cycles are called *acyclic networks*. Acyclic networks are extremely important, because the lack of cycles can greatly simplify analysis. Even when networks actually have cycles in reality, many efficient transportation network algorithms temporarily divide the network into a set of acyclic subnetworks.

A defining characteristic of acyclic networks is that a *topological order* can be established. A topological order is a labeling of each node with a number between 1 and n (n is the number of nodes), such that every link connects a node of lower topological order to a node of higher topological order. Every path, therefore, consists of a sequence of nodes in increasing topological order. As an example, in Figure 2.1, the nodes are labeled in a topological order: each path ($[1, 3, 4]$, $[1, 2, 4]$, and $[1, 2, 3, 4]$) traverses nodes in increasing numerical order. (The phrase “a” topological order is used because it may not be unique, and there may be several ways to label the nodes so that links always connect

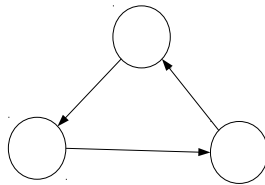


Figure 2.5: Cyclic network with no topological order.

lower-numbered nodes to higher-numbered ones.) In general networks, a topological order may not exist (see Figure 2.5, where a cycle makes it impossible to label nodes so that the numbers increase when traversing any link). The following theorem shows that the existence of a topological order is a *defining* characteristic of acyclic networks.

Theorem 2.1. *A topological order exists on a network if and only if it is acyclic.*

Proof. (\Rightarrow) Assume a topological order exists on G . Then every path contains a sequence of nodes whose topological order is strictly increasing (if this were not so, then there is a link connecting a higher topological order node to a lower one, contradicting the definition of topological order). Therefore, no path can repeat the same node more than once, and the network is acyclic.

(\Leftarrow) Assume the network G is acyclic. Then we can prove existence of a topological order by construction. Let $o(i)$ be the topological order of node i . We will describe a procedure to build such a topological order one step at a time. At any point in time, a “marked” node is one which has a topological order assigned, and an “unmarked” node is one which does not yet have its topological order. Because G is acyclic, there is at least one node r with no incoming links. (You will be asked to prove this statement as an exercise.) Let $o(r) = 1$. Again, because G is acyclic, there is at least one unmarked node with no incoming links from an unmarked node; mark this node with the topological order 2. This process can be repeated until all nodes are marked; if at any point every unmarked node has an incoming link from another unmarked node, then a cycle exists, contradicting the assumption of acyclicity. Otherwise, the topological order will have been constructed. The topological order so constructed is valid because when a node is assigned an order, the only incoming links are from nodes which have already been marked, and thus have a lower topological order. \square

A *tree* is a special type of acyclic network, which also shows up frequently in transportation network algorithms. A tree is defined as a network in which there is a unique node r (called the *root*) which has the property that exactly one path exists between r and every other node in the network.¹ Figure 2.6 shows a tree. Other useful properties of trees are:

¹In graph theory, this is usually called an *arborescence*, and trees are defined for undirected graphs. However, the transportation community generally uses the term “tree” in both cases, a convention followed in this book.

- A tree has at least two nodes with degree one.
- In a tree, there is at most one path between any two nodes.
- A tree would become disconnected if any of its links were deleted.
- Every node i in a tree (except the root) has a unique incoming link; the tail node of that unique link is called the *parent* of node i . Similarly, the head nodes of the links in the forward star of i are its *children*.

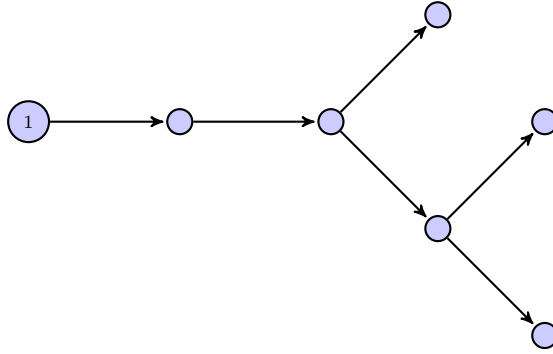


Figure 2.6: A tree

2.3 Data Structures

In practice, the methods and techniques in this book are carried out by computer programs. Thus, it is important not only to have a convenient way to theoretically represent a network (the network structure) but a way to store data so that it can be easily accessed and utilized by computer programs. To demonstrate, Figure 2.7 presents a network similar to those earlier, with additional information provided. In this case, the extra information is the travel time of the link from i to j , denoted t_{ij} .

The first data structure is the node-link incidence matrix, shown in Figure 2.8. The columns are indexed by the network links, and the rows by the network nodes. A “1” in the matrix indicates an outgoing link from that node, while a “−1” indicates an incoming link. A “0” indicates that the link is not incident to that particular node. This is a fairly inefficient data structure, especially for large networks, as in nearly any case there are going to be a large number of zeros in every row. Note that the forward and reverse stars of node i can respectively be identified with the columns in the matrix which have a 1 or −1 in the i -th row.

The second data structure, a node-node adjacency matrix is a simpler representation of a network but also suffers from inefficiencies. Each row represents a node i and each column a node j . A “1” indicates the existence of a link (i, j) .

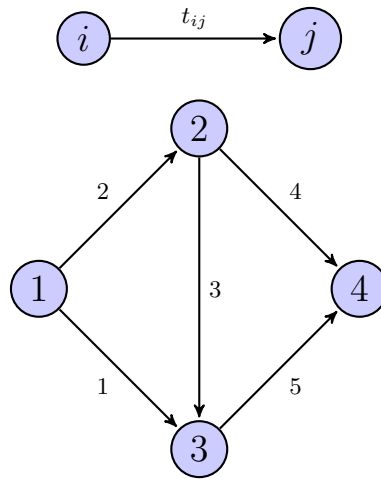


Figure 2.7: Graph used for network representation examples

$$\begin{array}{c}
 \begin{array}{ccccc}
 & (1, 2) & (1, 3) & (2, 3) & (2, 4) & (3, 4) \\
 \begin{array}{l} 1 \\ 2 \\ 3 \\ 4 \end{array} & \begin{bmatrix} 1 & 1 & 0 & 0 & 0 \\ -1 & 0 & 1 & 1 & 0 \\ 0 & -1 & -1 & 0 & 1 \\ 0 & 0 & 0 & -1 & -1 \end{bmatrix}
 \end{array}
 \end{array}$$

Figure 2.8: Node-link incidence matrix for the network in Figure 2.7.

$$\begin{array}{c}
\begin{array}{cccc}
& 1 & 2 & 3 & 4 \\
\begin{array}{c} 1 \\ 2 \\ 3 \\ 4 \end{array} & \begin{bmatrix} 0 & 1 & 1 & 0 \\ 0 & 0 & 1 & 1 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 0 & 0 \end{bmatrix}
\end{array}
\end{array}$$

Figure 2.9: Node-node adjacency matrix for the network in Figure 2.7.

In this structure, a non-zero entry provides both existence and direction information. An example of a node-node adjacency matrix is given in Figure 2.9. In very dense networks, with many more links than nodes ($m \gg n$), node-node adjacency matrices can be efficient.

A shortcoming of the previous data structures is that they do not contain information beyond the existence and direction of the links in a network. In transportation applications we are typically working with networks for which we want to store and use more information about each element. For example, we may want to store information about the travel time, number of lanes (capacity), speed limit, signalization, etc. for each link. Adjacency lists give us an opportunity to store this information efficiently in a manner that is compatible with programming languages and their built-in data structures.

Storing a network as an array is more efficient in terms of space, and is easier to work with in some programming languages. Storage is not wasted on a large number of zeros, as with the adjacency matrices. A “forward star” representation is presented below. In the forward star representation, links are sorted in order of their tail node, as shown in Table 2.2. A second array, called **point**, is used to indicate where the forward star of each node begins. If the network has n nodes and m links, the **point** array has $n + 1$ entries, and **point**($n + 1$) is always equal to $m + 1$; the reason why is discussed below. As shown in Table 2.2, **point**(2) = 3 because the link 3 is the first link adjacent to node 2; **point**(3) = 5 because link 5 is the first link adjacent to node 3, and so forth. Three conventions associated with the star representation are:

1. The number of entries in the **point** array is one more than the number of nodes.
2. We automatically set **point**($n+1$) = $m + 1$.
3. If there are no outgoing links, **point**(i) = **point** ($i+1$)

With these conventions, we can say that the forward star for any node i consists of all links whose IDs are greater than or equal to **point**(i), and strictly less than **point**($i+1$). This representation is most useful when we frequently need to loop over all of the links leaving a particular node, which can be accomplished by programming a “for” loop between **point**(i) and **point**($i+1$) - 1, and referring to the link arrays with these indices. We can make this statement universally, because we defined the **point** array to have one more entry than the number of nodes. If **point** only had n entries, then referring to **point**($n+1$)

Table 2.2: Forward star representation of a network.

Node	Point	Link	tail	head	cost	capacity
1	1	1	1	3	1	...
2	3	2	1	2	2	...
3	5	3	2	3	3	...
4	6	4	2	4	4	...
5	6	5	3	4	6	...

would be meaningless in the above statements. It is also possible to store a network in a “reverse star” representation along similar lines, which is most useful when we frequently need to loop over all the links entering a particular node. The exercises explore this further, along with an array-based method for easily iterating over both the forward and reverse stars of a node.

The purpose of this brief discussion of data structures is not to provide a comprehensive treatment of the topic. Rather, it is intended to present some basic ideas related to network storage and data representation. It should also prompt you to think about data structures carefully when working with networks. Transportation network problems tend to be large in size and complicated in nature. The difference between computer code that takes minutes to run, as opposed to hours, is often the way data is stored and passed.

2.4 Shortest Paths

As a first network algorithm, we’ll discuss how to find the shortest (least cost) path between any two nodes in a network in an efficient and easily-automated way. In shortest path algorithms, we are given a network $G = (N, A)$; each link $(i, j) \in A$ has a *fixed* cost c_{ij} . The word “cost” does not necessarily mean a monetary cost, and refers to whatever total quantity we are trying to minimize. In traffic applications, we often use the link travel times t_{ij} as the cost, to find the least travel-time path. If there are tolls in a network, the cost of each link may be the sum of the toll on that link, and the travel time on that link multiplied by a “value of time” factor converting time into monetary units. Costs may reflect still other factors, and some of these are explored in the exercises. Because we can solve all of these problems in the same way, we may as well use a single name for all of these quantities we are trying to minimize, and “cost” has become the standard term.

The word “fixed” in the last paragraph is emphasized because in many transportation problems the cost is dependent on the total flow and vehicle route choices, and it is not reasonable to assume fixed costs. However, even in such network problems the most practical solution methods involve solving several shortest path problems as part of an iterative framework, updating paths as travel times change. The context for a shortest path algorithm is to find the least travel-time path between two nodes, at the current travel times — with

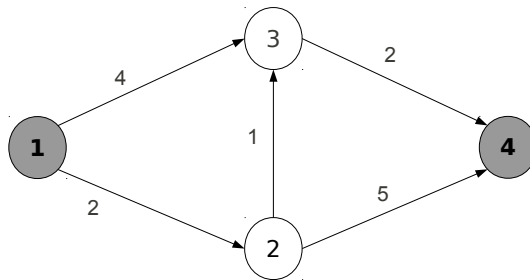


Figure 2.10: Example network for shortest paths (link costs indicated).

everybody's route choices held fixed. By separating the process of identifying the shortest path from the process of shifting path flows toward the shortest path, we simplify the problem and end up with something which is relatively easy to solve. Later, in Section 4.2, we will show how this problem can be phrased in the language of optimization, by identifying an objective function, decision variables, and constraints. Here, we develop specialized algorithms for the shortest path problem which are more efficient and more direct than general optimization techniques.

Although the shortest path problem most obviously fits into transportation networks, many other applications also exist in construction management, geometric design, operations research, and many other areas. For instance, the fastest way to solve a Rubik's Cube from a given position can be solved using a shortest path algorithm, as can the fewest number of links needed to connect an actor to Kevin Bacon when playing Six Degrees of Separation.

A curious fact of shortest path algorithms is that finding the shortest path from a single origin to *every other* node is only slightly harder than finding the shortest path from that origin to a single destination. This also happens to be the reason why we can find shortest paths without having to list all of the zillions of possible paths from an origin to a destination, and adding up their paths. This common reason is *Bellman's Principle*, which states that *any segment of a shortest path must itself be a shortest path between its endpoints*. For instance, consider the network in Figure 2.10, where the costs c_{ij} are printed next to each link. The shortest path from node 1 to node 4 is $[1, 2, 3, 4]$. Bellman's Principle requires that $[1, 2, 3]$ also be a shortest path from nodes 1 to 3, and that $[2, 3, 4]$ be a shortest path from nodes 2 to 4. It further requires that $[1, 2]$ be a shortest path from node 1 to node 2, $[2, 3]$ be a shortest path from node 2 to node 3, and $[3, 4]$ be a shortest path from node 3 to node 4. You should verify that this is true with the given link costs.

To see why this must be the case, assume that Bellman's Principle was violated. If the cost on link $(1, 3)$ was reduced to 2, then $[1, 2, 3]$ is no longer a shortest path from node 1 to node 3 (that path has a cost of 3, while the single-

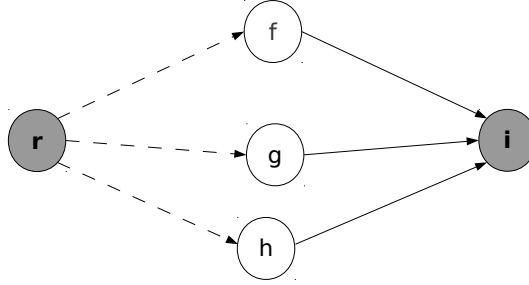


Figure 2.11: Illustration of Bellman's Principle (dashed lines indicated shortest paths between two nodes).

link path $[1, 3]$ has cost 2). Bellman's Principle then implies that $[1, 2, 3, 4]$ is no longer the shortest path between nodes 1 and 4. Why? The first part of the path can be replaced by $[1, 3]$ (the new shortest path between 1 and 3), reducing the cost of the path from 1 to 4: $[1, 3, 4]$ now has a cost of 4. In general, if a segment of a path does *not* form the shortest path between two nodes, we can replace it with the shortest path, and thus reduce the cost of the entire path. Thus, the shortest path must satisfy Bellman's Principle for all of its segments.

The implication of this is that *we can construct shortest paths one node at a time*, proceeding inductively. Let's say we want to find the shortest path from node r to a node i , and furthermore let's assume that we've already found the shortest paths from r to every node which is directly upstream of i (nodes f , g , and h in Figure 2.11). The shortest path from r to i must pass through either f , g , or h ; and according to Bellman's Principle, the shortest path from r to i must be either (a) the shortest path from r to f , plus link (f, i) ; the shortest path from r to g , plus link (g, i) ; or the shortest path from r to h , plus link (h, i) . This is efficient because, rather than considering all of the possible paths from r to i , we only have to consider three, which can be easily compared. Furthermore, we can re-use the information we found when finding shortest paths to f , g , and h , and don't have to duplicate the same work when finding the shortest path to i . This idea doesn't give a complete algorithm yet — how did we find the shortest paths to f , g , and h , for instance? — but gives the flavor of the shortest path algorithms presented next.

Bellman's Principle also gives us a compact way of expressing all of the shortest paths from an origin to every other node in the network: for each node, simply indicate the *last* node in the shortest path from that origin. This is called the backnode vector \mathbf{q}^r , where each component q_i^r is the node immediately preceding i in the shortest path from r to i . If $i = r$, then q_i^r is not well-defined (i is the origin itself; what is the shortest path from the origin to itself, and if we can define it, what node immediately precedes the origin?) so we say $q_r^r \equiv -1$ by definition. For the network in Figure 2.10 (with the original costs),

we thus have $q_1^1 = -1$, $q_2^1 = 1$, $q_3^1 = 2$, and $q_4^1 = 3$, or, in vector notation, $\mathbf{q}^1 = [-1 \ 1 \ 2 \ 3]$.

The backnode vector can be used as follows: say we want to look up the shortest path from node 1 to node 4. Starting at the destination, the backnode of 4 is 3, which means “the shortest path to node 4 is the shortest path to node 3, plus the link (3, 4).” To find the shortest path to node 3, consult its backnode: “the shortest path to node 3 is the shortest path to node 2, plus the link (2, 3).” For the shortest path to node 2, its backnode says: “the shortest path to node 2 is from node 1.” This is the origin, so we’ve found the start of the path, and can reconstruct the original path to node 4: [1, 2, 3, 4]. More briefly, we can use the backnodes to trace the shortest path back to an origin, by starting from the destination, and reading back one node at a time.

We will also define L_i^r to be the total cost on the shortest path from origin r to node i (the letter L is used because these values are often referred to as node labels), with $L_r^r \equiv 0$, so in this example we would have $L_1^1 = 0$, $L_2^1 = 2$, $L_3^1 = 3$, and $L_4^1 = 5$.

This chapter presents four shortest path algorithms. The first, in Section 2.4.1, only applies when the network is acyclic but is extremely fast and simple. Sections 2.4.2 and 2.4.3 next present the two most common approaches for solving shortest paths in networks with cycles: *label setting* and *label correcting*. Dijkstra’s algorithm and the Bellman-Ford algorithm are the quintessential label setting and label correcting algorithms, respectively, and they are used to show the difference between these approaches. These three algorithms actually find the shortest path from the origin r to *every* other node, not just the destination s ; this exploits Bellman’s Principle, because the shortest path from r to s must also contain the shortest path from r to every node in that path. Further, both of them contain a set of labels L_i^r associated with each node, representing the shortest path from r to node i . With a label setting approach, each node’s label is determined once and only once. On the other hand, with a label correcting approach, each node’s label can be updated multiple times.

The fourth algorithm, A^* is presented in Section 2.4.4. This algorithm can be faster than Dijkstra’s algorithm if we are only interested in the shortest path from a single origin to a single destination.

2.4.1 Shortest paths on acyclic networks

Recall from Section 2.2 that an *acyclic* network is one in which no cyclic paths exist; it is impossible to visit any node more than once on any path. Conversely, a *cyclic* network is one where a cycle does exist, and where (in theory) one could repeatedly drive around in a circle forever. Transportation networks are usually cyclic, and in fact any network where you can reach every node from every other node must be cyclic. (Why?) However, we’ll take a short diversion into acyclic networks for the purposes of finding shortest paths, for two reasons: (1) finding the shortest path on an acyclic network is much simpler and faster, and makes for a good first illustration; and (2) more advanced solution methods take advantage of the fact that people do *not* use cyclic paths (because of the shortest

path assumption), so we can only look at an acyclic portion of the network and thereby use the much faster shortest path algorithm for acyclic networks.

Once we have a topological order on a network, it becomes very easy to find the shortest path from any node r to any other node s (clearly r has lower topological order than s), using the following algorithm which is based directly on Bellman's Principle:

1. Initialize by setting $L_i^r = \infty$ and $q_i^r = -1$ for all nodes $i \in N$ to indicate that we have not found any shortest paths yet. Then set $L_r^r = 0$, because the distance from the origin to itself is zero.
2. Let i be the node topologically following r .
3. We can find the shortest path from r to i using Bellman's Principle, using these formulas:

$$L_i^r = \min_{(h,i) \in \Gamma^{-1}(i)} \{L_h^r + c_{hi}\} \quad (2.1)$$

$$q_i^r \in \arg \min_{(h,i) \in \Gamma^{-1}(i)} \{L_h^r + c_{hi}\} \quad (2.2)$$

(Essentially, these formulas have us search all of the links arriving at node i for the approach with minimum cost.)

4. Does $i = s$? If so, stop, and we have found the shortest path from r to s , which has cost L_s^r . Otherwise, let i be the next node topologically, and return to step 3.

Bellman's Principle lets us find all of the shortest paths in one pass over the nodes (in topological order), because there are no cycles which could make us loop back.

Here's how to apply this algorithm to the network in Figure 2.10, with node 1 as origin and node 4 as destination.

Step 1: Initialize: $\mathbf{L}^1 \leftarrow [0 \ \infty \ \infty \ \infty]$ and $\mathbf{q}^1 \leftarrow [-1 \ -1 \ -1 \ -1]$.

Step 2: We set $j \leftarrow 2$.

Step 3: The only approach to node 2 is (1,2), so $\Gamma^{-1}(2)$ is just $\{(1,2)\}$ and the minimization is trivial: $L_2^1 \leftarrow L_1^1 + c_{12} = 0 + 2 = 2$ and $q_2^1 \leftarrow 1$.

Step 4: $2 \neq 4$, so we advance j to the next node topologically: $j \leftarrow 3$ and return to step 3.

Step 3: There are two approaches to node 3, (1,3) and (2,3). So $L_3^1 \leftarrow \min\{L_1^1 + c_{13}, L_2^1 + c_{23}\} = \min\{0 + 4, 2 + 1\} = 3$, which corresponds to approach (2,3), so $q_3^1 \leftarrow 2$.

Step 4: $3 \neq 4$, so we advance j to the next node topologically: $j \leftarrow 4$ and return to step 3.

Step 3: There are two approaches to node 4, (2, 4) and (3, 4). So $L_4^1 \leftarrow \min\{L_2^1 + c_{24}, L_3^1 + c_{34}\} = \min\{2 + 5, 3 + 2\} = 5$, which corresponds to approach (3, 4), so $q_4^1 \leftarrow 3$.

Step 4: We have reached the destination (node 4) and terminate.

A simple induction proof shows that this algorithm must give the correct shortest paths to each node (except for those topologically before r , since such paths do not exist). By “give the correct shortest paths,” we mean that the labels \mathbf{L} give the lowest cost possible to each node when traveling from r , and that the backnodes \mathbf{q} yield shortest paths when traced back to r . Assume that the nodes are numbered in topological order. Clearly L_r^r and q_r^r are set to the correct values in the first step (the shortest path from r to itself is trivial), and are never changed again because the network proceeds in increasing topological order. Now assume that the algorithm has found correct L and q values for all nodes between r and k . In the next step, it updates labels for node $k + 1$. Let $(i, k + 1)$ be the last link in a shortest path from r to $k + 1$. By Bellman’s principle, the first part of this path must be a shortest path from r to i . Since i is topologically between r and k , by the induction hypothesis the L^r and q^r labels are correct for all nodes immediately upstream of k . Therefore, $L_i^r + c_{i,k+1} \leq L_j^r + c_{j,k+1}$ for all immediately upstream nodes j , and the algorithm makes the correct choices in step 3 for node $k + 1$.

2.4.2 Shortest paths on cyclic networks: label correcting

When there are cycles in the network, this previous approach can’t be applied, because there is no clear sequence in which to examine nodes and apply Bellman’s Principle. However, we can generalize the approach. Rather than scanning nodes in a rigid order, we can fan out from the origin, keeping in mind that we may have to scan a node multiple times in case of a cycle. To keep track of the nodes we need to examine, we define a scan eligible list SEL , a set of nodes that we still need to examine before we have found all of the shortest paths. This is a *label correcting* approach, because nodes may be scanned multiple times, and the labels updated.

1. Initialize every label L_i^r to ∞ , except for the origin, where $L_r^r \leftarrow 0$, and the backnode vector $\mathbf{q}^r \leftarrow -\mathbf{1}$.
2. Initialize the scan eligible list to contain all nodes immediately downstream of the origin, that is, $SEL \leftarrow \{i : (r, i) \in \Gamma(r)\}$.
3. Choose a node $i \in SEL$ and remove it from that list.
4. Scan node i by applying the same equations as in the previous algorithm:

$$L_i^r = \min_{(h,i) \in \Gamma^{-1}(i)} \{L_h^r + c_{hi}\} \quad (2.3)$$

$$q_i^r = \arg \min_{(h,i) \in \Gamma^{-1}(i)} \{L_h^r + c_{hi}\} \quad (2.4)$$

5. If the previous step changed the value of L_i^r , then add all nodes immediately downstream of i to the scan eligible list, that is, $SEL \leftarrow SEL \cup \{j : (i, j) \in \Gamma(i)\}$.
6. If SEL is empty, then terminate. Otherwise, return to step 3.

Repeating the same example, this algorithm works as follows:

- Step 1:** We set $\mathbf{L}^1 \leftarrow [0 \quad \infty \quad \infty \quad \infty]$ and $\mathbf{q}^1 \leftarrow [-1 \quad -1 \quad -1 \quad -1]$.
- Step 2:** We set $SEL \leftarrow \{2, 3\}$.
- Step 3:** Choose node 3 and remove it from SEL , so $i \leftarrow 3$, and $SEL = \{2\}$.
- Step 4:** The two possible approaches are $(1, 3)$ and $(2, 3)$, so $L_3^1 \leftarrow \min\{L_1^1 + c_{13}, L_2^1 + c_{23}\} = \min\{0 + 4, \infty + 1\} = 4$, which corresponds to approach $(1, 3)$, so $q_3^1 \leftarrow 1$. Note that, unlike in the case of the acyclic network, we are not claiming that this is the shortest path to node 3. We are simply claiming that this is the shortest path we have found thus far.
- Step 5:** The previous step reduced L_3^1 from ∞ to 4, so we add the downstream node 4 to SEL , so $SEL = \{2, 4\}$.
- Step 6:** SEL is nonempty, so we return to step 3.
- Step 3:** Choose node 4 and remove it from SEL , so $i \leftarrow 4$, and $SEL = \{2\}$.
- Step 4:** There are two approaches to node 4, $(2, 4)$ and $(3, 4)$. So $L_4^1 \leftarrow \min\{L_2^1 + c_{24}, L_3^1 + c_{34}\} = \min\{\infty + 5, 4 + 2\} = 6$, which corresponds to approach $(3, 4)$, so $q_4^1 \leftarrow 3$.
- Step 5:** The previous step reduced L_4^1 from ∞ to 6, but there are no downstream nodes to add, so $SEL = \{2\}$.
- Step 6:** SEL is nonempty, so we return to step 3.
- Step 3:** Choose node 2 and remove it from SEL , so $i \leftarrow 2$, and $SEL = \emptyset$.
- Step 4:** The only approach to node 2 is $(1, 2)$, so $\{(h, i) \in \Gamma^{-1}(i)\}$ is just $\{(1, 2)\}$ and the minimization is trivial: $L_2^1 \leftarrow L_1^1 + c_{12} = 0 + 2 = 2$ and $q_2^1 \leftarrow 1$.
- Step 5:** The previous step reduced L_2^1 from ∞ to 2, so we add the downstream nodes 3 and 4 to SEL , so $SEL = \{3, 4\}$.
- Step 6:** SEL is nonempty, so we return to step 3.
- Step 3:** Choose node 3 and remove it from SEL , so $i \leftarrow 3$, and $SEL = \{4\}$. Note that this is the *second* time we are scanning node 3. We have to do this, because we found a path to node 2, and there is a possibility that this could lead to a shorter path to node 3 as well.

Step 4: The two possible approaches are $(1, 3)$ and $(2, 3)$, so $L_3^1 \leftarrow \min\{L_1^1 + c_{13}, L_2^1 + c_{23}\} = \min\{0 + 4, 2 + 1\} = 3$, which corresponds to approach $(2, 3)$, so $q_3^1 \leftarrow 2$. Note that we have changed the labels again, showing that we have indeed found a better path through node 2.

Step 5: The previous step reduced L_3^1 from 4 to 3, so we add the downstream node 4 to SEL , so $SEL = \{4\}$.

Step 6: SEL is nonempty, so we return to step 3.

Step 3: Choose node 4 and remove it from SEL , so $i \leftarrow 4$, and $SEL = \emptyset$.

Step 4: There are two approaches to node 4, $(2, 4)$ and $(3, 4)$. So $L_4^1 \leftarrow \min\{L_2^1 + c_{24}, L_3^1 + c_{34}\} = \min\{2 + 5, 3 + 2\} = 5$, which corresponds to approach $(3, 4)$, so $q_4^1 \leftarrow 3$.

Step 5: The previous step reduced L_4^1 from 6 to 5, but there are no downstream nodes to add, so SEL remains empty.

Step 6: SEL is empty, so we terminate.

As you can see, this method required more steps than the algorithm for acyclic networks (because there is a possibility of revisiting nodes), but it does not rely on having a topological order and can work in any network. One can show that the algorithm converges no matter how you choose the node from SEL , but it is easier to prove if you choose a systematic rule. If you operate SEL in a “first-in, first-out” manner, always choosing a node which has been in the queue for the greatest number of iterations, it is possible to show that after m iterations, you have certainly found the shortest paths from r which consist of only a single link. (You’ve probably found quite a few more shortest paths, but even in the worst case you’ll have found at least these.) After $2m$ iterations, you will have certainly found the shortest paths from r which consist of one or two links only, and so on. So, after mn iterations, you will have found all of the shortest paths, since a shortest path cannot use more than n links. The exercises ask you to fill in the details of this proof sketch.

2.4.3 Shortest paths on general networks: label setting

In the label-correcting method, we accounted for cycles by allowing nodes to be scanned more than once, relieving us of having to scan nodes in the “right” order when there is no topological order to guide us. A label-setting method works in cyclic networks, yet only has to scan each node once. This requires a bit more care to scan nodes in the proper order, since there is no opportunity to revisit nodes to correct their labels later on. So, while label-setting methods only require one scan per node, determining which node to scan requires more effort at each iteration.

The quintessential label setting algorithm was developed by Edsger Dijkstra. Dijkstra’s algorithm finds the shortest path from the origin r to every other node,

when every link has a nonnegative cost. Furthermore, at each step it finds the shortest path to at least one additional node. It uses the concept of *finalized* nodes, that is, nodes to which the shortest path has already been found. It uses the same \mathbf{L} and \mathbf{q} labels as before. The primary distinction between label setting and label correcting is that label setting requires fewer iterations than label correcting, but each iteration takes a greater amount of effort. Which one is superior depends on the network topology, the specific implementation of each algorithm, and the skill of the programmer.

Dijkstra's algorithm can be stated as follows:

1. Initialize every label L_i^r to ∞ , except for the origin, where $L_r^r \leftarrow 0$.
2. Initialize the set of finalized nodes $F = \emptyset$ and the backnode vector $\mathbf{q}^r \leftarrow -\mathbf{1}$.
3. Find an unfinalized node i (i.e., not in F) for which L_i^r is minimal.
4. Finalize node i by adding it to set F ; if all nodes are finalized ($F = N$) then terminate.
5. Update labels for links leaving node i : for each $(i, j) \in A$, update $L_j = \min\{L_j^r, L_i^r + c_{ij}\}$.
6. Update backnodes: for each $(i, j) \in A$, if L_j^r was reduced in the previous step set $q_j^r \leftarrow i$.
7. If all nodes are finalized, stop. Otherwise, return to step 3

This algorithm is considered *label setting*, because once a node's label is updated, the node is finalized and never revisited again. We consider how this algorithm can be applied to the same example used for the other algorithm demonstrations.

Initialization. We initialize $\mathbf{L}^1 = [0 \ \infty \ \infty \ \infty]$, $F = \emptyset$, and $\mathbf{q}^1 = [-1 \ -1 \ -1 \ -1]$.

Iteration 1. The unfinalized node with least L value is the origin, so we set $i = 1$ and finalize it: $F = \{1\}$. We update the downstream labels: $L_2^1 = \min\{\infty, 0 + 2\} = 2$ and $L_3^1 = 4$. Both labels were reduced, so we update the backnodes: $q_2^1 = q_3^1 = 1$.

Iteration 2. Of the unfinalized nodes, $i = 2$ has the lowest label, so we finalize the node: $F = \{1, 2\}$. We update the downstream labels: $L_3^1 = \min\{4, 2 + 1\} = 3$ and $L_4^1 = 7$. Both labels were reduced, so we update the backnodes: $q_3^1 = q_4^1 = 2$.

Iteration 3. Of the unfinalized nodes, $i = 3$ has the lowest label, so we finalize it: $F = \{1, 2, 3\}$. We update the downstream label $L_4^1 = \min\{7, 3 + 2\} = 5$ and backnode $q_4^1 = 3$.

Iteration 4. Node 4 is the only unfinalized node, so $i = 4$. We finalize it, and since $F = N$ we are done.

The correctness of Dijkstra's algorithm can be proved by an induction argument, on the size of the set of finalized nodes $|F|$. Specifically, we show that at every step of the algorithm the backnodes to every finalized node represent a shortest path. The algorithm always finalizes the origin r first, so when $|F| = 1$ the only finalized node is the origin, and the shortest path to this node is trivial.

Now assume that when $|F| = k$, the algorithm has correctly found shortest paths to every finalized node, and consider what happens when the $k + 1$ -th node is added to this set. This occurs in step 4 of the algorithm. Call this newly finalized node j , and say that at this point the backnode is $q_j^r = i$. By contradiction, assume that there is a shorter path from r to j than the one through node i . Notice that the entries of the backnode vector which are not -1 all point to finalized nodes, because step 6 is the only step which changes this vector, and only changes entries of this vector to a newly finalized node. Furthermore, from the induction hypothesis, the backnode vectors from all of these previously-finalized nodes describe shortest paths.

So, if there is a shorter path to j than the one through i , it must differ in the last link, that is, this shorter path must end with a link (h, j) where $h \neq i$. Either h is a finalized node, or it is not. If h were finalized, then by the induction hypothesis L_h^r is the correct shortest path cost to r , and since it leads to a shorter path to j we must have $L_h^r + c_{hj} < L_i^r + c_{ij}$. But if this were the case, then steps 5 and 6 would have set q_j^r to h when node h was finalized, so q_j^r could not equal i , a contradiction. Or, if h were not finalized, then $L_h^r > L_j^r$, since in step 3 we always finalize a node with the smallest L value. Since all links have nonnegative costs, we cannot have $L_h^r + c_{hj} < L_j^r$, which is what we need if h led to a shorter path to j . Since both cases lead to contradictions, we conclude that a correct shortest path to node j has been found, and by induction, this argument extends to every node.

2.4.4 Shortest paths from one origin to one destination

The previous sections gave algorithms to find the shortest path from one origin to all destinations. Slight modifications to these algorithms can find the shortest path from all origins to one destination (see Exercise 47). As discussed above, Bellman's principle lets us re-use information from one shortest path when finding another, and as a result even finding a shortest path from one origin to one destination can provide many other shortest paths "for free." In traffic assignment, where we are modeling the flows of many drivers, it often makes sense to use an algorithm that finds many shortest paths simultaneously.

However, there are times when we are only concerned with a single origin and destination, and do not care about the "free" shortest paths to other destinations we get. Examples are in-vehicle routing systems, or if the number of origin-destination pairs with positive demand is small compared to the overall size of the matrix (i.e., the OD matrix is sparse) so there are relatively few paths we need to find. In such cases we can use a more focused algorithm to find a single shortest path in less time than it takes to find shortest paths to all destinations.

The A^* algorithm is a simple modification to label-setting, which in some

cases can dramatically reduce the running time, in exchange for limiting the scope to one origin r and one destination s . This algorithm requires an additional value g_i^s for each node, representing an estimate of the cost of the shortest path from i to s . This estimate (often called the “heuristic”) should be a *lower bound* on the actual cost of the shortest path from i to s . Some examples of how these estimates are chosen are discussed below.

Once the estimates g_i^s are chosen, the label-setting algorithm from Section 2.4.3 proceeds as before with a small modification in Step 3: rather than choosing a node i in E which minimizes L_i^r , we choose a node which minimizes $L_i^r + g_i^s$. Everything else proceeds exactly as before. The intuition is that we are taking into account our estimates of which nodes are closer to the destination. The algorithm in Section 2.4.3 fans out in all directions from the origin (by simply looking at L_i^r), rather than directing the search towards a particular destination.

It can be shown that this algorithm will always yield the correct shortest path from r to s as long as the g_i^s are lower bounds on actual shortest path costs from i to s . If this is not the case, A^* is not guaranteed to find the shortest path from r to s . Some care must be taken in how these estimates are found. Two extreme examples are:

- Choose $g_i^s = 0$ for all i . This is certainly a valid lower bound on the shortest path costs (recall that label-setting methods assume nonnegative link costs), so A^* will find the shortest path from r to s . However, zero is a very poor estimate of the actual shortest path costs. With this choice of g_i^s , A^* will run exactly the same as Dijkstra’s algorithm, and there is no time savings.
- Choose g_i^s to be the *actual* shortest path cost from i to s . This is the tightest possible “lower bound,” and will make A^* run extremely quickly — in fact, it will only scan nodes along the shortest path, the best possible performance that can be achieved. However, coming up with these estimates is just as hard as solving the original problem! So in the end we aren’t saving any effort; what we gain from A^* is more than lost by the extra effort we need to compute g_i^s in the first place.

So, there is a tradeoff between choosing tight bounds (the closer g_i^s to the true costs, the faster A^* will be) while not spending too long in computing the estimates (which might swamp the savings in A^* itself). Luckily, in transportation networks, there are several good bounds available which can be computed fairly quickly. For instance:

- The Euclidean (“as the crow flies”) distance between i and s , divided by the fastest travel speed in the network, is a lower bound on the travel time between i and s .
- Replace every link with a lower bound on its cost (say, free-flow travel time) and find shortest paths between all nodes and all destinations (repeatedly using one of the previous algorithms from this chapter). This

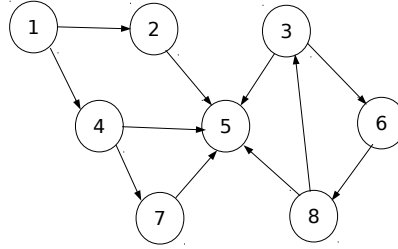


Figure 2.12: Network for Exercises 1–3

takes more time, but only needs to be done once and can be done as a preprocessing step. As we will see in Chapter 7, solving traffic assignment requires many shortest path computations. The extra time spent finding these costs once might result in total time savings over many iterations.

You may find it instructive to think about other ways we can estimate g_i^s values, and how they might be used in transportation settings.

2.5 Exercises

1. [9] In the network in Figure 2.12, list the indegree, outdegree, degree, forward star, and reverse star of each node.
2. [3] In the network in Figure 2.12, list all of the paths between nodes 1 and 5.
3. [4] State whether the network in Figure 2.12 is or is not (a) cyclic; (b) a tree; (c) connected; (d) strongly connected.
4. [15] For each of the following, either draw a network with the stated properties, or explain why no such network can exist: (a) connected, but not strongly connected; (b) strongly connected, but not connected; (c) cyclic, but not strongly connected.
5. [25] If m and n are the number of links and nodes in a network, show that $m < n^2$.
6. [25] If a network is connected, show that $m \geq n - 1$.
7. [25] If a network is strongly connected, show that $m \geq n$.
8. [10] Show that $\sum_{i \in N} |\Gamma(i)| = \sum_{i \in N} |\Gamma^{-1}(i)| = m$.
9. [1] Why must transportation infrastructure networks be cyclic?
10. [15] Although we do not usually expect drivers to use cyclic paths, there are some exceptions. Name one.

23. [14] Write down the node-node adjacency matrix of the network in Figure 2.12.
24. [16] Consider the network defined by this node-node adjacency matrix:

$$\begin{bmatrix} 0 & 0 & 0 & 1 & 1 & 0 & 0 \\ 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 1 & 1 & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 \end{bmatrix}$$

- (a) Draw the network.
- (b) How many links does the network have?
- (c) How many links enter node 1? How many links leave node 1?
25. [17] Is the network represented by the following node-node adjacency matrix strongly connected?

$$\begin{bmatrix} 0 & 0 & 0 & 1 & 1 & 0 & 0 \\ 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 \end{bmatrix}$$

26. [10] If the nodes in an acyclic network are numbered in a topological order, show that the node-node adjacency matrix is upper triangular.
27. [37] Let \mathbf{A} be the node-node adjacency matrix for a network. What is the interpretation of the matrix product \mathbf{A}^2 ?
28. [42] Let \mathbf{A} be the node-node adjacency matrix for an acyclic network. First show that $\sum_{n=1}^{\infty} \mathbf{A}^n$ exists, and give an interpretation of this sum.
29. [65] A *unimodular* matrix is a square matrix whose elements are integers and whose determinant is either $+1$ or -1 . A matrix is *totally unimodular* if every nonsingular square submatrix is unimodular. (Note that a totally unimodular matrix need not be square). Show that every node-link incidence matrix is totally unimodular.
30. [10] One disadvantage of the forward star representation is that it is time-consuming to identify the reverse star of a node — one must search through the entire array to find every link with a particular head node. Describe a “reverse star” data structure using arrays, where the reverse star can be easily identified.

Table 2.3: Backnode vector for Exercise 33.

Node	Backnode
1	4
2	6
3	-1
4	7
5	4
6	5
7	3
8	10
9	5
10	6

31. [52] By combining the forward star representation from the text and the reverse star representation from the previous exercise, we can quickly identify both the forward and reverse stars of every node. However, a naive implementation will have two different sets of arrays, one sorted according to the forward star representation, and the other sorted according to the reverse star representation. This duplication wastes space, especially if there are many attributes associated with each link (travel time, cost, capacity, etc.) Identify a way to easily identify the forward and reverse stars of every node, with only *one* set of arrays of link data, by adding an appropriate attribute to each link.
32. [58] In the language of your choice, write computer code to do the following:
- Produce the node-node adjacency matrix of a network when given the node-link incidence matrix.
 - Produce the node-link incidence matrix of a network when given the node-node adjacency matrix.
 - Produce the node-node adjacency matrix of a network when given the forward star representation of the network.
 - Produce the forward star representation of a network when given the node-node adjacency matrix.
33. [13] After solving a shortest path problem from node 3 to every other node, I obtain the backnode vector shown in Table 2.3. Write the shortest paths (a) from node 3 to node 5; (b) from node 3 to node 7; (c) from node 4 to node 8.
34. [26] Find the shortest path from node 1 to every other node in the network shown in Figure 2.14. Report the final labels and backnodes (L and q values) for all nodes.

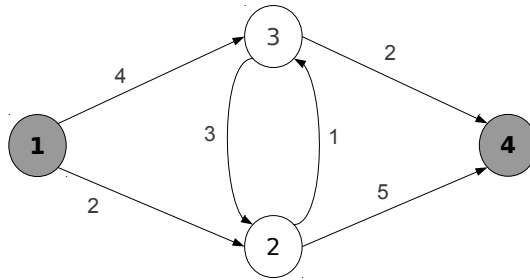


Figure 2.14: Network for Exercises 34 and 41.

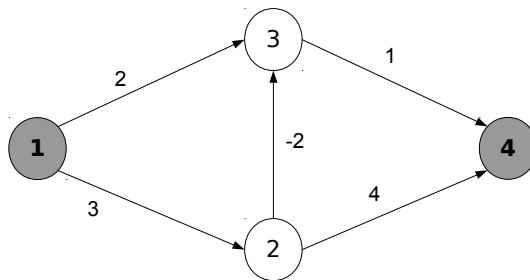


Figure 2.15: Network for Exercise 35.

35. [25] The network in Figure 2.15 has a link with a negative cost. Show that the label-correcting algorithm still produces the correct shortest paths in this network, while the label-setting algorithm does not.
36. [57] Prove or disprove the following statement: “Any network with negative costs can be transformed into a network with nonnegative costs by adding a large enough constant to every link’s cost. We can then use the label-setting algorithm on this new network. Therefore, the label-setting algorithm can find the shortest paths on any network, even if it has negative-cost links.” (Proving this statement means showing that it is true in any network; to disprove it, it is sufficient to find a single counterexample.)
37. [37]. Find the shortest paths on the network in Figure 2.16, using both the label-correcting and label-setting algorithms. Sketch the *shortest path tree* (that is, the network with only the links implied by the backnode vector) produced by each algorithm. Which algorithm required fewer iterations?
38. [33]. You and your friends are camping at Yellowstone Park, when you

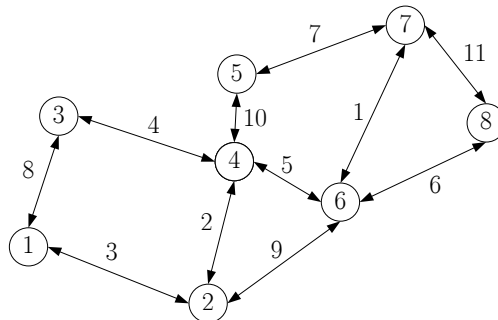


Figure 2.16: Network for Exercise 37, link labels are costs.

suddenly realize you have to be at the Denver airport in ten hours to catch a flight. You don't have Internet access, but you do have an atlas showing the distance and travel time between selected cities (Figure 2.17). Assuming these travel times are accurate, can you make it to Denver in time for your flight?

39. [43]. Consider the following variation of the shortest path problem: instead of the costs c_{ij} being fixed, instead assume that they are random, and the cost of link (i, j) takes the values $\{c_{ij}^1, c_{ij}^2, \dots, c_{ij}^s\}$ with probabilities $\{p_{ij}^1, p_{ij}^2, \dots, p_{ij}^s\}$ independent of the cost of any other link. How can you adapt the shortest path algorithms in this chapter to find the path with the least *expected* cost?
40. [53]. Assume that each link (i, j) in the network fails with probability p_{ij} , and that link failures are independent. (Failure can represent some kind of damage or disruption, the probability of being detected in a military routing problem, etc.) A path fails if any link in that path fails. Explain how you can find the most reliable path (that is, the one with the least failure probability) with a shortest path algorithm.
41. [23] Instead of trying to find the *shortest* path between two nodes, let's try to find the *longest* path between two nodes. (As we will see later, there are actually cases when this is useful.)
 - (a) Modify the algorithm presented in Section 2.4.1 to find the longest path in an acyclic network.
 - (b) If you modify the label-correcting or label-setting algorithms for general networks in a similar way, will they find the longest paths successfully? Try them on the network in Figure 2.14.
42. [25]. As a modification to the node-node adjacency matrix you might use to represent the network in a computer program (cf. Section 2.3), you

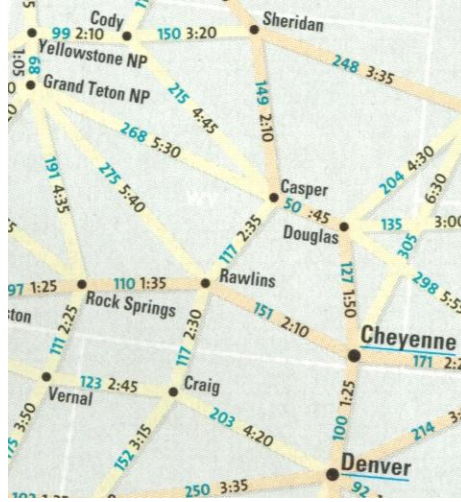


Figure 2.17: Atlas page for Exercise 38.

could store the cost of the links in this matrix, with ∞ where no link exists (as opposed to ‘1’ where links exist and ‘0’ where they do not). Find the shortest path between nodes 7 and 4 using the modified adjacency matrix below, where the entry in row i and column j is the cost c_{ij} if this link exists, and ∞ if it does not.

$$\begin{bmatrix} \infty & \infty & \infty & 5 & 3 & \infty & \infty \\ 4 & \infty & 7 & \infty & \infty & \infty & \infty \\ \infty & \infty & \infty & 6 & \infty & \infty & \infty \\ \infty & \infty & \infty & \infty & 2 & \infty & \infty \\ \infty & \infty & \infty & \infty & \infty & 4 & \infty \\ 3 & 2 & \infty & \infty & \infty & \infty & \infty \\ \infty & 1 & \infty & \infty & \infty & 4 & \infty \end{bmatrix}$$

43. [22]. In the game “Six Degrees of Kevin Bacon,” players are given the name of an actor or actress, and try to connect them to Kevin Bacon in as few steps as possible, winning if they can make the connection in six steps or less. Two actors or actresses are “connected” if they were in the same film together. For example, Alfred Hitchcock is connected to Kevin Bacon in three steps: Hitchcock was in *Show Business at War* with Orson Welles, who was in *A Safe Place* with Jack Nicholson, who was in *A Few Good Men* with Kevin Bacon. Beyoncé Knowles is connected to Kevin Bacon in two steps, since she was in *Austin Powers: Goldmember*, where Tom Cruise had a cameo, and Cruise was in *A Few Good Men* with Bacon. Assuming that you have total, encyclopedic knowledge of celebrities and films, show how you can solve “Six Degrees of Kevin Bacon” as a shortest path problem. Specify what nodes, links, costs, origins, and destinations

represent in the network you construct.

44. [14] In a network with no negative-cost cycles, show that $-nC$ is a lower bound on the shortest path cost between any two nodes in a network.
45. [57] Show that if the label-correcting algorithm is performed, and that at each iteration you choose a node in SEL which has been in the list the longest, at the end of kn iterations the cost and backnode labels correctly reflect all shortest paths from r which are no more than k links long. (Hint: try an induction proof.)
46. [44] Assume that the label-correcting algorithm is terminated once the label for some node i falls below $-mC$. Show that the following the current backnode labels from i will lead to a negative-cost cycle.
47. [32] Modify three of the shortest path algorithms in this chapter so that they find shortest paths from all origins to one destination, rather than one origin to all destinations:
 - (a) The acyclic shortest path algorithm from Section 2.4.1.
 - (b) The label-correcting algorithm from Section 2.4.2.
 - (c) The label-setting algorithm from Section 2.4.3.
48. [73] It is known that the label correcting algorithm will find the correct shortest paths as long as the initial labels \mathbf{L} correspond to the distance of *some* path from the origin (they do not necessarily need to be initialized to $+\infty$). Assume that we are given a vector of backnode labels \mathbf{q} which represents some tree (not necessarily the shortest paths) rooted at the origin r . Develop a one-to-all shortest path algorithm that uses this vector to run more efficiently. In particular, if the given backnode labels \mathbf{q} *do* correspond to a shortest path tree, your algorithm should recognize this fact and terminate in a number of steps linear in the number of network links.

Chapter 3

Mathematical Concepts

This chapter provides a survey of mathematical techniques used in network analysis. There is no attempt to be comprehensive; indeed, many books have been written about each of the sections in this chapter and the next. Rather, the intent is to cover topics which are used frequently in transportation network problems. Where the coverage is brief, the assumption is that the reader has seen the ideas before and simply needs a refresher; if this is not the case, you are advised to seek out additional resources beyond this book. Concepts which we expect are new to most readers are explained in more detail. Section 3.1 begins by reviewing definitions and facts related to vectors, matrices, sets, and functions which are needed, including the topics of convex functions, convex sets, and multivariable calculus (the gradient vector, and the Jacobian and Hessian matrices will play a particularly important role). Section 3.2 introduces the fixed point problem and variational inequality, two methods for mathematically modeling equilibrium problems. A third tool, optimization, is introduced in the following chapter.

3.1 Basic Definitions

This section reviews the mathematical machinery needed to understand equilibrium in transportation problems. This involves understanding certain properties of matrices, sets, and functions. The descriptions are relatively terse, because there are countless textbooks and online resources with additional explanations and examples. It is also a very incomplete overview, only covering concepts which will be used later in the book. The section also collects a number of useful results related to matrices, sets, and functions. The proofs of most of these are left as exercises.

3.1.1 Indices and summation

We commonly work with attributes of links or nodes, such as the total number of trips starting or ending at a node, or the travel time on a link. In such cases, it makes sense to use the same variable letter for the number of trips produced at a node (say P), and to indicate the productions at a specific node with a subscript or superscript (P_1 or P^3 .) Unfortunately, the superscript notation can be ambiguous, since the superscript in P^3 can be interpreted both as an exponent, and as the index referring to a particular node. In most cases it will be clear which meaning is intended, and parentheses can be used in cases which are unclear: $(P^3)^2$ is the square of the productions at node 3. To avoid confusion, subscripts are generally preferred to superscripts, but superscripts can make the notation more compact when there are multiple indices. An example is δ_{ij}^π from Section 1.5, where the variable δ is indexed both by a link (i, j) and a path π . We have striven to be consistent with which indices appear as subscripts and which appear as superscripts. However, in a few occasions, being absolutely rigid in this regard would create cluttered formulas. In such cases, we have sacrificed rigor for readability, and hope that you are not thrown off by an index typically appearing as a subscript in the superscript position, and vice versa.

Subscript or superscript indices also allow formulas to be written more compactly. A common example is the summation notation, such as

$$\sum_{i=1}^5 P_i = P_1 + P_2 + P_3 + P_4 + P_5. \quad (3.1)$$

The left-hand side of Equation (3.1) is used as shorthand for the right-hand side. More formally, the left-hand side instructs us to choose all the values of i between 1 and 5 (inclusive); collect the terms P_i for all of these values of i ; and to add them together. A variant of this notation is

$$\sum_{i \in N} P_i, \quad (3.2)$$

which expresses the sum of the productions from all nodes in the set N . Here i ranges over all elements in the set N , rather than between two numbers as in (3.1). We can also add conditions to this range by using a colon. If we only wanted to sum over nodes whose products are less than, say, 500, we can write

$$\sum_{i \in N: P_i < 500} P_i. \quad (3.3)$$

When it is exceptionally clear what values or set i is ranging over, we can simply write

$$\sum_i P_i, \quad (3.4)$$

but this “abbreviated” form should be used sparingly, and avoided if there is any ambiguity at all as to what values i should take in the sum.

When there are multiple indices, a double summation can be used:

$$\sum_{i=1}^3 \sum_{j=2}^3 x_{ij} = x_{12} + x_{22} + x_{32} + x_{13} + x_{23} + x_{33}. \quad (3.5)$$

You can think of a double sum either as a “nested” sum:

$$\sum_{i=1}^3 \sum_{j=2}^3 x_{ij} = \sum_{i=1}^3 \left(\sum_{j=2}^3 x_{ij} \right) = \sum_{i=1}^3 (x_{i2} + x_{i3}), \quad (3.6)$$

which expands to the same thing as the right-hand side of (3.5), or as summing over all combinations of i and j such that i is between 1 and 3, and j is between 2 and 3. Triple sums, quadruple sums, and so forth behave in exactly the same way, and are common when there are many indices.

The summation index is often called a *dummy variable*, because the indices do not have an absolute meaning. Rather, they are only important insofar as they point to the correct numbers to add up. For instance, $\sum_{i=1}^5 P_i$ and $\sum_{i=0}^4 P_{i+1}$ are exactly the same, because both expressions have you add up P_1 through P_5 . Likewise, $\sum_{i=1}^5 P_i$ and $\sum_{j=1}^5 P_j$ are exactly the same, the fact that we are counting from 1 to 5 using the variable j instead of i is of no consequence.

Related to this, it is *wrong* to refer to a summation index outside of the sum itself. A formula such as $x_i + \sum_{i=1}^5 y_i$ is incorrect. For the formula to make sense, x_i needs to refer to one specific value of i . But using i as the index in the sum $\sum_{i=1}^5$ means that i must range over all the values between 1 and 5. Does y_i refer to the index of summation (which ranges from 1 to 5), or to the one specific value of i used outside the sum? If you want to refer to a specific node, as well as to index over all nodes, you can add a prime to one of them, as in

$$x_i + \sum_{i'=1}^5 y_{i'}, \quad (3.7)$$

or you can either use a different letter altogether, as in

$$x_i + \sum_{j=1}^5 y_j. \quad (3.8)$$

Both conventions are common.

In transportation network analysis, we frequently have to sum over all of the links which are in the forward or reverse star of a node (Section 2.3). If the link flows are denoted by a variable x , the total flow entering node i is the sum of the link flows in its reverse star. This can be written in several ways:

$$\sum_{(h,i) \in \Gamma^{-1}(i)} x_{hi} \qquad \sum_{(h,i) \in A} x_{hi} \quad (3.9)$$

The notation on the right can be a bit confusing at first glance, since it looks like we are summing the flow of *every* link (the whole set A), rather than just the links entering node i . The critical point is that in this formula, i refers to one specific node which was previously chosen and defined outside of the summation. *The only variable we are summing over in Equation (3.9) is h .* In this light, i is fixed, and the right-most sum is over all the values of h such that (h, i) is a valid link (that is, $(h, i) \in A$). These are exactly the links which form the reverse star of i .

Almost all of the sums we will see in this book involve only a finite number of terms. These sums are much easier to work with than infinite sums, and have the useful properties listed below. (Some of these properties do not always apply to sums involving infinitely many terms.)

1. You can factor constants out of a sum:

$$\sum_i c x_i = c \sum_i x_i, \quad (3.10)$$

no matter what values i ranges over. This follows from the distributive property for sums: $a(b + c) = ab + ac$. (A “constant” here is any term which does not depend on i .)

2. If what you are summing is itself a sum, you can split them up:

$$\sum_i (x_i + y_i) = \sum_i x_i + \sum_i y_i. \quad (3.11)$$

This follows from the commutative property for sums: $a + b = b + a$, so you can rearrange the order in which you add up the terms.

3. You can exchange the order of a double summation:

$$\sum_i \sum_j x_{ij} = \sum_j \sum_i x_{ij} \quad (3.12)$$

no matter what values i and j range over. This also follows from the commutative property.

None of these properties are anything new; they simply formalize how we can operate with the \sum notation using the basic properties of addition. These properties can also be combined. It is common to exchange the order of a double sum, and then factor out a term, to perform a simplification:

$$\sum_i \sum_j c_j x_{ij} = \sum_j \sum_i c_j x_{ij} = \sum_j \left(c_j \sum_i x_{ij} \right). \quad (3.13)$$

The last step is permissible because c_j is a “constant” relative to a sum over i . We could not do this step at the beginning, because c_j is not a constant

relative to a sum over j . This kind of manipulation is helpful if, say, $\sum_i x_{ij}$ has a convenient form or a previously-calculated value.

Occasionally it is useful to refer to products over an index. In analogy with the summation notation \sum used for sums over an index, products over an index are written with the notation \prod . As an example

$$\prod_{i=1}^5 P_i = P_1 P_2 P_3 P_4 P_5 \quad (3.14)$$

Exercise 3 asks you to investigate which of the properties of summation notation carry over to the product notation.

3.1.2 Vectors and matrices

A *scalar* is a single real number, such as 2, -3 , $\sqrt{2}$, or π . The notation $x \in \mathbb{R}$ means that x is a real number (a scalar). A *vector* is a collection of scalars; the *dimension* of a vector is the number of scalars it contains. For instance, $[3 \ -5]$ is a two-dimensional vector, $[6 \ 0 \ \pi]$ is a three-dimensional vector, and so forth. It is possible to write vectors either horizontally, with its component scalars in a row, or vertically, with its component scalars in a column. For the most part it doesn't matter whether vectors are written in a row or in a column; the exception is in formulas involving multiplication between vectors and matrices, as described below, where row and column vectors must be distinguished. Vectors are usually denoted with boldfaced lower-case letters, like \mathbf{x} ; the notation $\mathbf{x} \in \mathbb{R}^n$ indicates that \mathbf{x} is an n -dimensional vector.¹ (So, $[3 \ -5] \in \mathbb{R}^2$ and $[6 \ 0 \ \pi] \in \mathbb{R}^3$.) Individual components of vectors are often denoted with subscripts or superscripts, as in x_1 or x^1 . The *zero vector* is a vector with zeros for all of its components, and is denoted $\mathbf{0}$.

Two vectors of the same dimension can be added together by adding the corresponding components of each vector. If $\mathbf{x} = [1 \ 2]$ and $\mathbf{y} = [3 \ 4]$, then $\mathbf{x} + \mathbf{y} = [4 \ 6]$. Multiplying a vector by a scalar means multiplying each component of the vector by the scalar, so $3\mathbf{x} = [3 \ 6]$ and $-\mathbf{y} = [-3 \ -4]$. The *dot product* of two vectors of the same dimension is defined as

$$\mathbf{x} \cdot \mathbf{y} = \sum_i x_i y_i, \quad (3.15)$$

where the sum is taken over all vector components; with the example above, $\mathbf{x} \cdot \mathbf{y} = 1 \times 3 + 2 \times 4 = 11$. Note that the dot product of two vectors is a scalar. The magnitude of a vector \mathbf{x} is given by its *norm* $|\mathbf{x}| = \sqrt{\mathbf{x} \cdot \mathbf{x}}$. This norm provides a measure of distance between two vectors; the distance between \mathbf{x} and \mathbf{y} is given by $|\mathbf{x} - \mathbf{y}|$.

The dot product can also be written

$$\mathbf{x} \cdot \mathbf{y} = |\mathbf{x}| |\mathbf{y}| \cos \theta \quad (3.16)$$

¹There are some exceptions; for instance, shortest-path labels are traditionally denoted with an upper-case L .

where θ is the angle between the vectors \mathbf{x} and \mathbf{y} if they are both drawn from a common point. In particular, \mathbf{x} and \mathbf{y} are perpendicular if $\mathbf{x} \cdot \mathbf{y} = 0$.

A matrix is a rectangular array of scalars. If a matrix has m rows and n columns, it is called an $m \times n$ matrix, and is an element of $\mathbb{R}^{m \times n}$. A matrix is *square* if it has the same number of rows and columns. In this book, matrices are denoted by boldface capital letters, such as \mathbf{X} or \mathbf{Y} . In the examples that follow, let \mathbf{X} , \mathbf{Y} , and \mathbf{Z} be defined as follows:

$$\mathbf{X} = \begin{bmatrix} 1 & 2 \\ 3 & 4 \end{bmatrix} \quad \mathbf{Y} = \begin{bmatrix} 0 & -1 \\ -3 & 5 \end{bmatrix} \quad \mathbf{Z} = \begin{bmatrix} -1 & 1 & -2 \\ 3 & -5 & 8 \end{bmatrix}.$$

Elements of matrices are indexed by their row first and column second, so $x_{11} = 1$ and $x_{12} = 2$. Addition and scalar multiplication of matrices works in the same way as with vectors, so

$$\mathbf{X} + \mathbf{Y} = \begin{bmatrix} 1+0 & 2-1 \\ 3-3 & 4+5 \end{bmatrix} = \begin{bmatrix} 1 & 1 \\ 0 & 9 \end{bmatrix}$$

and

$$5\mathbf{X} = \begin{bmatrix} 5 & 10 \\ 15 & 20 \end{bmatrix}.$$

The transpose of a matrix \mathbf{A} , written \mathbf{A}^T , is obtained by interchanging the rows and columns, so that if \mathbf{A} is an $m \times n$ matrix, \mathbf{A}^T is an $n \times m$ matrix. As examples we have

$$\mathbf{X}^T = \begin{bmatrix} 1 & 3 \\ 2 & 4 \end{bmatrix} \quad \mathbf{Z}^T = \begin{bmatrix} -1 & 3 \\ 1 & -5 \\ -2 & 8 \end{bmatrix}.$$

Matrix multiplication is somewhat less intuitive. Two matrices can only be multiplied together if the number of columns in the first matrix is the same as the number of rows in the second. (The reason for this will become clear when the operation is defined.) For instance, you can multiply a 2×3 matrix by a 3×3 matrix, but you can't multiply a 2×3 matrix by another 2×3 matrix. This immediately suggests that the order of matrix multiplication is important, since two matrices may have compatible dimensions in one order, but not in the other. Multiplying an $m \times n$ matrix by a $n \times p$ matrix creates an $m \times p$ matrix, defined as follows. Let $\mathbf{A} \in \mathbb{R}^{m \times n}$ and $\mathbf{B} \in \mathbb{R}^{n \times p}$. Then the product $\mathbf{C} = \mathbf{AB} \in \mathbb{R}^{m \times p}$ has elements defined as

$$c_{ij} = \sum_{k=1}^n a_{ik} b_{kj}. \quad (3.17)$$

If you imagine that each row of the first matrix is treated as a vector, and that each column of the second matrix is treated as a vector, then each component of the product matrix is the dot product of a row from the first matrix and a column from the second. For this dot product to make sense, these row

and column vectors have to have the same dimension, that is, the number of columns in the first matrix must equal the number of rows in the second. Using the matrices defined above, we have

$$\mathbf{XY} = \begin{bmatrix} 1 \times 0 + 2 \times -3 & 1 \times -1 + 2 \times 5 \\ 3 \times 0 + 4 \times -3 & 3 \times -1 + 4 \times 5 \end{bmatrix} = \begin{bmatrix} -6 & 9 \\ -12 & 17 \end{bmatrix}.$$

Observe that the element in the first row and first column of the product matrix is the dot product of the first row from \mathbf{X} and the first column from \mathbf{Y} ; the element in the first row and second column of the product is the dot product of the first row from \mathbf{X} and the second column of \mathbf{Y} ; and so forth. You should be able to verify that

$$\mathbf{XZ} = \begin{bmatrix} 5 & -9 & 14 \\ 9 & -17 & 26 \end{bmatrix}.$$

A matrix and a vector can be multiplied together, if you interpret a row vector as an $1 \times n$ matrix or a column vector as a $m \times 1$ matrix. Here, the distinction between row and column vectors is important, because matrices can only be multiplied if their dimensions are compatible. In matrix multiplications, the convention used in this text is that vectors such as \mathbf{x} are column vectors, and a row vector is denoted by \mathbf{x}^T . Again, this distinction is only relevant in matrix multiplication, and for other purposes row and column vectors can be treated interchangeably. *It is worth repeating that matrix multiplication is not commutative, so that usually $\mathbf{XY} \neq \mathbf{YX}$ (and in fact both products may not even exist, depending on the dimensions of \mathbf{X} and \mathbf{Y}), although there are some exceptions.* You may wonder why this seemingly-complicated definition of matrix multiplication is used instead of other, seemingly simpler, approaches. One reason is that this definition actually ends up representing real-world calculations more frequently than other definitions.

A square matrix \mathbf{A} is *symmetric* if $a_{ij} = a_{ji}$ (in other words, for a symmetric matrix $\mathbf{A} = \mathbf{A}^T$), and it is *diagonal* if $a_{ij} = 0$ unless $i = j$ (that is, all its elements are zero except on the diagonal from upper-left to lower-right). A very special diagonal matrix is the *identity matrix*, which has 1's along the diagonal and 0's everywhere else. The notation \mathbf{I} denotes an identity matrix of any size, so we can write

$$\mathbf{I} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \quad \text{or} \quad \mathbf{I} = \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{bmatrix}.$$

In practice, the dimension is usually obvious because of the requirements of matrix multiplication. The identity matrix has the unique property that $\mathbf{AI} = \mathbf{IA} = \mathbf{A}$ for any matrix \mathbf{A} .

A square matrix \mathbf{A} is *invertible* if there is another square matrix (call it \mathbf{A}^{-1}) such that $\mathbf{AA}^{-1} = \mathbf{A}^{-1}\mathbf{A} = \mathbf{I}$. Multiplying by an inverse matrix should be thought of as the equivalent of “matrix division.” Just as it is not possible to divide by all numbers (division by zero is undefined), not all matrices have

inverses. However, the matrices seen in this book will all be invertible. Computing the matrix inverse is a bit tedious, and is rarely needed in transportation network analysis. However one special case is worth mentioning:

Proposition 3.1. *A diagonal matrix \mathbf{A} is invertible if and only if all its diagonal entries are nonzero; in this case, its inverse \mathbf{A}^{-1} is also a diagonal matrix, whose diagonal entries are the reciprocal of the diagonal entries in \mathbf{A} .*

So, for instance,

$$\left(\begin{bmatrix} 5 & 0 \\ 0 & -3 \end{bmatrix} \right)^{-1} = \begin{bmatrix} 1/5 & 0 \\ 0 & -1/3 \end{bmatrix}.$$

Finally, a symmetric matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is *positive definite* if the matrix product $\mathbf{x}^T \mathbf{A} \mathbf{x}$ is strictly positive for any nonzero vector $\mathbf{x} \in \mathbb{R}^n$, and *positive semidefinite* if $\mathbf{x}^T \mathbf{A} \mathbf{x} \geq 0$ for any $\mathbf{x} \in \mathbb{R}^n$ whatsoever. As examples, consider the matrices

$$\mathbf{A} = \begin{bmatrix} 10 & 1 \\ 1 & 5 \end{bmatrix} \quad \mathbf{B} = \begin{bmatrix} 0 & 0 \\ 0 & 1 \end{bmatrix} \quad \mathbf{C} = \begin{bmatrix} 1 & 10 \\ 10 & 5 \end{bmatrix}.$$

The matrix \mathbf{A} is positive definite, because for any vector $\mathbf{x} = [x_1 \ x_2]$, the matrix product is

$$\mathbf{x}^T \mathbf{A} \mathbf{x} = [x_1 \ x_2] \begin{bmatrix} 10 & 1 \\ 1 & 5 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = 10x_1^2 + 2x_1x_2 + 5x_2^2.$$

The expression on the right is always positive, because it can be rewritten as

$$9x_1^2 + 4x_2^2 + (x_1 + x_2)^2.$$

None of those three terms can be negative, and since $\mathbf{x} \neq \mathbf{0}$, at least one of these terms is strictly positive.

The matrix \mathbf{B} is not positive definite, since if $\mathbf{x} = [1 \ 0]$ then $\mathbf{x}^T \mathbf{B} \mathbf{x} = 0$, which is not strictly positive. However, it is positive semidefinite, since the matrix product is

$$\mathbf{x}^T \mathbf{B} \mathbf{x} = [x_1 \ x_2] \begin{bmatrix} 0 & 0 \\ 0 & 1 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = x_2^2$$

which is surely nonnegative.

The matrix \mathbf{C} is neither positive definite nor positive semidefinite, since if $\mathbf{x} = [1 \ -1]$ then

$$\mathbf{x}^T \mathbf{C} \mathbf{x} = [1 \ -1] \begin{bmatrix} 1 & 10 \\ 10 & 5 \end{bmatrix} \begin{bmatrix} 1 \\ -1 \end{bmatrix} = -14 < 0.$$

Notice that for a matrix to be positive definite or semidefinite, we have to check whether a condition holds for all possible nonzero vectors. If the condition fails for even one vector, the matrix is not positive definite or semidefinite.

Checking positive definiteness or semidefiniteness can be tedious. One test is to find all the eigenvalues of the matrix; if they are all positive, the matrix is positive definite, and if all are non-negative it is positive semidefinite. However, diagonal matrices arise fairly often in transportation network analysis, and this case is easy.

Proposition 3.2. *A diagonal matrix is positive definite if and only if all its diagonal entries are strictly positive. A diagonal matrix is positive semidefinite if and only if all its diagonal entries are nonnegative.*

There are also times where we will apply the concept of positive definiteness to non-symmetric matrices. The idea is the same – we want $\mathbf{x}^T \mathbf{A} \mathbf{x}$ to be strictly positive for all nonzero \mathbf{x} – but the eigenvalue test does not apply, and there are non-symmetric matrices which have strictly positive eigenvalues but do not satisfy $\mathbf{x}^T \mathbf{A} \mathbf{x} > 0$ for all nonzero \mathbf{x} . However, we can form the *symmetric part* of the matrix \mathbf{A} by calculating $\frac{1}{2}(\mathbf{A} + \mathbf{A}^T)$. It is easy to show that this is always a symmetric matrix, and that $\mathbf{x}^T \mathbf{A} \mathbf{x} > 0$ if and only if $\mathbf{x}^T (\frac{1}{2}(\mathbf{A} + \mathbf{A}^T)) \mathbf{x} > 0$. So, when we refer to a non-symmetric matrix being positive definite, what we will mean is that its symmetric part $\frac{1}{2}(\mathbf{A} + \mathbf{A}^T)$ is positive definite.

3.1.3 Sets

A *set* is a collection of any type of objects, denoted by a plain capital letter, such as X or Y . In transportation network analysis, we work with sets of numbers, sets of nodes, sets of links, sets of paths, sets of origin-destination pairs, and so forth. Sets can contain either a finite or infinite number of elements. As examples, let's work with the sets

$$X = \{1, 2, 4\} \quad Y = \{1, 3, 5, 7\},$$

using curly braces to denote a set, and commas to list the elements. Set membership is indicated with the notation \in , so $1 \in X$, $1 \in Y$, $2 \in X$, but $2 \notin Y$. The *union* of two sets $X \cup Y$ is the set consisting of elements either in X or in Y (or both), so

$$X \cup Y = \{1, 2, 3, 4, 5, 7\}.$$

The *intersection* of two sets $X \cap Y$ is the set consisting only of elements in both X and Y , so

$$X \cap Y = \{1\}.$$

In this book, sets will take one of three forms:

1. Sets that consist of a finite number of elements, which can be listed as with X and Y above. This includes the set of nodes in a network, the set of acyclic paths in a network, and so forth. For a finite set, the notation $|X|$ indicates the number of elements in X .
2. Sets which are intervals on the real line, such as $[-1, 1]$ or $(0, 2]$. These intervals are sets containing all real numbers between their endpoints; a

square bracket next to an endpoint means that the endpoint is included in the set, while parentheses mean that the endpoint is not included. Intervals usually contain infinitely many elements.

3. Sets which contain all objects satisfying one or more conditions. For instance, the set $\{x \in \mathbb{R} : x^2 < 4\}$ contains all real numbers whose square is less than four; in this case it can be written simply as the interval $(-2, 2)$. A more complicated set is $\{(x, y) \in \mathbb{R}^2 : x + y \leq 3, |x| \leq |y|\}$. This set contains all two-dimensional vectors which (if x and y are the two components of the vector) satisfy both the conditions $x + y \leq 3$ and $|x| \leq |y|$. It can also be thought of as the intersection of the sets $\{(x, y) \in \mathbb{R}^2 : x + y \leq 3\}$ and $\{(x, y) \in \mathbb{R}^2 : |x| \leq |y|\}$. If there are *no* vectors which satisfy all of the conditions, the set is *empty*, denoted \emptyset .

Sets of scalars or vectors can be described in other ways. Given any vector $\mathbf{x} \in \mathbb{R}^n$, the *ball of radius r* is the set

$$B_r(\mathbf{x}) = \{\mathbf{y} \in \mathbb{R}^n : |\mathbf{x} - \mathbf{y}| < r\}, \quad (3.18)$$

that is, the set of all vectors whose distance to \mathbf{x} is less than r , where r is some positive number. A ball in one dimension is an interval; a two-dimensional ball is a circle; a three-dimensional ball is a sphere; a four-dimensional ball a hypersphere, and so on.

Given some set of real numbers X , the vector \mathbf{x} is a *boundary point* of X if *every* ball $B_r(\mathbf{x})$ contains both elements in X and elements not in X , no matter how small the radius r . Notice that the boundary points of a set need not belong to the set: 2 is a boundary point of the interval $(-2, 2)$. A set is *closed* if it contains all of its boundary points. A set X is *bounded* if every element of X is contained in a sufficiently large ball centered at the origin, that is, if there is some r such that $\mathbf{x} \in B_r(\mathbf{0})$ for all $\mathbf{x} \in X$. A set is *compact* if it is both closed and bounded.

These facts will prove useful:

Proposition 3.3. *Let $f(x_1, x_2, \dots, x_n)$ be any linear function, that is, $f(\mathbf{x}) = a_1x_1 + a_2x_2 + \dots + a_nx_n$ for some constants a_i , and let b be any scalar.*

- (a) *The set $\{\mathbf{x} \in \mathbb{R}^n : f(\mathbf{x}) = b\}$ is closed.*
- (b) *The set $\{\mathbf{x} \in \mathbb{R}^n : f(\mathbf{x}) \leq b\}$ is closed.*
- (c) *The set $\{\mathbf{x} \in \mathbb{R}^n : f(\mathbf{x}) \geq b\}$ is closed.*

Proposition 3.4. *Let X and Y be any sets of scalars or vectors.*

- (a) *If X and Y are closed, so are $X \cap Y$ and $X \cup Y$*
- (b) *If X and Y are bounded, so are $X \cap Y$ and $X \cup Y$*
- (c) *If X and Y are compact, so are $X \cap Y$ and $X \cup Y$.*

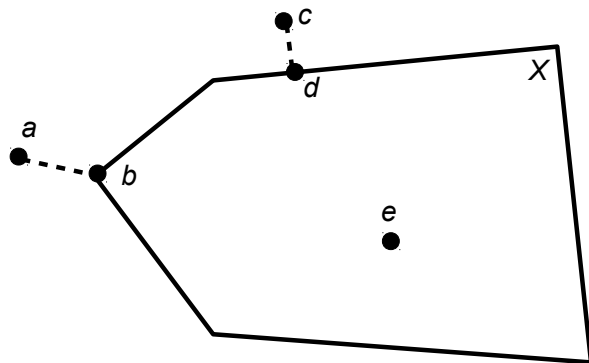


Figure 3.1: Some examples of the projection operation.

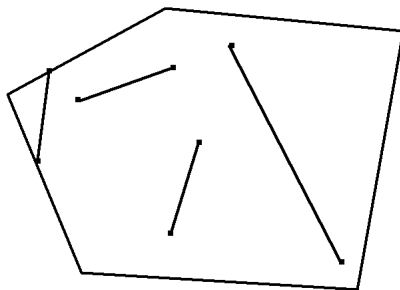


Figure 3.2: A convex set.

Combining Propositions 3.3 and 3.4, we can see that any set defined solely by linear equality or weak inequality constraints (any number of them) is closed.

If X is a closed set of n -dimensional vectors, and $\mathbf{y} \in \mathbb{R}^n$ is any n -dimensional vector, the *projection of \mathbf{y} onto X* , written $\text{proj}_X(\mathbf{y})$, is the vector \mathbf{x} in X which is “closest” to \mathbf{y} in the sense that $|\mathbf{x} - \mathbf{y}|$ is minimal. In Figure 3.1, we have $\text{proj}_X(a) = b$, $\text{proj}_X(c) = d$, and $\text{proj}_X(e) = e$.

This subsection concludes with a discussion of what it means for a set to be *convex*. This notion is very important, and is described in more detail than the other concepts, starting with an intuitive definition. If X is convex, geometrically this means that line segments connecting points of X lie entirely within X . For example, the set in Figure 3.2 is convex, while those in Figure 3.3 and Figure 3.4 are not. Intuitively, a convex set cannot have any “holes” punched into it, or “bites” taken out of it.

Mathematically, we write this as follows:

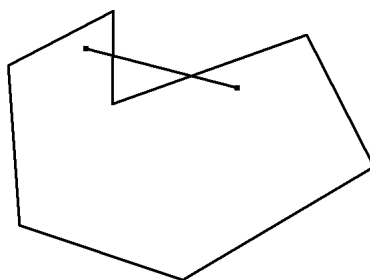


Figure 3.3: A nonconvex set with a “bite” taken out of it.

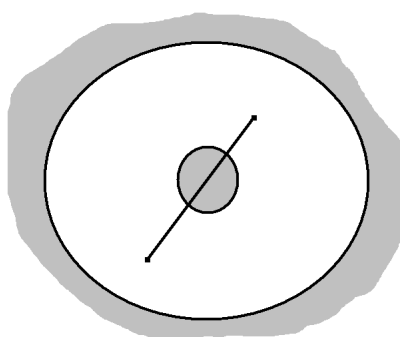


Figure 3.4: A nonconvex set with a “hole” in it.

Definition 3.1. A set $X \subseteq \mathbb{R}^n$ is convex if, for all $x_1, x_2 \in X$ and all $\lambda \in [0, 1]$, the point $\lambda x_2 + (1 - \lambda)x_1 \in X$.

If this definition is not clear, notice that one way to express the line segment between any two points is $\lambda x_2 + (1 - \lambda)x_1$, and that you cover the entire line between x_1 and x_2 as λ varies between 0 and 1, regardless of how close or far apart these two points are located.

Example 3.1. Show that the one-dimensional set $X = \{x : x \geq 0\}$ is convex.

Solution. Pick any $x_1, x_2 \geq 0$ and any $\lambda \in [0, 1]$. Because x_1, x_2 , and λ are all nonnegative, so are λx_2 and $(1 - \lambda)x_1$, and therefore so is $\lambda x_2 + (1 - \lambda)x_1$. Therefore $\lambda x_2 + (1 - \lambda)x_1$ belongs to X as well. ■

Example 3.2. Show that the hyperplane $X = \{\mathbf{x} \in \mathbb{R}^n : \sum_{i=1}^n a_i x_i - b = 0\}$ is convex.

Solution. This set is the same as $\{\mathbf{x} \in \mathbb{R}^n : \sum_{i=1}^n a_i x_i = b\}$. Pick any $\mathbf{x}, \mathbf{y} \in X$ and any $\lambda \in [0, 1]$. Then

$$\begin{aligned} \sum_{i=1}^n a_i (\lambda y_i + (1 - \lambda)x_i) &= \lambda \sum_{i=1}^n a_i y_i + (1 - \lambda) \sum_{i=1}^n a_i x_i \\ &= \lambda b + (1 - \lambda)b \\ &= b \end{aligned}$$

so $\lambda \mathbf{y} + (1 - \lambda)\mathbf{x} \in X$ as well. ■

Sometimes, more complicated arguments are needed.

Example 3.3. Show that the two-dimensional ball $B = \{[x \ y] \in \mathbb{R}^2 : x^2 + y^2 \leq 1\}$ is convex.

Solution. Pick any vectors $\mathbf{a}, \mathbf{b} \in B$ and any $\lambda \in [0, 1]$. We will write the components of these as $\mathbf{a} = [a_x \ a_y]$ and $\mathbf{b} = [b_x \ b_y]$. The point $\lambda \mathbf{b} + (1 - \lambda)\mathbf{a}$ is the vector $[\lambda b_x + (1 - \lambda)a_x \ \lambda b_y + (1 - \lambda)a_y]$. To show that it is in B , we must show that the sum of the squares of these components is no greater than 1.

$$\begin{aligned} &(\lambda b_x + (1 - \lambda)a_x)^2 + (\lambda b_y + (1 - \lambda)a_y)^2 \\ &= \lambda^2(b_x^2 + b_y^2) + (1 - \lambda)^2(a_x^2 + a_y^2) + 2\lambda(1 - \lambda)(a_x b_x + a_y b_y) \\ &\leq \lambda^2 + (1 - \lambda)^2 + 2\lambda(1 - \lambda)(a_x b_x + a_y b_y) \end{aligned}$$

because $\mathbf{a}, \mathbf{b} \in B$ (and therefore $a_x^2 + a_y^2 \leq 1$ and $b_x^2 + b_y^2 \leq 1$). Notice that $a_x b_x + a_y b_y$ is simply the dot product of \mathbf{a} and \mathbf{b} , which is equal to $|\mathbf{a}||\mathbf{b}| \cos \theta$, where θ is the angle between the vectors \mathbf{a} and \mathbf{b} . Since $|\mathbf{a}| \leq 1$, $|\mathbf{b}| \leq 1$ (by definition of B), and since $\cos \theta \leq 1$ regardless of θ , $a_x b_x + a_y b_y \leq 1$. Therefore

$$\begin{aligned} &\lambda^2 + (1 - \lambda)^2 + 2\lambda(1 - \lambda)(a_x b_x + a_y b_y) \\ &\leq \lambda^2 + (1 - \lambda)^2 + 2\lambda(1 - \lambda) = (\lambda + (1 - \lambda))^2 = 1 \end{aligned}$$

so the point $\lambda \mathbf{b} + (1 - \lambda)\mathbf{a}$ is in B regardless of the values of \mathbf{a} , \mathbf{b} , or λ . Thus B is convex. ■

Example 3.4. Show that the set $C = \{[x \ y] \in \mathbb{R}^2 : x^2 + y^2 > 1\}$ is not convex.

Solution. Let $\mathbf{b}^1 = [2 \ 0]$, $\mathbf{b}^2 = [-2 \ 0]$, $\lambda = 1/2$. Then $\lambda \mathbf{b}^2 + (1 - \lambda)\mathbf{b}^1 = [0 \ 0] \notin C$ even though $\mathbf{b}^1, \mathbf{b}^2 \in C$ and $\lambda \in [0, 1]$. Therefore C is not convex. ■

Proving that a set is convex requires showing that something is true for *all* possible values of $x_1, x_2 \in X$, and $\lambda \in [0, 1]$, whereas disproving convexity only requires you to pick one combination of these values where the definition fails. Lastly, this result will prove useful in showing that a set is convex:

Proposition 3.5. If X and Y are convex sets, so is $X \cap Y$.

3.1.4 Functions

A function is a mapping between sets. If the function f maps set X to set Y , then f associates *every* element of X with *some* single element of Y . (Note that not every element of Y needs to be associated with an element of X .) The set X is known as the *domain* of f . Examples include the familiar functions $f(x) = x^2$ and $g(x, y) = x^2 + y^2$. The function f maps \mathbb{R} to \mathbb{R} , while g maps \mathbb{R}^2 to \mathbb{R} . An example of a function which maps \mathbb{R}^2 to \mathbb{R}^2 is the vector-valued function

$$\mathbf{h}(x_1, x_2) = \begin{bmatrix} 3x_1 + 2x_2 \\ -x_1x_2 \end{bmatrix}$$

The inverse of a function f , denoted f^{-1} , “undoes” the mapping f in the sense that if $f(x) = y$, then $f^{-1}(y) = x$. As an example, if $f(x) = x^3$, then $f^{-1}(x) = \sqrt[3]{x}$. Not every function has an inverse, and inverse functions may need to be restricted to subsets of X and Y . The composition of two functions f and g , denoted $f \circ g$ involves substituting function g into function f . If $f(x) = x^3$ and $g(x) = 2x + 4$, then $f \circ g(x) = (2x + 4)^3$. This can also be written as the function $f(g(x))$.

A function is *continuous* if, for any point $\hat{x} \in X$, the limit $\lim_{x \rightarrow \hat{x}} f(x)$ exists and is equal to $f(\hat{x})$. Intuitively, continuous functions can be drawn without lifting your pen from the paper. A function is *differentiable* if, for any point $x \in X$, the limit

$$\lim_{\Delta x \rightarrow 0} \frac{f(x + \Delta x) - f(x)}{|\Delta x|} \quad (3.19)$$

exists; that limit is then called the *derivative* of f , and gives the slope of the tangent line at any point. It can be shown that any differentiable function must be continuous.

Proposition 3.6. Let f and g be continuous functions. Then we have the following:

- The multiple αf is a continuous function for any scalar α .
- The sum $f + g$ is a continuous function.
- The product fg is a continuous function.
- The composition $f \circ g$ is a continuous function.

Furthermore, all of these results hold if “continuous” is replaced by “differentiable.”

Proposition 3.7. *Let X be a nonempty, convex set of n -dimensional vectors. Then the projection function $\text{proj}_X(\mathbf{x})$ is defined and continuous for all $\mathbf{x} \in \mathbb{R}^n$.*

Differentiability is more complicated when dealing with functions of multiple variables (that is, functions whose domain is \mathbb{R}^2 , \mathbb{R}^3 , and so forth). The basic notion is the *partial derivative*, in which all variables except one are assumed constant, and an ordinary derivative is taken with respect to the remaining variable. For instance, if $f(x, y) = x^2 + 3xy + y^3$, the partial derivatives with respect to x and y are

$$\frac{\partial f}{\partial x} = 2x + 3y \quad \frac{\partial f}{\partial y} = 3x + 3y^2.$$

Second partial derivatives are found in the same way, taking partial derivatives of $\frac{\partial f}{\partial x}$ and $\frac{\partial f}{\partial y}$. If all of the second partial derivatives of a function are continuous at a point, the order of differentiation does not matter and

$$\frac{\partial^2 f}{\partial x \partial y} = \frac{\partial^2 f}{\partial y \partial x} \quad (3.20)$$

at that point.

All of these partial derivatives can be organized into different kinds of vectors and matrices. The *gradient* of a function is the vector of all its partial derivatives. For $f(x, y) = x^2 + 3xy + y^3$, this is

$$\nabla f(x, y) = \begin{bmatrix} \frac{\partial f}{\partial x} \\ \frac{\partial f}{\partial y} \end{bmatrix} = \begin{bmatrix} 2x + 3y \\ 3x + 3y^2 \end{bmatrix}.$$

A point \mathbf{x} where $\nabla f(\mathbf{x}) = \mathbf{0}$ is called a *stationary point*. (If f is a function of a single variable, this is just a point where the derivative vanishes.) The *Hessian* of a function is the matrix of all of its second partial derivatives. For this same function, we have

$$Hf(x, y) = \begin{bmatrix} \frac{\partial^2 f}{\partial x^2} & \frac{\partial^2 f}{\partial x \partial y} \\ \frac{\partial^2 f}{\partial y \partial x} & \frac{\partial^2 f}{\partial y^2} \end{bmatrix} = \begin{bmatrix} 2 & 3 \\ 3 & 6y \end{bmatrix}$$

By equation (3.20), if the Hessian matrix contains continuous functions, it is symmetric.

Now consider a vector-valued function of multiple variables $f : \mathbb{R}^n \rightarrow \mathbb{R}^m$. (That is, f has n input variables and produces an m -dimensional vector as output.) The *Jacobian* of f is a matrix of all of its *first* partial derivatives. For a concrete example, let

$$\mathbf{g}(x, y, z) = \begin{bmatrix} g_1(x, y, z) \\ g_2(x, y, z) \end{bmatrix} = \begin{bmatrix} x^2 + y^2 + z^2 \\ xyz \end{bmatrix}.$$

The Jacobian is the matrix

$$J\mathbf{g}(x, y, z) = \begin{bmatrix} \frac{\partial g_1}{\partial x} & \frac{\partial g_1}{\partial y} & \frac{\partial g_1}{\partial z} \\ \frac{\partial g_2}{\partial x} & \frac{\partial g_2}{\partial y} & \frac{\partial g_2}{\partial z} \end{bmatrix} = \begin{bmatrix} 2x & 2y & 2z \\ yz & xz & xy \end{bmatrix}$$

Be sure to note the difference between Hessians and Jacobians. The Hessian is defined for scalar-valued functions, and contains all the second partial derivatives. The Jacobian is defined for vector-valued functions, and contains all the first partial derivatives.

Finally, there is an important notion of function *convexity*. Confusingly, this is a different idea than set convexity discussed in Section 3.1.3, although there are some relationships and similar ideas between them. Set and function convexity together play a pivotal role in optimization and network equilibrium problems, so function convexity is discussed at length here. Geometrically, a convex function lies below its secant lines. Remember that a secant line is the line segment joining two points on the function. As we see in Figure 3.5, no matter what two points we pick, the function always lies below its secant line. On the other hand, in Figure 3.6, not every secant line lies above the function: some lie below it, and some lie both above and below it. Even though we can draw some secant lines which are above the function, this isn't enough: *every* possible secant must lie above the function. For this concept to make sense, the domain X of the function must be a convex set, an assumption which applies for the remainder of this section

The following definition makes this intuitive notion formal:

Definition 3.2. A function $f : X \rightarrow \mathbb{R}$ is convex if, for every $x_1, x_2 \in X$ and every $\lambda \in [0, 1]$,

$$f((1 - \lambda)x_1 + \lambda x_2) \leq (1 - \lambda)f(x_1) + \lambda f(x_2) \quad (3.21)$$

and strictly convex if

$$f((1 - \lambda)x_1 + \lambda x_2) < (1 - \lambda)f(x_1) + \lambda f(x_2) \quad (3.22)$$

for all distinct $x_1, x_2 \in X, \lambda \in (0, 1)$

Essentially, x_1 and x_2 are the two endpoints for the secant line. Since this entire line segment must be above the function, we need to consider every point between x_1 and x_2 . This is what λ does: as λ varies between 0 and 1, the points $\lambda x_2 + (1 - \lambda)x_1$ cover every point between x_1 and x_2 . You can think of

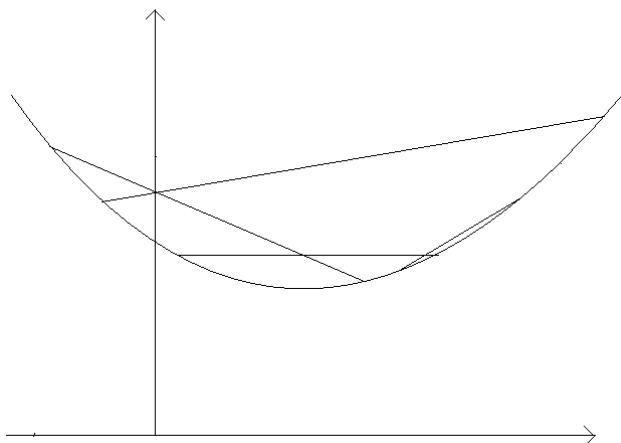


Figure 3.5: A convex function lies below all of its secants.

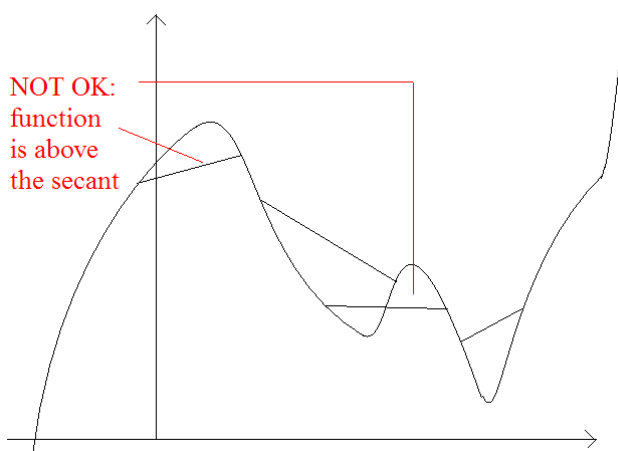


Figure 3.6: A nonconvex function does not lie below all of its secants.

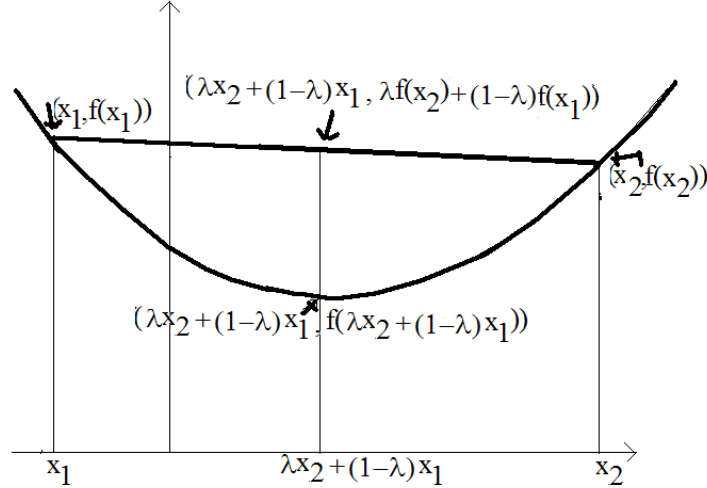


Figure 3.7: All of the relevant points for the definition of convexity.

$\lambda x_2 + (1 - \lambda)x_1$ as a “weighted average,” where λ is the weight put on x_2 . For $\lambda = 0$, all the weight is on x_1 . For $\lambda = 1$, all the weight is on x_2 . For $\lambda = 1/2$, equal weight is put on the two points, so the weighted average is the midpoint. $\lambda = 1/3$ corresponds to the point a third of the way between x_1 and x_2 .

So we need to say that, at all such intermediate x values, the value of f is lower than the y -coordinate of the secant. The value of the function at this point is simply $f((1 - \lambda)x_1 + \lambda x_2)$. Because the secant is a straight line, its y -coordinate can be seen as a weighted average of the y -coordinates of its endpoints, that is, $f(x_1)$ and $f(x_2)$. This weighted average can be written as $(1 - \lambda)f(x_1) + \lambda f(x_2)$, so requiring the function to lie below the secant line is exactly the same as enforcing condition (3.21) for all possible secant lines: that is, for all $x_1, x_2 \in X$ and all $\lambda \in [0, 1]$.

Figure 3.7 explains this in more detail. Along the horizontal axis, the secant endpoints x_1 and x_2 are shown, along with an intermediate point $\lambda x_2 + (1 - \lambda)x_1$. The y -coordinates are also shown: at the endpoints, these are $f(x_1)$ and $f(x_2)$. At the intermediate point, the y -coordinate of the function is $f(\lambda x_2 + (1 - \lambda)x_1)$, while the y -coordinate of the secant is $\lambda f(x_2) + (1 - \lambda)f(x_1)$. Because the function is convex, the former can be no bigger than the latter. Time spent studying this diagram is very well spent. Make sure you understand what each of the four points marked on the diagram represents, and why the given mathematical expressions correctly describe these points. Make sure you see what role λ plays: as λ increases from 0 to 1, the central vertical line moves from x_1 to x_2 . (What would happen if we picked x_1 and x_2 such that $x_1 > x_2$? What if $x_1 = x_2$?)

Example 3.5. Is the function $f(x) = |x|$, $x \in \mathbb{R}$ convex? Is it strictly convex?

Solution. To see if f is convex, we need to see if (3.21) is true; to see if it is strictly convex, we need to check (3.22). Furthermore, these inequalities have to be true for *every* $x_1, x_2 \in \mathbb{R}$, and *every* $\lambda \in [0, 1]$. It is not enough to simply pick a few values randomly and check the equations. So, we have to work symbolically. In this case,

$$\begin{aligned} f((1-\lambda)x_1 + \lambda x_2) &= |(1-\lambda)x_1 + \lambda x_2| \\ &\leq |(1-\lambda)x_1| + |\lambda x_2| \quad \text{by the triangle inequality} \\ &= (1-\lambda)|x_1| + \lambda|x_2| \quad \text{because } \lambda, 1-\lambda \geq 0 \\ &= (1-\lambda)f(x_1) + \lambda f(x_2) \end{aligned}$$

Therefore (3.21) is satisfied, so f is convex. To show that it is strictly convex, we would have to show that the inequality

$$|(1-\lambda)x_1 + \lambda x_2| < |(1-\lambda)x_1| + |\lambda x_2|$$

can be replaced by a *strict* inequality $<$. However, we can't do this: for example, if $x_1 = 1$, $x_2 = 2$, $\lambda = 0.5$, the left side of the inequality ($|1/2 + 2/2| = 3/2$) is exactly equal to the right side ($|1/2| + |2/2| = 3/2$). So f is not strictly convex. ■

Note that proving that $f(x)$ is convex requires a *general* argument, where proving that $f(x)$ was not strictly convex only required a single counterexample. This is because the definition of convexity is a “for all” or “for every” type of argument. To prove convexity, you need an argument that allows for all possible values of x_1 , x_2 , and λ , whereas to disprove it you only need to give one set of values where the necessary condition doesn't hold.

Example 3.6. Show that every linear function $f(x) = ax + b$, $x \in \mathbb{R}$ is convex, but not strictly convex.

Solution.

$$\begin{aligned} f((1-\lambda)x_1 + \lambda x_2) &= a((1-\lambda)x_1 + \lambda x_2) + b \\ &= a((1-\lambda)x_1 + \lambda x_2) + ((1-\lambda) + \lambda)b \\ &= (1-\lambda)(ax_1 + b) + \lambda(ax_2 + b) \\ &= (1-\lambda)f(x_1) + \lambda f(x_2) \end{aligned}$$

So we see that inequality (3.21) is in fact satisfied as an *equality*. That's fine, so every linear function is convex. However, this means we can't replace the inequality \leq with the strict inequality $<$, so linear functions are not strictly convex. ■

Sometimes it takes a little bit more work, as in the following example:

Example 3.7. Show that $f(x) = x^2$, $x \in \mathbb{R}$ is strictly convex.

Solution. Pick x_1, x_2 so that $x_1 \neq x_2$, and pick $\lambda \in (0, 1)$.

$$\begin{aligned} f((1-\lambda)x_1 + \lambda x_2) &= ((1-\lambda)x_1 + \lambda x_2)^2 \\ &= (1-\lambda)^2 x_1^2 + \lambda^2 x_2^2 + 2(1-\lambda)\lambda x_1 x_2 \end{aligned}$$

What to do from here? Since $x_1 \neq x_2$, $(x_1 - x_2)^2 > 0$. Expanding, this means that $x_1^2 + x_2^2 > 2x_1 x_2$. This means that

$$\begin{aligned} (1-\lambda)^2 x_1^2 + \lambda^2 x_2^2 + 2(1-\lambda)\lambda x_1 x_2 &< (1-\lambda)^2 x_1^2 + \lambda^2 x_2^2 + (1-\lambda)(\lambda)(x_1^2 + x_2^2) \\ &= (1-\lambda)x_1^2 + \lambda x_2^2 \text{ (after some algebra)} \\ &= (1-\lambda)f(x_1) + \lambda f(x_2) \end{aligned}$$

which proves strict convexity. ■

This last example shows that proving convexity can be difficult and nonintuitive, even for simple functions like x^2 . The good news is that there are often simpler conditions that we can check. These conditions involve the first and second derivatives of a function.

Proposition 3.8. *Let $f : X \rightarrow \mathbb{R}$ be a differentiable function, where X is a subset of \mathbb{R} . Then f is convex if and only if*

$$f(x_2) \geq f(x_1) + f'(x_1)(x_2 - x_1)$$

for all $x_1, x_2 \in X$.

Proposition 3.9. *Let $f : X \rightarrow \mathbb{R}$ be twice differentiable, where X is a subset of \mathbb{R} , and let f be twice differentiable on X . Then f is convex if and only if $f''(x) \geq 0$ for all $x \in X$*

Equivalent conditions for strict convexity can be obtained in a natural way, changing \geq to $>$ and requiring that x_1 and x_2 be distinct in Proposition 3.8. Proposition 3.9 also changes slightly in this case; $f''(x) > 0$ is sufficient for strict convexity but is not necessary. You are asked to prove these statements in the exercises. Essentially, Proposition 3.8 says that f lies *above* its tangent lines (Figure 3.8), while Proposition 3.9 says that f is always “curving upward.” (A convex function lies *above* its tangents, but *below* its secants.)

These conditions are usually easier to verify than that of Definition 3.2.

Example 3.8. *Show that $f(x) = x^2$ is strictly convex using Proposition 3.8*

Solution. Pick any $x_1, x_2 \in \mathbb{R}$ with $x_1 \neq x_2$. We have $f'(x_1) = 2x_1$, so we need to show that

$$x_2^2 > x_1^2 + 2x_1(x_2 - x_1)$$

Expanding the right-hand side and rearranging terms, we see this is equivalent to

$$x_1^2 - 2x_1 x_2 + x_2^2 > 0$$

or

$$(x_1 - x_2)^2 > 0$$

which is clearly true since $x_1 \neq x_2$. Thus f is strictly convex. ■

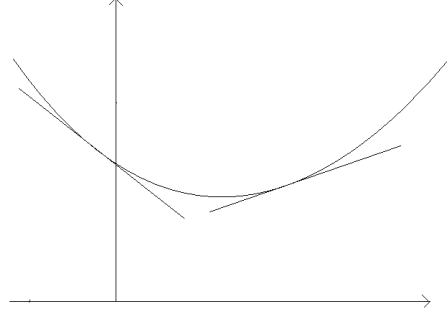


Figure 3.8: A convex function lies above its tangents.

Example 3.9. Show that $f(x) = x^2$ is strictly convex using Proposition 3.9

Solution. $f''(x) = 2 > 0$ for all $x \in \mathbb{R}$, so f is strictly convex. ■

Much simpler! If f is differentiable (or, better yet, twice differentiable) checking these conditions is almost always easier.

Furthermore, once we know that some functions are convex, we can use this to show that many other combinations of these functions must be convex as well.

Proposition 3.10. If f and g are convex functions, and α and β are positive real numbers, then $\alpha f + \beta g$ is convex as well.

Proposition 3.11. If f and g are convex functions, then $f \circ g$ is convex as well.

Some common convex functions are $|x|$, x^2 , e^x , and $ax + b$. So, Proposition 3.10 tells us that $3x^2 + 4|x|$ is convex. It also tells us that any quadratic function $ax^2 + bx + c$ is convex as long as $a > 0$. Proposition 3.11 says that the composition of two convex functions is convex as well. For instance, e^{x^2} is convex, and $x^4 = (x^2)^2$ is convex as well.

What about functions of more than one variable? The “shortcut” conditions in Propositions 3.8 and 3.9 only apply if the domain of f is one-dimensional. It turns out that very similar conditions can be given for the multivariable case. The multi-dimensional equivalent of the first derivative is the gradient, and the equivalent of the second derivative is the Hessian.

The equivalent conditions on convexity are

Proposition 3.12. Let $f : X \rightarrow \mathbb{R}$ be a function whose gradient exists everywhere on X . Then f is convex if and only if

$$f(\mathbf{x}_2) \geq f(\mathbf{x}_1) + \nabla f(\mathbf{x}_1) \cdot (\mathbf{x}_2 - \mathbf{x}_1)$$

for all $\mathbf{x}_1, \mathbf{x}_2 \in X$.

Proposition 3.13. *Let $f : X \rightarrow \mathbb{R}$ be a function whose Hessian exists everywhere on X . Then f is convex if and only if $H(f)$ is positive semidefinite for all $\mathbf{x} \in X$; and f is strictly convex if $H(f)$ is positive definite for all $\mathbf{x} \in X$.*

Unfortunately, neither of these is as easy to check as the single-dimension equivalents. In particular, it is rather tedious to check whether or not a matrix is positive semidefinite or not. Fortunately, in many important cases in transportation network analysis, the Hessian is diagonal and Proposition 3.2 applies.

As a final note, one useful link between convex sets and convex functions is the following result:

Proposition 3.14. *If $g : \mathbb{R}^n \rightarrow \mathbb{R}$ is a convex function, then the set $X = \{x : g(x) \leq 0\}$ is a convex set.*

3.2 Tools for Modeling Equilibrium

Chapter 1 characterized equilibrium as a “consistent” state in which no driver can improve his or her satisfaction by unilaterally changing routes. Under the assumption of continuous flow variables, we can use calculus to greatly simplify the problem. This section explores the equilibrium concept further, looking at three different ways to formulate this principle mathematically. These are the *fixed point* and *variational inequality* approaches. A third approach, *convex optimization*, is undertaken in Chapter 4. Each approach has its own advantages from the standpoint of this book. The fixed point formulation is perhaps the most intuitive and generally-applicable definition, but does not give any indication as to how one might actually find this equilibrium. The variational inequality formulation lends itself to physical intuition and can also accommodate a number of variations on the equilibrium problem. The convex optimization approach provides an intuitive interpretation of solution methods and provides an elegant proof of equilibrium uniqueness in link flows, but the connection between the equilibrium concept and optimization requires more mathematical explanation and is less apparent at first glance.

3.2.1 Fixed-point problems

A **fixed point** of a function f is a value x such that $f(x) = x$, that is, the value x is unchanged by f . As an example, the function $f_1(x) = 2x - 1$ has only one fixed point at $x = 1$, because $f_1(1) = 2 \times 1 - 1 = 1$, the function $f_2(x) = x^2$ has two fixed points at 0 and 1 ($0^2 = 0$ and $1^2 = 1$), while the function $f_3(x) = x^2 + 1$ has no fixed points at all. A helpful visual illustration is that the fixed points of a function are the points of intersection between the function’s graph and the 45-degree line $y = x$ (Figure 3.9).

Equilibrium solutions can often be formulated as fixed points of a suitable function. For instance, in the traffic assignment problem, a route choice model and a congestion model are mutually interdependent: drivers choose routes to

avoid congestion, but congestion is determined by the routes drivers choose. (Figure 1.3). An equilibrium solution is consistent in the sense that drivers are satisfied with the travel times calculated by the paths they chose. Feeding the route choices into the congestion model, then feeding the resulting travel times into the route choice model, one can obtain the original route choices back. This is reminiscent of fixed point problems: when you evaluate a function at its fixed-point, after performing whatever calculations the function requires you obtain the fixed-point again. Fixed points thus arise naturally when dealing with these kinds of “circular” dependencies.

As a first example, consider the problem of trying to estimate the number of bus riders in a dense urban area. The bus system is subject to congestion; when there are x riders, the average delay to customers is given by the function $t = T(x)$ which is assumed continuous and increasing. (Assuming that the fleet of buses and timetables are fixed, more riders mean longer boarding and offloading times, more frequent stops, and the possibility of full buses skipping waiting passengers.) However, the number of bus riders depends on the congestion in the system — as the buses become more crowded, some riders may switch to alternate modes of transportation or combine trips, so we can write $x = X(t)$ for some function X which is continuous and decreasing. For a concrete example, assume that the system is designed such that $T(x) = 10 + x$, and that the ridership function is given by $X(t) = [20 - t]^+$, when x and t are measured in appropriate units (say, thousands of passengers and minutes).²

The goal is to find the ridership x ; but $x = X(t)$ and $t = T(x)$, which means that we need to find some value of x such that $x = X(T(x))$. This is a fixed point problem! Here the function f is the composition of X and T :

$$\begin{aligned} x &= X(T(x)) \\ &= [20 - T(x)]^+ \\ &= [20 - (10 + x)]^+ \\ &= [10 - x]^+ \end{aligned}$$

Assuming that $10 - x$ is nonnegative, we can replace the right-hand side of the last equation by $10 - x$. Solving the resulting equation $x = 10 - x$, we obtain $x = 5$ as the fixed point. Checking, when there are 5 riders the average travel time will be 15 minutes (based on T); and when the travel time is 15 minutes, there will indeed be 5 riders (based on X). Finally, if $x = 5$, our assumption that $10 - x$ was nonnegative was true so this solution is valid.

For this example, it is relatively easy to find the fixed point by substituting the definitions of the functions and performing some algebra. However, in more complex problems it will be difficult or impossible to calculate the fixed point directly. Despite this, fixed points are important because of so-called “fixed-point theorems” which guarantee the existence of a fixed point under certain

²The notation $[.]^+$ is used to mean the *positive part* of the term in brackets, that is, $[y]^+ = \max\{y, 0\}$. If the term in brackets is negative, it is replaced by zero; otherwise it is unchanged.

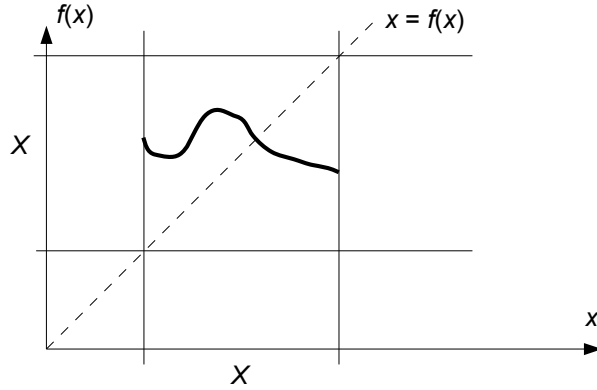


Figure 3.9: Visualizing fixed-point problems; the intersection of $f(x)$ and the 45-degree line is a point where $x = f(x)$.

conditions of the function f . A fixed point theorem by Brouwer is provided below, and another by Kakutani is provided in Section 6.1.1. These fixed point theorems are *non-constructive* because they give us no indication of how to find a fixed point, they simply guarantee that at least one exists. Brouwer's theorem can be stated as:

Theorem 3.1. (Brouwer). *Let f be a continuous function from the set K to itself, where K is convex and compact. Then there is at least one point $x \in K$ such that $f(x) = x$.*

The exercises ask you to show that each of these conditions is necessary; you might find it helpful to visualize these conditions geometrically similar to Figure 3.9.

Let us apply Brouwer's theorem to the transit ridership example. Both X and T are continuous functions, so their composition $X \circ T$ is continuous as well. (Alternately, by substituting one function into the other we obtain $X(T(x)) = [10 - x]^+$, which is evidently continuous.) What are the domain and range of $X(T(x))$? Since $X(T(x))$ is the positive part of $20 - T(x)$, then $x = X(T(x)) \geq 0$. Further note that because $x \geq 0$, $T(x) \geq 10$, so $x = X(T(x)) \leq 10$. That is, we have shown that x must lie between 0 and 10, so the function $X(T(x)) = [10 - x]^+$ can be defined from the set $[0, 10]$ to itself. This set is convex and compact, so Brouwer's theorem would have told us that at least one fixed point must exist, even if we didn't know how to find it.

3.2.2 Variational inequalities

Fixed point problems often lend themselves to elegant theorems like Brouwer's, which prove that a fixed point must exist. However, such problems often lack easy solution methods. The variational inequality can be more convenient to

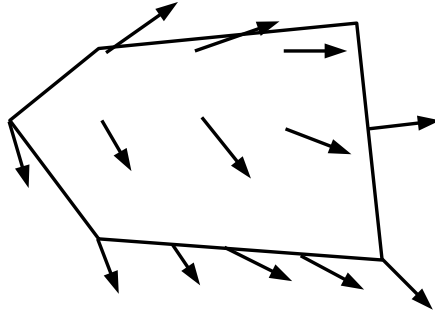


Figure 3.10: A container with a force field.

work with in this respect. Variational inequalities can be motivated with a physical analogy. Imagine an object, initially stationary, which is confined to move within some frictionless container (Figure 3.10) and cannot leave. This object is acted on by a force whose magnitude and direction can be different at each point. If the object is in the interior of the container, the object will begin to move in the same direction as the force at that point. If the object starts at the edge of the container, it may not be able to move in the same direction as the force, but it might slide along the side of the object. The variational inequality problem is to determine where in the container (if anywhere) the object can be placed so that it will not move under the action of the force field.

Figure 3.11 shows some examples. In this figure, the direction of the force is drawn with black arrows, shown only at the points under consideration for clarity. At three of the points (A, B, and C) the point will move: at A in the direction of the force, and at B and C sliding along the edge of the container in the general direction of the force. At the other two points (D and E), the object will not move under the action of the force, being effectively resisted by the container wall. So, D and E are solutions to the variational inequality problem created by the container shape and force field, while A, B, and C are not.

A little thought should convince you that (1) if the object is on the boundary of the container, but not at a corner point, it will be unmoved if and only if the force is perpendicular to the boundary (point D in Figure 3.11), and (2) if the object is at a corner of the container, it will be unmoved if and only if the force makes a right or obtuse angle with all of the boundary directions (point E). These two cases can be combined together: a point on the boundary is an equilibrium if and only if the force makes a right or obtuse angle with all boundary directions. In fact, if the force makes such an angle with all boundary directions, it will do so with any other direction pointing into the feasible set (Figure 3.12). So, we see that *a point is a solution to the variational inequality if and only if the direction of the force at that point makes a right or obtuse angle with any possible direction the object could move in.*

The mathematical definition of a variational inequality is little more than

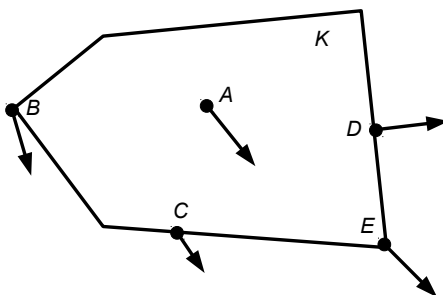


Figure 3.11: Two solutions, and three non-solutions, to a variational inequality.

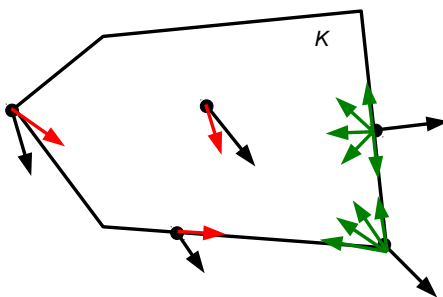


Figure 3.12: Stable points (in green) make an obtuse or right angle with all feasible directions; unstable points (in red) make an acute angle with some feasible directions.

translating the above physical problem into algebraic terminology. The “container” is replaced by a set K of n -dimensional vectors, which for our purposes can be assumed compact and convex (as in all of the figures so far). The “force field” is replaced by a vector-valued function $\mathbf{F} : \mathbb{R}^n \rightarrow \mathbb{R}^n$ which depends on n variables and produces an n -dimensional vector as a result. Recalling vector operations, specifically equation (3.16), saying that two vectors make a right or obtuse angle is equivalent to saying that their dot product is negative. Therefore, rewriting the condition in the previous paragraph with this mathematical notation, we have the following definition:

Definition 3.3. *Given a convex set $K \subseteq \mathbb{R}^n$ and a function $\mathbf{F} : K \rightarrow \mathbb{R}^n$, we say that the vector $\hat{\mathbf{x}} \in K$ solves the $\text{VI}(K, \mathbf{F})$ if, for all $\mathbf{x} \in K$, we have*

$$\mathbf{F}(\hat{\mathbf{x}}) \cdot (\mathbf{x} - \hat{\mathbf{x}}) \leq 0. \quad (3.23)$$

In other words, as \mathbf{x} ranges over all possible points in K , $\mathbf{x} - \hat{\mathbf{x}}$ represents all possible directions the object can move from $\hat{\mathbf{x}}$. The point $\hat{\mathbf{x}}$ solves the variational inequality precisely when the dot product of the force and all possible directions is negative.

There is a relationship between fixed point problems and variational inequalities, which can be motivated again by the physical analogy of a force acting within a container. The solutions to the variational inequality are quite literally “fixed points” in the sense that an object placed there will not move. Consider some point \mathbf{x} under the action of the force \mathbf{F} . Assume furthermore that this force is constant and does not change magnitude or direction as this point moves. Then, since K is convex, the trajectory of the object can be identified with the curve $\text{proj}_K(\mathbf{x} + \alpha \mathbf{F}(\mathbf{x}))$ where α ranges over all positive numbers and proj_K means projection onto the set K as defined in Section 3.1.3. See Figure 3.13 for a few examples. If a point is a solution to $\text{VI}(K, \mathbf{F})$, then the corresponding “trajectory” will simply be the same point no matter what α is. So, for the sake of convenience we arbitrarily choose $\alpha = 1$ and look at the location of the point $\text{proj}_K(\mathbf{x} + \mathbf{F}(\mathbf{x}))$. If this is the same as the initial point \mathbf{x} , then \mathbf{x} is a solution to the variational inequality. So, finally, if we let

$$\mathbf{f}(\mathbf{x}) = \text{proj}_K(\mathbf{x} + \mathbf{F}(\mathbf{x})), \quad (3.24)$$

then the fixed points of \mathbf{f} coincide exactly with solutions to $\text{VI}(K, \mathbf{F})$.

In many cases of practical interest, \mathbf{F} will be a continuous function. Furthermore, it can be shown that the projection mapping onto a convex set (such as K) is a well-defined (i.e., single-valued), continuous function. Then by Proposition 3.6, the function \mathbf{f} defined by equation (3.24) is a continuous function. So, if the set K is compact in addition to being convex, then Brouwer’s theorem shows that the variational inequality must have at least one solution:

Theorem 3.2. *If $K \subseteq \mathbb{R}^n$ is a compact, convex set and $\mathbf{F} : \mathbb{R}^n \rightarrow \mathbb{R}^n$ is a continuous function, then the variational inequality $\text{VI}(K, \mathbf{F})$ has at least one solution.*

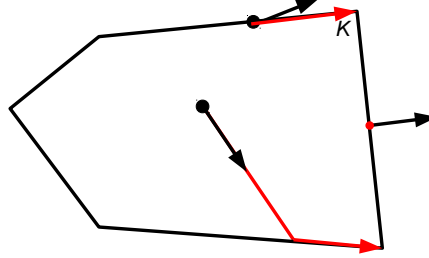


Figure 3.13: Trajectories of different points (red) assuming the force (black) is constant.

You should convince yourself that all of these conditions are necessary: if the container K is not bounded or not closed, or if the force field \mathbf{F} is not continuous, then it is possible that an object placed at any point in the container will move under the action of the force field.

3.3 Exercises

1. [22] For this exercise, let $x_{11} = 5$, $x_{12} = 6$, $x_{13} = 7$, $x_{21} = 4$, $x_{22} = 3$, and $x_{23} = 9$. Evaluate each of these sums.
 - (a) $\sum_{i=1}^2 \sum_{j=1}^3 x_{ij}$
 - (b) $\sum_{i=1}^2 \sum_{j=2}^3 x_{ij}$
 - (c) $\sum_{j=1}^3 x_{1j}$
 - (d) $\sum_{j=1}^2 \sum_{i=1}^3 x_{ji}$
 - (e) $\sum_{j=1}^3 \sum_{k=0}^1 x_{k+1,j}$
 - (f) $\sum_{i=1}^2 \sum_{j=i}^3 x_{ij}$
2. [22] Repeat Exercise 1, but with products \prod instead of sums \sum .
3. [33] Section 3.1.1 lists three properties of the summation notation \sum . Formulate and prove analogous properties for the product notation \prod . You can assume that the product involves finitely many factors.
4. [22] Prove Proposition 3.1.
5. [24] Prove Proposition 3.2.
6. [53] Prove that if a matrix is invertible, its inverse matrix is unique.

7. [26] For each of the sets below, identify its boundary points and indicate whether or not it is closed, whether or not it is bounded, and whether or not it is convex.
 - (a) $(5, 10]$
 - (b) $[4, 6]$
 - (c) $(0, \infty)$
 - (d) $\{x : |x| > 5\}$
 - (e) $\{x : |x| \leq 5\}$
 - (f) $\{(x, y) : 0 \leq x \leq 4, -3 \leq y \leq 3\}$
 - (g) $\{(x, y) : 4x - y = 1\}$
8. [34] Prove Proposition 3.3.
9. [34] Prove Proposition 3.4.
10. [53] Identify the projections of the following points on the corresponding sets. You may find it helpful to draw sketches.
 - (a) The point $x = 3$ on the set $[0, 1]$.
 - (b) The point $x = \frac{1}{2}$ on the set $[0, 1]$.
 - (c) The point $(2, 5)$ on the unit circle $x^2 + y^2 = 1$.
 - (d) The point $(2, 5)$ on the line $x + y = 1$.
 - (e) The point $(2, 5)$ on the line segment between $(0, 1)$ and $(1, 0)$.
 - (f) The point $(2, 3)$ on the line segment between $(0, 1)$ and $(1, 0)$.
 - (g) The point $(1, 2, 3)$ on the sphere $x^2 + y^2 + z^2 = 2$.
11. [20] Prove Proposition 3.5.
12. [23] Find the inverses of the following functions.
 - (a) $f(x) = 3x + 4$
 - (b) $f(x) = 5e^x - 1$
 - (c) $f(x) = 1/x$
13. [57] Prove Proposition 3.6 from first principles, using the definitions of continuity and differentiability in the text.
14. [65] Prove Proposition 3.7, using the formal definition of continuity.
15. [22] Calculate the gradients of the following functions:
 - (a) $f(x_1, x_2, x_3) = x_1^2 + 2x_2x_3 + x_3^2$
 - (b) $f(x, y) = \frac{xy}{x^2 + y^2 + 1}$
 - (c) $f(x_1, x_2, d_1, d_2) = 3x_1 + 5x_2^2 + 3(d_1 + 4)^2 + 5(d_2 - 1)^3$

16. [25] Calculate the Jacobian matrices of the following functions:

(a) $f(x, y) = \begin{bmatrix} y \\ -x \end{bmatrix}$

(b) $f(x_1, x_2, x_3) = \begin{bmatrix} x_1^2 + x_2 + x_3 \\ 3x_2 - 5x_3 \end{bmatrix}$

(c) $f(x, y) = \begin{bmatrix} 3x + y \\ 3y + x \\ xy \end{bmatrix}$

17. [26] Calculate the Hessian matrices of functions from Exercise 15.

18. [35] Determine which of the following functions are convex, strictly convex, or neither. Justify your answer rigorously.

(a) $f(x) = \log x$ where $x \geq 1$.

(b) $f(x) = -x^3$ where $x \leq 0$.

(c) $f(x) = \sin x$ where $x \in \mathbb{R}$.

(d) $f(x) = e^{x^2-2x-3}$ where $x \in \mathbb{R}$.

(e) $f(x) = \begin{cases} 0 & \text{if } x \leq 0 \\ x & \text{if } x \geq 0 \end{cases}$

19. [35] Determine which of the following sets is convex. Justify your answer rigorously.

(a) $X = \{(x_1, x_2) \in \mathbb{R}^2 : 4x_1 - 3x_2 \leq 0\}$

(b) $X = \{x \in \mathbb{R} : e^x \leq 4\}$

(c) $X = \{(x_1, x_2, x_3) \in \mathbb{R}^3 : x_1 = 0\}$

(d) $X = \{(x_1, x_2, x_3) \in \mathbb{R}^3 : x_1 \neq 0\}$

(e) $X = \{1\}$

20. [31] Show that the integral of an increasing function is a convex function.

21. [54] Show that a differentiable function f of a single variable is convex if and only if $f(x) + f'(x)(y - x) \leq f(y)$ for all x and y .

22. [54] Show that a twice-differentiable function f of a single variable is convex if and only if f'' is everywhere nonnegative.

23. [68] Prove the claims of Exercises 21 and 22 for functions of multiple variables.

24. [32] Page 70 states that “ $f''(x) > 0$ is sufficient for strict convexity but not necessary.” Give an example of a strictly convex function where $f''(x) = 0$ at some point.

25. [31] For each of the following functions, find all of its fixed points or state that none exist.
- (a) $f(x) = x^2$, where $X = \mathbb{R}$
 - (b) $f(x) = 1 - x^2$, where $X = [0, 1]$
 - (c) $f(x) = e^x$, where $X = \mathbb{R}$.
 - (d) $f(x_1, x_2) = \begin{bmatrix} -x_1 \\ x_2 \end{bmatrix}$ where $X = \{(x_1, x_2) : -1 \leq x_1 \leq 1, 0 \leq x_2 \leq 1\}$.
 - (e) $f(x, y) = \begin{bmatrix} -y \\ x \end{bmatrix}$ where $X = \{(x, y) : x^2 + y^2 \leq 1\}$ is the unit disc.
26. [54] Brouwer's theorem guarantees the existence of a fixed point for the function $f : X \rightarrow X$ if f is continuous and X is closed, bounded, and convex. Show that each of these four conditions is necessary by creating examples of a function f and set X which satisfy only three of those conditions but do not have a fixed point. Come up with such examples with each of the four conditions missing. Hint: you will probably find it easiest to work with simple functions and sets whenever possible, e.g. something like $X = [0, 1]$. Visualizing fixed points as intersections with the diagonal line through the origin may help you as well.
27. [45] Find all of the solutions of each of the following variational inequalities $\text{VI}(K, \mathbf{F})$.
- (a) $F(x) = x + 1$, $K = [0, 1]$
 - (b) $F(x) = \frac{x}{2}$, $K = [-1, 1]$
 - (c) $F(x, y) = \begin{bmatrix} 0 \\ -y \end{bmatrix}$, $K = \{(x, y) : x^2 + y^2 \leq 1\}$
28. [58] Theorem 3.2 guarantees the existence of a solution to the variational inequality $\text{VI}(K, \mathbf{F})$ if K is closed, bounded, and convex, and if \mathbf{F} is continuous. Show that each of these three conditions is necessary by creating counterexamples of functions \mathbf{F} and sets K which satisfy only three of these conditions, but for which $\text{VI}(K, \mathbf{F})$ has no solution. It may be helpful to include sketches.

Chapter 4

Optimization

Optimization is a broad and powerful tool used in many areas of engineering and business. This chapter discusses optimization techniques as they can be applied to transportation network problems and, like the previous chapter, makes no pretensions about being comprehensive. Together with fixed point problems and variational inequalities, discussed in the previous chapter, the language of convex optimization forms the third and most powerful way to represent equilibrium problems. However, this connection is deeper and less intuitive. The chapter begins by introducing the elements of an optimization problem in Section 4.1, along with some examples of different kinds of optimization problems in Section 4.2. Section 4.3 bridges the previous chapter to the current one, providing some useful results on minimizing convex functions defined on convex sets. Section 4.4 presents a few techniques for solving simple optimization problems, including line search methods, and the use of Lagrange multipliers. Finally, the optional Section 4.5 presents two heuristics which can be used to approximately solve almost any kind of optimization problem, and are common for solving certain logistics and network problems.

4.1 Introduction to Optimization

Every optimization problem has three components: an objective function, decision variables, and constraints. When one talks about *formulating* an optimization problem, it means translating a “real-world” problem into the mathematical equations and variables which comprise these three components.

The objective function, often¹ denoted f or z , reflects a single quantity to be either maximized or minimized. Examples in the transportation world include “minimize congestion”, “maximize safety”, “maximize accessibility”, “minimize cost”, “maximize pavement quality”, “minimize emissions,” “maxi-

¹But not always! Often problem-specific notation will be used, such as c for cost, and so forth. The notation given in this section is what is traditionally used if we are referring to a generic optimization problem outside of a specific context.

mize revenue,” and so forth. You may object to the use of a *single* objective function, since real-world problems typically involve many different and conflicting objectives from different stakeholders, and this objection is certainly valid. There are several reasons why the scope of this book is restricted to single objectives. From a historical perspective, if we don’t know how to optimize a single objective function, then we have no hope of being able to optimize multiple objectives simultaneously, since the latter build on the former. From a pedagogical perspective, the methods of multi-objective optimization are much more complex, and the basic concepts of optimization are best learned in a simpler context first. From a mathematical perspective, the definition of “simultaneously optimize” is very tricky, since there is probably no plan which will, say, simultaneously minimize congestion and agency cost — therefore multiobjective optimization is “fuzzier,” and this fuzziness can be confusing when first taught. So, don’t be alarmed by this restriction to a single objective, but do keep it in the back of your head if you use optimization beyond this book.

The decision variables (often, but not always, denoted as the vector \mathbf{x}) reflect aspects of the problem that you (or the decision maker) can control. This can include both variables you can directly choose, as well as variables which you indirectly influence by the choice of other decision variables. For example, if you are a private toll road operator trying to maximize your profit, you can directly choose the toll, so that is a decision variable. You can’t directly choose how many people drive on the road — but because that’s influenced by the toll you chose, the toll road volume should be a decision variable as well. However, you should avoid including extraneous decision variables. Every decision variable in your formulation should either directly influence the objective function, or influence another decision variable that affects the objective function.

Constraints represent any kind of limitation on the values that the decision variables can take. The most intuitive types of constraints are those which directly and obviously limit the choices you can make: you can’t exceed a budget, you are required by law to provide a certain standard of maintenance, you are not allowed to change the toll by more than \$1 from its current value, and so forth. One very frequent mistake is to omit “obvious” constraints (e.g., the toll can’t be negative). An optimization formulation must be complete and not leave out any constraint, no matter how obvious it may seem to you. Real-world optimization problems are solved by computers, for which nothing is “obvious.” The second type of constraint is required to ensure consistency among the decision variables. Following the toll road example from the previous paragraph, while you can influence both the toll (directly) and the roadway volume (indirectly), the roadway volume must be consistent with the toll you chose. This relationship must be reflected by an appropriate constraint; say, a demand function $d(\tau)$ giving the demand d for the tollway in terms of the toll τ , which can be estimated in a variety of econometric ways. Constraints are the primary way these linkages between decision variables can be captured.

This section concludes with some notation for optimization problems in general, when we don’t want to refer to a specific problem context but instead make blanket statements about larger classes of problems. The most generic specifi-

cation for an optimization problem is to denote all the decision variables with the vector \mathbf{x} (whose dimension is equal to the number of decision variables), the objective function as $f(\mathbf{x})$, and the set of all decision variables which satisfy all the constraints as X . This latter set is called the *feasible set* or *feasible region*, and a vector $\mathbf{x} \in X$ is called a *feasible solution* or simply *feasible*. A feasible solution $\hat{\mathbf{x}} \in X$ is said to be a *global minimum* of $f(\mathbf{x})$ if $f(\hat{\mathbf{x}}) \leq f(\mathbf{x})$ for all feasible \mathbf{x} , and a *global maximum* of $f(\mathbf{x})$ if $f(\hat{\mathbf{x}}) \geq f(\mathbf{x})$. A global minimum for a minimization problem, or a global maximum for a maximization problem, are called *optimal solutions* or simply *optima*. In general, optima do not always exist; and if they exist they may not be unique. Can you think of some of these cases?

The following results are very useful in showing cases when two superficially different optimization problems may in fact be the same. The first result shows that it is easy to convert any maximization problem to a minimization problem (or vice versa) by negating the objective function. The second shows that constants may be freely added or subtracted to objective functions without changing the optimal solutions. The third shows that multiplication or division by a positive constant does not change the optimal solutions.

Proposition 4.1. *If the feasible set is X , the solution $\hat{\mathbf{x}}$ is a global maximum of f if and only if $\hat{\mathbf{x}}$ is a global minimum of $-f$ for the same feasible set X .*

Proof. For the first part, assume that $\hat{\mathbf{x}}$ is the global maximum of f on the set X . Then $f(\hat{\mathbf{x}}) \geq f(\mathbf{x})$ for all $\mathbf{x} \in X$. Multiplying both sides by -1 reverses the sign of the inequality, giving $-f(\hat{\mathbf{x}}) \leq -f(\mathbf{x})$ for all $\mathbf{x} \in X$, so $\hat{\mathbf{x}}$ is a global minimum for $-f$. For the second part, assume that $\hat{\mathbf{x}}$ is the global minimum of $-f$, so $-f(\hat{\mathbf{x}}) \leq -f(\mathbf{x})$ for all $\mathbf{x} \in X$. Again multiplying both sides by -1 gives $f(\hat{\mathbf{x}}) \geq f(\mathbf{x})$ for all $\mathbf{x} \in X$, so $\hat{\mathbf{x}}$ is a global maximum of f . \square

This result is convenient, because we do not need to develop two different theories for maximization and minimization problems. Instead, we can focus just on one. In this book, we develop results for minimization problems, a fairly common convention in engineering.² Whenever you encounter a maximization problem, you can convert it to a minimization problem by negating the objective function, then using the minimization results and procedures. (Of course, this choice is completely arbitrary; we could just as well have chosen to develop results only for maximization problems and in fact some other fields do just this.)

Proposition 4.2. *Let $\hat{\mathbf{x}}$ be an optimal solution when the objective function is $f(\mathbf{x})$ and the feasible set is X . Then for any constant b not depending on \mathbf{x} , $\hat{\mathbf{x}}$ is also an optimal solution when the objective function is $f(\mathbf{x}) + b$ and the feasible set is X .*

Proof. By Proposition 4.1, with no loss of generality we can assume that the optimization problem is a minimization problem and that $\hat{\mathbf{x}}$ is a global minimum.

²Economists often use maximization problems as their standard convention.

Therefore $f(\hat{\mathbf{x}}) \leq f(\mathbf{x})$ for all $x \in X$. We can add the constant b to both sides of the equation, so $f(\hat{\mathbf{x}}) + b \leq f(\mathbf{x}) + b$ for all $x \in X$. Therefore $\hat{\mathbf{x}}$ is also a global minimum (and thus optimal) when the objective function is $f(\mathbf{x}) + b$. \square

Proposition 4.3. *Let $\hat{\mathbf{x}}$ be an optimal solution when the objective function is $f(\mathbf{x})$ and the feasible set is X . Then for any constant $c > 0$ not depending on \mathbf{x} , $\hat{\mathbf{x}}$ is also an optimal solution when the objective function is $cf(\mathbf{x})$ and the feasible set is X .*

Proof. See Exercise 6. \square

4.2 Examples of Optimization Problems

The discussion to this point may be a bit abstract. The best way to demonstrate these concepts is through a few tangible examples. These examples will also demonstrate, *en passant*, the stylistic conventions for writing optimization problems. In each of these examples, the objective function, decision variables, and constraints will first be described intuitively, using plain English. Only then will we introduce mathematical notation, equations, inequalities, and so forth. You should make this a habit, no matter how simple the problem may appear at first glance.

This first example may be familiar to you from introductory calculus:

Example 4.1. *You have 60 feet of fence available, and wish to enclose the largest rectangular area possible. What dimensions should you choose for the fenced-off area?*

Solution. The objective is clear from the problem statement: you wish to maximize the area enclosed by the fence. The decision variables are not directly given in the problem. Rather, you are told that you must enclose a rectangular area. To determine a rectangle, you need to make two decisions: its length and its width. These are both decision variables you can control directly, and there are no indirect decision variables because the length and width directly determine its area. There is one obvious constraint — the perimeter of the fence cannot exceed 60 feet — and two less obvious ones: the length and width must be nonnegative. Since the length and width are independent of each other (the perimetric constraint notwithstanding), there is no need to add a “consistency constraint” linking them.

Now we can proceed to introduce mathematical notation. Starting with the decision variables, let L represent the length and W the width. Then the objective function is to maximize the area, which is equal to LW . The perimeter of a rectangle with length L and width W is $2L + 2W$, so the perimeter constraint can be expressed $2L + 2W \leq 60$. Finally, the nonnegativity constraints are $L \geq 0$ and $W \geq 0$. All of this information is concisely represented in the following way:

$$\begin{array}{ll}
\max_{L,W} & LW \\
\text{s.t.} & 2L + 2W \leq 60 \\
& L \geq 0 \\
& W \geq 0
\end{array}$$

Note that the decision variables are indicated underneath the word “max” (which would naturally be replaced by “min” if this was a minimization problem), and that each constraint is listed on a separate row underneath the objective function. The letters “s.t.” stand for “subject to.” ■

Example 4.2. (Transit frequency setting.) *You are working for a public transit agency in a city, and must decide the frequency of service on each of the bus routes. The bus routes are known and cannot change, but you can change how the city’s bus fleet is allocated to each of these routes. (The more buses assigned to a route, the higher the frequency of service.) Knowing the ridership on each route, how should buses be allocated to routes to minimize the total waiting time?*

Solution. To formulate this as an optimization problem, we need to identify an objective, decision variables, and constraints. For this problem, as we do this we will be faced with other assumptions which must be made. With this and other such problems, there may be more than one way to write down an optimization problem, what matters is that we *clearly state all of the assumptions made*. Another good guiding principle is to *start with the simplest model which captures the important behavior*, which can then be refined by relaxing assumptions or replacing simple assumptions with more realistic ones.

We need some notation, so let R be the set of bus routes. We know the current ridership on each route d_r . Here we will make our first assumption: that the demand on each route is inelastic and will not change based on the service frequency. Like all assumptions, it is not entirely true but in some cases may be close enough to the truth to get useful results; if not, you should think about what you would need to replace this assumption, which is always a good exercise. We are also given the current fleet size (which we will denote N), and must choose the number of buses associated with each route (call this n_r) — thus the decision variables in this problem are the n_r values.

The objective is to minimize the total waiting time, which is the sum of the waiting time for the passengers on each route. How long must passengers wait for a bus? If we assume that travelers arrive at a uniform rate, then the average waiting time will be half of the service headway. How is the headway related to the number of buses on the route n_r ? Assuming that the buses are evenly dispersed throughout the time period we are modeling, and assuming that each bus is always in use, then the headway on route r will be the time required to traverse this route (T_r) divided by the number of buses n_r assigned to this route. So, the average delay per passenger is half of the headway, or $T_r/(2n_r)$, and the total passenger delay on this route is $(d_r T_r)/(2n_r)$. This leads us to the

objective function

$$D(\mathbf{n}) = \sum_{r \in R} \frac{d_r T_r}{2n_r} \quad (4.1)$$

in which the total delay is calculated by summing the delay associated with each route.

What constraints do we have? Surely we must run at least one bus on each route (or else we would essentially be canceling a route), and in reality as a matter of policy there may be some lower limit on the number of buses assigned to each route; for route r , call this lower bound L_r . Likewise, there is some upper bound U_r on the number of buses assigned to each route as well. So, we can introduce the constraint $L_r \leq n_r \leq U_r$ for each route r .

Putting all of these together, we have the optimization problem

$$\begin{aligned} \min_{\mathbf{n}} \quad & D(\mathbf{n}) = \sum_{r \in R} \frac{d_r T_r}{2n_r} \\ \text{s.t.} \quad & n_r \geq L_r \quad \forall r \in R \\ & n_r \leq U_r \quad \forall r \in R \\ & \sum_{r \in R} n_r \leq N \end{aligned}$$

■

Example 4.3. (Scheduling maintenance.) *You are responsible for scheduling routine maintenance on a set of transportation facilities (such as pavement sections or bridges.) The state of these facilities can be described by a condition index which ranges from 0 to 100. Each facility deteriorates at a known, constant rate (which may differ between facilities). If you perform maintenance during a given year, its condition will improve by a known amount. Given an annual budget for your agency, when and where you should perform maintenance to maximize the average condition of these facilities? You have a 10 year planning horizon.*

Solution. In contrast to the previous example, where the three components of the optimization problem were described independently, from here on problems will be formulated in a more organic way, describing a model built from the ground up. (This is how optimization models are usually described in practice.) After describing the model in this way, we will identify the objective function, decision variables, and constraints to write the optimization problem in the usual form. We start by introducing notation based on the problem statement.

Let F be the set of facilities, and let c_f^t be the condition of facility f at the **end**³ of year t , where t ranges from 1 to 10. Let d_f be the annual deterioration on facility f , and i_f the amount by which the condition will improve if maintenance is performed. So, if no maintenance is performed during year t , then

$$c_f^t = c_f^{t-1} - d_f \quad \forall f \in F, t \in \{1, 2, \dots, 10\} \quad (4.2)$$

³This word is intentionally emphasized. In this kind of problem it is very easy to get confused about what occurs at the start of period t , at the end of period t , during the middle of period t , etc.

and if maintenance is performed we have

$$c_f^t = c_f^{t-1} - d_f + i_f \quad \forall f \in F, t \in \{1, 2, \dots, 10\} \quad (4.3)$$

Both of these cases can be captured in one equation with the following trick: let x_f^t equal one if maintenance is performed on facility f during year t , and 0 if not. Then

$$c_f^t = c_f^{t-1} - d_f + x_f^t i_f \quad \forall f \in F, t \in \{1, 2, \dots, 10\} \quad (4.4)$$

Finally, the condition can never exceed 100 or fall below 0, so the full equation for the evolution of the state is

$$c_f^t = \begin{cases} 100 & \text{if } c_f^{t-1} - d_f + x_f^t i_f > 100 \\ 0 & \text{if } c_f^{t-1} - d_f + x_f^t i_f < 0 \\ c_f^{t-1} - d_f + x_f^t i_f & \text{otherwise} \end{cases} \quad (4.5)$$

for all $f \in F$ and $t \in \{1, 2, \dots, 10\}$. Of course, for this to be usable we need to know the current conditions of the facilities, c_f^0 .

The annual budget can be represented this way: let k_f be the cost of performing maintenance on facility f , and B^t the budget available in year t . Then

$$\sum_{f \in F} k_f x_f^t \leq B^t \quad \forall t \in \{1, \dots, 10\}. \quad (4.6)$$

For the objective, we need the average condition across all facilities and all years; this is simply $\frac{1}{10|F|} \sum_{f \in F} \sum_{t=1}^{10} c_f^t$. The obvious decision variables are the maintenance variables x_f^t , but we also have to include c_f^t because these are influenced by the maintenance variables. As constraints, we need to include the state evolution equations (4.5), the budget constraints (4.6), and, less obviously the requirement that x_f^t be either 0 or 1. Putting it all together, we have the optimization problem

$$\begin{aligned} \max_{\mathbf{x}, \mathbf{c}} \quad & \frac{1}{10|F|} \sum_{f \in F} \sum_{t=1}^{10} c_f^t \\ \text{s.t.} \quad & \sum_{f \in F} k_f x_f^t \leq B^t \quad \forall t \in \{1, \dots, 10\} \\ & c_f^t \text{ is given by (4.5)} \quad \forall f \in F, t \in \{1, 2, \dots, 10\} \\ & x_f^t \in \{0, 1\} \quad \forall f \in F, t \in \{1, 2, \dots, 10\} \end{aligned}$$

■

In this example, pay close attention to the use of formulas like (4.6), which show up very frequently in optimization. It is important to make sure that every “index” variable in the formula is accounted for in some way. Equation (4.6) involves the variables x_f^t , but for which facilities f and time periods t ? The facility index f is summed over, while the time index t is shown at right as $\forall t \in \{1, \dots, 10\}$. This means that a copy of (4.6) exists for each time period, and in each of these copies, the left-hand side involves a sum over all facilities at that time. Therefore the one line (4.6) actually includes 10 constraints, one

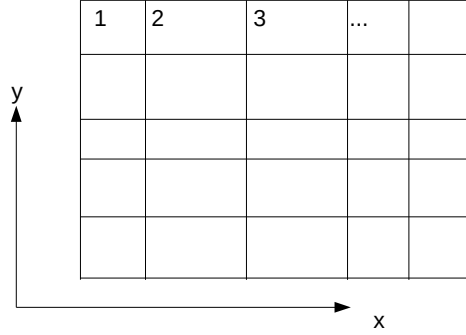


Figure 4.1: Coordinates and labeling of intersections for Example 4.4

for each year. It is common to forget to include things like $\forall t \in \{1, \dots, 10\}$, or to try to use an index of summation outside of the sum as well (e.g., an expression like $B_f - \sum_{f \in F} k_f x_f^t$), which is meaningless. Make sure that all of your indices are properly accounted for!

In this next example, the objective function is less obvious.

Example 4.4. (A facility location problem.) *In a city with a grid network, you need to decide where to locate three bus terminals. Building the terminal at different locations costs a different amount of money. Knowing the home locations of customers throughout the city who want to use the bus service, where should the terminals be located to minimize the construction cost and walking distance customers have to walk? Assume that each customer will walk from their home location to the nearest terminal.*

Solution. This problem could easily become very complicated if we take into account the impact of terminal locations on bus routes, so let's focus on what the problem is asking for: simply locating terminals to minimize walking distance from customers' home locations.

Number each of the intersections in the grid network (Figure 4.1) from 1 to I , the total number of intersections. Assume that the terminals will be located at these intersections, and let the variables L_1 , L_2 , and L_3 denote the numbers of the intersections where terminals will be built. Let $C(i)$ be the cost of building a terminal at location i , so the total cost of construction is $C(L_1) + C(L_2) + C(L_3)$. Let P be the set of customers, and let H_p denote the intersection that is the home location of customer p .

How can we calculate the walking distance between two intersections (say, i and j)? Figure 4.1 shows a coordinate system superimposed on the grid. Let $x(i)$ and $y(i)$ be the coordinates of intersection i in this system. Then the walking distance between points i and j is

$$d(i, j) = |x(i) - x(j)| + |y(i) - y(j)|. \quad (4.7)$$

This is often called the *Manhattan distance* between two points, after one of the densest grid networks in the world.

So what is the walking distance $D(p)$ for customer p ? The distance from p to the first terminal is $d(H_p, L_1)$, to the second terminal is $d(H_p, L_2)$, and to the third is $d(H_p, L_3)$. The passenger will walk to whichever is closest, so $D(p, L_1, L_2, L_3) = \min\{d(H_p, L_1), d(H_p, L_2), d(H_p, L_3)\}$ and the total walking distance is $\sum_{p \in P} D(p, L_1, L_2, L_3)$.

For this problem, the decision variables and constraints are straightforward: the only decision variables are L_1 , L_2 , and L_3 and the only constraint is that these need to be integers between 1 and I . The tricky part is the objective function: we are instructed both to minimize total cost as well as total walking distance. We have equations for each of these, but we can only have one objective function. In these cases, it is common to form a *convex combination* of the two objectives, introducing a weighting parameter $\lambda \in [0, 1]$. That is, let

$$f(L_1, L_2, L_3) = \lambda[C(L_1) + C(L_2) + C(L_3)] + (1 - \lambda) \left[\sum_{p \in P} D(p, L_1, L_2, L_3) \right] \quad (4.8)$$

Look at what happens as λ varies. If $\lambda = 1$, then the objective function reduces to simply minimizing the cost of construction. If $\lambda = 0$, the objective function is simply minimizing the total walking distance. For a value in between 0 and 1, the objective function is a weighted combination of these two objectives, where λ indicates how important the cost of construction is relative to the walking distance.

For concreteness, the optimization problem is

$$\begin{array}{ll} \min_{L_1, L_2, L_3} & \lambda[C(L_1) + C(L_2) + C(L_3)] + (1 - \lambda) \left[\sum_{p \in P} D(p, L_1, L_2, L_3) \right] \\ \text{s.t.} & L_f \in \{1, 2, \dots, I\} \quad \forall f \in \{1, 2, 3\} \end{array}$$

■

In the final example in this section, finding a mathematical representation of a solution is more challenging.

Example 4.5. (Shortest path problem.) *Figure 4.2 shows a road network, with the origin and destination marked. Given the travel time on each roadway link, what is the fastest route connecting the origin to the destination?*

Solution. We've already presented some algorithms for solving this problem in Section 2.4, but here we show how it can be placed into the general framework of optimization problems.

Notation first: let the set of nodes be N , and let r and s represent the origin and destination nodes. Let the set of links be A , and let t_{ij} be the travel time on link (i, j) . So far, so good, but how do we represent a route connecting two intersections?

Following Example 4.3, introduce binary variables $x_{ij} \in \{0, 1\}$, where $x_{ij} = 1$ if link (i, j) is part of the route, and $x_{ij} = 0$ if link (i, j) is not part of the route.

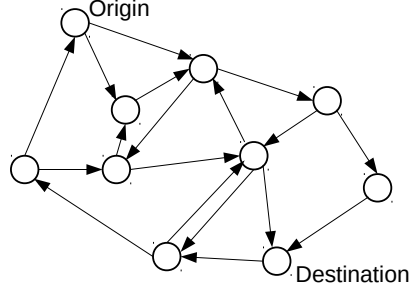


Figure 4.2: Roadway network for Example 4.5

The travel time of a route is simply the sums of the travel times of the links in the route, which is $\sum_{(i,j) \in A} t_{ij}x_{ij}$.

We now have an objective function and decision variables, but what of the constraints? Besides the trivial constraint $x_{ij} \in \{0,1\}$, we need constraints which require that the x_{ij} values actually form a contiguous path which starts at the origin r and ends at the destination s . We do this by introducing a *flow conservation constraint* at each intersection. For node i , recall that $\Gamma(i)$ denotes the set of links which leave node i , and $\Gamma^{-1}(i)$ denotes the set of links which enter node i .

Consider any contiguous path connecting intersection r and s , and examine any node i . One of four cases must hold:

Case I: Node i does not lie on the path at all. Then all of the x_{ij} values should be zero for $(i,j) \in \Gamma(i) \cup \Gamma^{-1}(i)$.

Case II: Node i lies on the path, but is not the origin or destination. Then $x_{ij} = 1$ for exactly one $(i,j) \in \Gamma(i)$, and for exactly one $(i,j) \in \Gamma^{-1}(i)$.

Case III: Node i is the origin r . Then all x_{ij} values should be zero for $(i,j) \in \Gamma^{-1}(i)$, and $x_{ij} = 1$ for exactly one $(i,j) \in \Gamma(i)$.

Case IV: Node i is the destination s . Then all x_{ij} values should be zero for $(i,j) \in \Gamma(i)$, and $x_{ij} = 1$ for exactly one $(i,j) \in \Gamma^{-1}(i)$.

An elegant way to combine these cases is to look at the differences $\sum_{(i,j) \in \Gamma(i)} x_{ij} - \sum_{(h,i) \in \Gamma^{-1}(i)} x_{hi}$. For cases I and II, this difference will be 0; for case III, the difference will be +1, and for case IV, the difference will be -1.

So, this leads us to the optimization formulation of the shortest path prob-

lem:

$$\begin{aligned}
 \min_{\mathbf{x}} \quad & \sum_{(i,j) \in A} t_{ij} x_{ij} \\
 \text{s.t.} \quad & \sum_{(i,j) \in \Gamma(i)} x_{ij} - \sum_{(h,i) \in \Gamma^{-1}(i)} x_{hi} = \begin{cases} 1 & \text{if } i = r \\ -1 & \text{if } i = s \\ 0 & \text{otherwise} \end{cases} \quad \forall i \in N \\
 & x_{ij} \in \{0, 1\} \quad \forall (i, j) \in A
 \end{aligned} \tag{4.9}$$

■

A careful reader will notice that if the four cases are satisfied for a solution, then the equations (4.9) are satisfied, but the reverse may not be true. Can you see why, and is that a problem?

There is often more than one way to formulate a problem: for instance, we might choose to minimize congestion by spending money on capacity improvements, subject to a budget constraint. Or, we might try to minimize the amount of money spent, subject to a maximum acceptable limit on congestion. Choosing the “correct” formulation in this case may be based on which of the two constraints is harder to adjust (is the budget constraint more fixed, or the upper limit on congestion?). Still, you may be troubled by this seeming imprecision. This is one way in which modeling is “as much art as science,” which is not surprising — since the human and political decision-making processes optimization tries to formalize, along with the inherent value judgments, (what truly is the objective?) are not as precise as they seem on the surface. One hallmark of a mature practitioner of optimization is a flexibility with different formulations of the same underlying problem, and a willingness to engage in “back-and-forth” with the decision maker as together you identify the best formulation for a particular scenario. In fact, in some cases it may not matter. The theory of duality (which is beyond the scope of this book) shows that these alternate formulations sometimes lead to the same ultimate decision, which is comforting.

4.3 Convex Optimization

In general, it is difficult to find the optimum value with a nonlinear objective function. For instance, the function in Figure 4.3 has many local minima and is unbounded below as $x \rightarrow -\infty$, both of which can cause serious problems if we’re trying to minimize this function. Usually, the best that a software program can do is find a local minimum. If it finds one of the local minima for this function, it may not know if there is a better one somewhere else (or if there is, how to find it). Or if it starts seeking x values which are negative, we could run into the unbounded part of this function.

On the other hand, some functions are very easy to minimize. The function in Figure 4.4 only has one minimum point, is not unbounded below, and there are many algorithms which can find that minimum point efficiently.

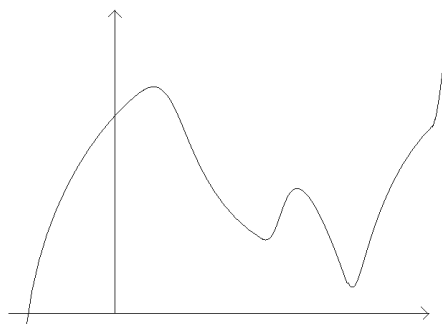


Figure 4.3: A function which is not convex.

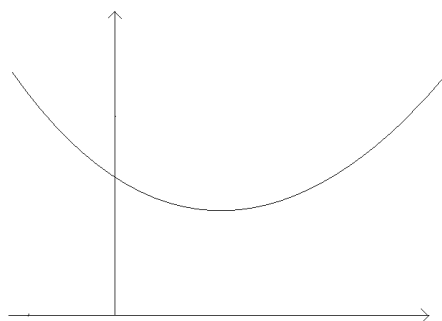


Figure 4.4: A convex function.

What distinguishes these is a property called *convexity*, which was defined in Chapter 3. If the feasible region is a convex set, and if the objective function is a convex function, then it is much easier to find the optimal solution. Checking convexity of the objective function is not usually too difficult. To check convexity of the feasible region, the following result is often useful.

Theorem 4.1. *Consider an optimization program whose constraints all take the form $g_i(\mathbf{x}) \leq 0$ or $h_j(\mathbf{x}) = 0$, where $i = 1, 2, \dots, k$ indexes the inequality constraints and $j = 1, 2, \dots, \ell$ indexes the equality constraints. If each function $g_1(\mathbf{x}), \dots, g_k(\mathbf{x})$ is convex, and if each function $h_1(\mathbf{x}), \dots, h_\ell(\mathbf{x})$ is linear, then the feasible region $X = \{\mathbf{x} \in \mathbb{R}^n : g_i(\mathbf{x}) \leq 0, h_j(\mathbf{x}) = 0, i \in \{1, \dots, k\}, j \in \{1, \dots, \ell\}\}$ is a convex set.*

Proof. Let $Y_i = \{\mathbf{x} \in \mathbb{R}^n : g_i(\mathbf{x}) \leq 0\}$ represent the values of \mathbf{x} which satisfy the i -th inequality constraint, and let $Z_j = \{\mathbf{x} \in \mathbb{R}^n : h_j(\mathbf{x}) = 0\}$ be the values of \mathbf{x} which satisfy the j -th equality constraint. From Proposition 3.14, all of the sets Y_i are convex. From Example 3.2, all of the sets Z_j are convex. The feasible region X is the set of vectors \mathbf{x} which satisfy *all* of the inequality and equality constraints, that is, the intersection of all of the sets Y_i and Z_j . By Proposition 3.5, therefore, X is convex. \square

This is a very common situation, where the functions representing the inequality constraints are convex, and the functions representing equality constraints are linear. From this theorem, this means that the feasible region must be convex.

This subsection collects a few useful results justifying why we're spending some time on convexity: minimizing a convex function over a feasible region that is a convex set is extremely advantageous. Every local minimum is a global minimum, every stationary point is a local minimum, and the set of global minima is a convex set. Furthermore, if the function is strictly convex, the global minimum is unique. In other words, unlike in elementary calculus, you don't have to perform any "second derivative tests" on solutions to ensure they are truly minima, or distinguish between local and global minima. The set of global minima being a convex set is useful because it means that all solutions are "connected" or "adjacent" in some sense — there are no far-flung optimal solutions.

First, a few definitions; in everything that follows, we are trying to minimize a function $f(\mathbf{x})$ over a feasible region X . (That is, X is the set of all \mathbf{x} which satisfy all of the constraints.)

Definition 4.1. *The point \mathbf{x} is a local minimum of f if $\mathbf{x} \in X$, and if there exists $\epsilon > 0$ such that $f(\mathbf{x}) \leq f(\mathbf{x}')$ for all $\mathbf{x}' \in X$ such that $\|\mathbf{x} - \mathbf{x}'\| < \epsilon$.*

Proposition 4.4. *If f is a convex function and X is a convex set, then every local minimum of f is also a global minimum.*

Proof. By contradiction, assume that \mathbf{x}_1 is a local minimum of f , but not a global minimum. Then there is some $\mathbf{x}_2 \in X$ such that $f(\mathbf{x}_2) < f(\mathbf{x}_1)$. Because

f is a convex function, for all $\lambda \in (0, 1]$ we have

$$f((1 - \lambda)\mathbf{x}_1 + \lambda\mathbf{x}_2) \leq (1 - \lambda)f(\mathbf{x}_1) + \lambda f(\mathbf{x}_2) = f(\mathbf{x}_1) + \lambda(f(\mathbf{x}_2) - f(\mathbf{x}_1))$$

and furthermore all points $(1 - \lambda)\mathbf{x}_1 + \lambda\mathbf{x}_2$ are feasible since X is a convex set and \mathbf{x}_1 and \mathbf{x}_2 are feasible. Since $f(\mathbf{x}_2) < f(\mathbf{x}_1)$, this means that

$$f((1 - \lambda)\mathbf{x}_1 + \lambda\mathbf{x}_2) < f(\mathbf{x}_1)$$

even as $\lambda \rightarrow 0$, contradicting the assumption that \mathbf{x}_1 is a local minimum. \square

Proposition 4.5. *If f is a convex function and X is a convex set, then the set of global minima is convex.*

Proof. Let \hat{X} be the set of global minima of f over the feasible region X . Choose any two global optima $\hat{\mathbf{x}}_1 \in \hat{X}$ and $\hat{\mathbf{x}}_2 \in \hat{X}$, and any $\lambda \in [0, 1]$. Since $\hat{\mathbf{x}}_1$ and $\hat{\mathbf{x}}_2$ are global minima, $f(\hat{\mathbf{x}}_1) = f(\hat{\mathbf{x}}_2)$; let \hat{f} denote this common value. Because X is a convex set, the point $(1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2$ is also feasible. Because f is a convex function,

$$\begin{aligned} f((1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2) &\leq (1 - \lambda)f(\hat{\mathbf{x}}_1) + \lambda f(\hat{\mathbf{x}}_2) \\ &= (1 - \lambda)\hat{f} + \lambda\hat{f} \\ &= \hat{f} \end{aligned}$$

Therefore $f((1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2) \leq \hat{f}$. But at the same time, $f((1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2) \geq \hat{f}$ because \hat{f} is the global minimum value of f . So we must have $f((1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2) = \hat{f}$ which means that this point is also a global minimum and $(1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2 \in \hat{X}$ as well, proving its convexity. \square

Proposition 4.6. (*Uniqueness.*) *If f is a strictly convex function and X is a convex set, the set \hat{X} contains at most one element.*

Proof. By contradiction, assume that \hat{X} contains two distinct elements $\hat{\mathbf{x}}_1$ and $\hat{\mathbf{x}}_2$. Repeating the proof of the previous proposition, because f is strictly convex, the first inequality becomes strict and we must have $f((1 - \lambda)\hat{\mathbf{x}}_1 + \lambda\hat{\mathbf{x}}_2) < \hat{f}$. This contradicts the assumption that $\hat{\mathbf{x}}_1$ and $\hat{\mathbf{x}}_2$ are global minima. \square

4.4 Basic Optimization Techniques

This section discusses techniques which are common in transportation network analysis. We begin with simple problems: optimization problems with a single decision variable, and optimization problems with no constraints, before moving to more general problem classes. In this section, we'll assume that both the objective function and feasible region are convex and differentiable, which greatly simplifies matters.

4.4.1 One-dimensional problems

To start off, consider a simple minimization problem in one variable with no constraints:

$$\min_x f(x)$$

Because f is convex, any local minimum is also a global minimum. So, all we need is to know when we've reached a local minimum. Thus, we're looking for a set of *optimality conditions*, that is, a set of equations and/or inequalities in x which are true if and only if x is optimal. For the unconstrained case, this is easy: we know from basic calculus that \hat{x} is a local minimum if

$$f'(\hat{x}) = 0$$

(We don't have to check whether \hat{x} is a local minimum or a local maximum because f is convex.)

Example 4.6. Find the value of x which minimizes $f(x) = x^2 - 3x + 5$.

Solution. $f'(x) = 2x - 3$, which vanishes if $x = 3/2$. Therefore $x = 3/2$ minimizes $f(x)$. ■

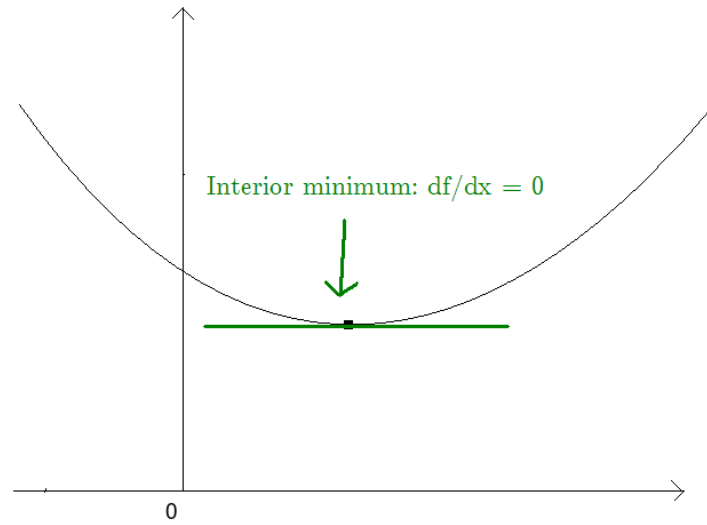
It's a little bit more complicated if we add constraints to the picture. For instance, consider the function in Figure 4.4 (which could very well be the function from Example 4.6), but with the added constraint $x \geq 0$. In this case, nothing is different, and the optimum still occurs where $f'(x)$ vanishes, that is, at $x = 3/2$. But what if the constraint was $x \geq 2$? In this case, $x = 3/2$ is infeasible, and $f'(x)$ is always strictly positive in the entire feasible region. This means that f is strictly increasing over the entire feasible region, so the minimum value is obtained at the smallest possible value of x , that is, $x = 2$. So we see that sometimes the local minimum of a constrained optimization problem can be at a point where $f'(x)$ is nonzero.

To simplify things a little bit, assume that the constraint is of the form $x \geq 0$, that is, we are trying to solve

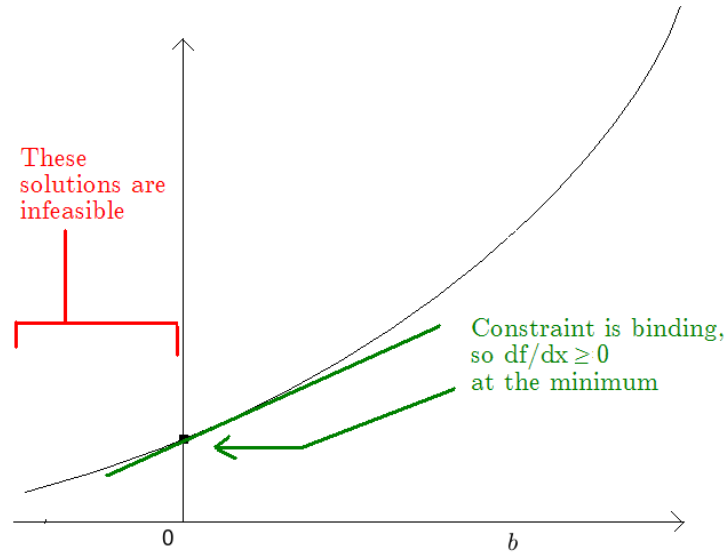
$$\begin{array}{ll} \min_x & f(x) \\ \text{s.t.} & x \geq 0 \end{array}$$

As Figure 4.5 shows, there are only two possibilities. In the first case, the minimum occurs when x is strictly positive. We can call this an *interior* minimum, or we can say that the constraint $x \geq 0$ is *nonbinding* at this point. In this case, clearly $f'(x)$ must equal zero: otherwise, we could move slightly in one direction or the other, and reduce f further. The other alternative is that the minimum occurs for $x = 0$, as in Figure 4.5(b). For this to be a minimum, we need $f'(0) \geq 0$ — if $f'(0) < 0$, f is decreasing at $x = 0$, so we could move to a slightly positive x , and thereby reduce f .

Let's try to draw some general conclusions. For the interior case of Figure 4.5(a), we needed $x \geq 0$ for feasibility, and $f'(x) = 0$ for optimality. For the



(a) Convex function where the constraint is not binding at the minimum.



(b) Convex function where the constraint is binding at the minimum.

Figure 4.5: Two possibilities for minimizing a convex function with a constraint.

boundary case of Figure 4.5(b), we had $x = 0$ exactly, and $f'(x) \geq 0$. So we see that in both cases, $x \geq 0$ and $f'(x) \geq 0$, and furthermore that at least one of these has to be exactly equal to zero. To express the fact that either x or $f'(x)$ must be zero, we can write $xf'(x) = 0$. So a solution \hat{x} solves the minimization problem if and only if

$$\begin{aligned}\hat{x} &\geq 0 \\ f'(\hat{x}) &\geq 0 \\ \hat{x}f'(\hat{x}) &= 0.\end{aligned}$$

Whenever we can find \hat{x} that satisfies these three conditions, we know it is optimal. These are often called *first-order conditions* because they are related to the first derivative of f . The condition $\hat{x}f'(\hat{x}) = 0$ is an example of a *complementarity constraint* because it forces either \hat{x} or $f'(\hat{x})$ to be zero.

4.4.2 Bisection method

The bisection method allows us to solve one-dimensional problems over bounded feasible regions. Consider the one-dimensional optimization problem

$$\begin{aligned}\min_x \quad & f(x) \\ \text{s.t.} \quad & a \leq x \leq b\end{aligned}$$

where f is continuously differentiable. The bisection method works by constantly narrowing down the region where the optimal solution lies. After the k -th iteration, the bisection method will tell you that the optimum solution lies in the interval $[a_k, b_k]$, with this interval shrinking over time (that is, $b_k - a_k < b_{k-1} - a_{k-1}$). A natural termination criterion is to stop when the interval is sufficiently small, that is, when $b_k - a_k < \epsilon$, where ϵ is the precision you want for the final solution.

Here's how the algorithm works.

Step 0: Initialize. Set the iteration counter $k = 0$, $a_0 = a$, $b_0 = b$.

Step 1: Evaluate midpoint. Calculate the derivative of f at the middle of this interval, $d_k = f'((a_k + b_k)/2)$

Step 2: Bisect. If $d_k > 0$, set $a_{k+1} = a_k$, $b_{k+1} = (a_k + b_k)/2$. Otherwise, set $a_{k+1} = (a_k + b_k)/2$, $b_{k+1} = b_k$.

Step 3: Iterate. Increase the counter k by 1 and check the termination criterion. If $b_k - a_k < \epsilon$, then terminate; otherwise, return to step 1.

Example 4.7. Use the bisection algorithm to find the minimum of $f(x) = (x - 1)^2$ on the interval $x \in [0, 3]$, to within a tolerance of 0.1.

Solution. You may find it useful to follow along in Table 4.1. We start off with $k = 0$, $a_0 = 0$, and $b_0 = 3$. We calculate the derivative at the midpoint:

Table 4.1: Demonstration of the bisection algorithm with $f(x) = (x - 1)^2$, $x \in [0, 3]$.

k	a_k	b_k	$(a_k + b_k)/2$	d_k
0	0	3	1.5	$1 > 0$
1	0	1.5	0.75	$-0.5 < 0$
2	0.75	1.5	1.125	$0.25 > 0$
3	0.75	1.125	0.938	$-0.125 < 0$
4	0.938	1.125	1.031	$0.062 > 0$
5	0.938	1.031	0.984	$-0.032 < 0$

$f'(x) = 2(x - 1)$, so $f'(1.5) = 1$, which is positive. Since f' is positive at $x = 1.5$, the minimum must occur to the *left* of this point, that is, somewhere in the interval $[0, 1.5]$. We set a_1 and b_1 equal to these new values, and repeat. The new midpoint is 0.75, and $f'(0.75) = -0.5$ is negative, so the minimum must occur to the *right* of this point. Thus, the new interval is $[0.75, 1.5]$, and we repeat as shown in Table 4.1. At the end of the fifth iteration, the width of the interval is $1.031 - 0.938 = 0.093$, which is less than the tolerance of 0.1; therefore we stop, and return our best guess of the optimum as the midpoint of this interval: $\hat{x} \approx 0.984$. The true minimum point occurs at $\hat{x} = 1$; if we had chosen a smaller tolerance ϵ , the algorithm would have narrowed the interval further, with both ends converging towards 1.

4.4.3 Multidimensional problems with nonnegativity constraints

Most interesting optimization problems have more than one decision variable. However, we will keep the assumption that the only constraint on the decision variables is nonnegativity. Using the vector \mathbf{x} to refer to all of the decision variables, we solve the problem

$$\begin{array}{ll} \min_{\mathbf{x}} & f(\mathbf{x}) \\ \text{s.t.} & \mathbf{x} \geq \mathbf{0} \end{array}$$

Using the same logic as before, we can show that $\hat{\mathbf{x}}$ solves this problem if and only if the following conditions are satisfied for *every* decision variable x_i :

$$\begin{aligned} \hat{x}_i &\geq 0 \\ \frac{\partial f(\hat{\mathbf{x}})}{\partial x_i} &\geq 0 \\ \hat{x}_i \frac{\partial f(\hat{\mathbf{x}})}{\partial x_i} &= 0 \end{aligned}$$

You should convince yourself that if these conditions are not met for each decision variable, then $\hat{\mathbf{x}}$ cannot be optimal: if the first condition is violated, the

solution is infeasible; if the second is violated, the objective function can be reduced by increasing x_i ; if the first two are satisfied but the third is violated, then $\hat{x}_i > 0$ and $\frac{\partial f(\hat{x})}{\partial x_i} > 0$, and the objective function can be reduced by decreasing x_i .

This can be compactly written in vector form as

$$\mathbf{0} \leq \hat{\mathbf{x}} \perp \nabla f(\hat{\mathbf{x}}) \geq \mathbf{0} \quad (4.10)$$

where the \perp symbol indicates orthogonality, i.e. that the dot product of $\hat{\mathbf{x}}$ and $\nabla f(\hat{\mathbf{x}})$ is zero.

Unfortunately, the bisection algorithm does not work nearly as well in higher dimensions. It is difficult to formulate an extension that always works, and those that do are inefficient. We'll approach solution methods for higher-dimensional problems somewhat indirectly, tackling a few other topics first: addressing constraints other than nonnegativity, and a few highlights of linear optimization.

4.4.4 Constrained problems

Constraints can take a very large number of forms. For the purposes of this book, we can restrict attention to only two types of constraints: *linear equality* constraints and *nonnegativity* constraints. The previous two sections showed you how to deal with nonnegativity constraints; this section discusses linear equality constraints. A linear equality constraint is of the form

$$\sum_i a_i x_i = b, \quad (4.11)$$

where the x_i are decision variables, and the a_i and b are constants.

We can handle these using the technique of Lagrange multipliers. This technique is demonstrated in the following example for the case of a single linear equality constraint.

$$\begin{array}{ll} \min_{x_1, x_2} & x_1^2 + x_2^2 \\ \text{s.t.} & x_1 + x_2 = 5 \end{array}$$

(It is a useful exercise to verify that $x_1^2 + x_2^2$ is a strictly convex function, and that $\{x_1, x_2 : x_1 + x_2 = 5\}$ is a convex set.)

The main idea behind Lagrange multipliers is that unconstrained problems are easier than constrained problems. The technique is an ingenious way of nominally removing a constraint while still ensuring that it holds at optimality. The equality constraint is “brought into the objective function” by multiplying the difference between the right- and left-hand sides by a new decision variable κ (called the Lagrange multiplier), adding the original objective function. This creates the Lagrangian function

$$\mathcal{L}(x_1, x_2, \kappa) = x_1^2 + x_2^2 + \kappa(5 - x_1 - x_2) \quad (4.12)$$

It is possible to show that *the optimal solutions of the original optimization problem correspond to stationary points of the Lagrangian function*, that is, to

values of x_1 , x_2 , and κ such that $\nabla \mathcal{L}(x_1, x_2, \kappa) = \mathbf{0}$. To find this stationary point, take partial derivatives with respect to each variable and set them all equal to zero:

$$\frac{\partial \mathcal{L}}{\partial x_1} = 2x_1 - \kappa = 0 \quad (4.13)$$

$$\frac{\partial \mathcal{L}}{\partial x_2} = 2x_2 - \kappa = 0 \quad (4.14)$$

$$\frac{\partial \mathcal{L}}{\partial \kappa} = 5 - x_1 - x_2 = 0 \quad (4.15)$$

Notice that the third optimality condition (4.15) *is simply the original constraint*, so this stationary point must be feasible. Equations (4.13) and (4.14) respectively tell us that $x_1 = \kappa/2$ and $x_2 = \kappa/2$; substituting these expressions into (4.15) gives $\kappa = 5$, and therefore the optimal solution occurs for $x_1 = x_2 = 5/2$.

This technique generalizes perfectly well to the case of multiple linear equality constraints. Consider the general optimization problem

$$\begin{array}{ll} \min_{x_1, \dots, x_n} & f(x_1, \dots, x_n) \\ \text{s.t.} & \sum_{i=1}^n a_{1i}x_i = b_1 \\ & \sum_{i=1}^n a_{2i}x_i = b_2 \\ & \vdots \\ & \sum_{i=1}^n a_{mi}x_i = b_m \end{array}$$

where f is convex. The corresponding Lagrangian is

$$\begin{aligned} \mathcal{L}(x_1, \dots, x_n, \kappa_1, \dots, \kappa_m) = & f(x_1, \dots, x_n) + \kappa_1 \left(b_1 - \sum_{j=1}^n a_{1j}x_j \right) + \\ & \kappa_2 \left(b_2 - \sum_{j=1}^n a_{2j}x_j \right) + \dots + \kappa_m \left(b_m - \sum_{j=1}^n a_{mj}x_j \right) \end{aligned}$$

For an optimization problem that has *both* linear equality constraints and nonnegativity constraints, we form the optimality conditions by combining the Lagrange multiplier technique with the complementarity technique from the previous section. Thinking back to Section 4.4.3, in the same way that we replaced the condition $f'(\hat{x}) = 0$ for the unconstrained case with the three conditions $\hat{x} \geq 0$, $f'(\hat{x}) \geq 0$, and $\hat{x}f'(\hat{x}) = 0$ when the nonnegativity constraint was added, we'll adapt the Lagrangian optimality conditions. If the optimization

problem has the form

$$\begin{aligned} \min_{x_1, \dots, x_n} \quad & f(x_1, \dots, x_n) \\ \text{s.t.} \quad & \sum_{i=1}^n a_{1i}x_i = b_1 \\ & \sum_{i=1}^n a_{2i}x_i = b_2 \\ & \vdots \\ & \sum_{i=1}^n a_{mi}x_i = b_m \\ & x_1, \dots, x_n \geq 0 \end{aligned}$$

where $m \leq n$, then the Lagrangian is

$$\begin{aligned} \mathcal{L}(x_1, \dots, x_n, \kappa_1, \dots, \kappa_m) = & f(x_1, \dots, x_n) + \kappa_1 \left(b_1 - \sum_{j=1}^n a_{1j}x_j \right) + \\ & \kappa_2 \left(b_2 - \sum_{j=1}^n a_{2j}x_j \right) + \dots + \kappa_m \left(b_m - \sum_{j=1}^n a_{mj}x_j \right) \end{aligned}$$

and the optimality conditions are

$$\begin{aligned} \frac{\partial \mathcal{L}}{\partial x_i} &\geq 0 & \forall i \in \{1, \dots, n\} \\ \frac{\partial \mathcal{L}}{\partial \kappa_j} &= 0 & \forall j \in \{1, \dots, m\} \\ x_i &\geq 0 & \forall i \in \{1, \dots, n\} \\ x_i \frac{\partial \mathcal{L}}{\partial x_i} &= 0 & \forall i \in \{1, \dots, n\} \end{aligned}$$

Be sure you understand what each of these formulas implies. Each decision variable must be nonnegative, the partial derivative of \mathcal{L} with respect to this variable must be nonnegative, and their product must equal zero (for the same reasons as discussed in Section 4.4.3). For the Lagrange multipliers $(\kappa_1, \dots, \kappa_m)$, the corresponding partial derivative of \mathcal{L} must be zero. Notice how this is a combination of the two techniques.

For small optimization problems, we can write down each of these conditions and solve for the optimal solution, as above. However, for large-scale problems this process can be very inconvenient. Later chapters in the book explain methods which work better for large problems. As a final note, the full theory of Lagrange multipliers is more involved than what is discussed here. However, it suffices for the case of a convex objective function and linear equality constraints.

4.5 Metaheuristics (*)

(This section is optional and can be skipped. However if you are interested in techniques that can be applied to optimization problems that are not convex, read this section to learn about simulated annealing and genetic algorithms.)

In optimization, there is often a tradeoff between how widely applicable a solution method is, and how efficient or effective it is at solving specific problems. For any one specific problem, a tailor-made solution process is likely much faster than a generally-applicable method, but generating such a method requires more effort and specialized knowledge, and is less “rewarding” in the sense that the method can only be more narrowly applied. Some of the most general techniques are *heuristics*, which are not guaranteed to find the global optimum solution, but tend to work reasonably well in practice. In practice, they only tend to be applied for very large or very complicated problems which cannot be solved exactly in a reasonable amount of time with our current knowledge.

Many engineers are initially uncomfortable with the idea of a heuristic. After all, the goal of an optimization problem is to find an optimal solution, so why should we settle for something which is only approximately “optimal,” often without any guarantees of how approximate the solution is? First, for very complicated problems a good heuristic can often return a reasonably good solution in much less time than it would take to find the exact, global optimal solution. For many practical problems, the cost improvement from a reasonably good solution to an exactly optimal one is not worth the extra expense (both time and computational hardware) needed, particularly if the heuristic gets you within the margin of error based on the input data.

Heuristics are also very, very common in psychology and nature. If I give someone a map and ask them to find the shortest-distance route between two points in a city by hand, they will almost certainly not formulate an mathematical model and solve it to provable optimality. Instead, they use mental heuristics (rules of thumb based on experience) and can find paths which are actually quite good. Many of the heuristics are inspired by things seen in nature.

An example is how ant colonies find food. When a lone wandering ant encounters a food source, it returns to the colony and lays down a chemical pheromone. Other ants who stumble across this pheromone begin to follow it to the food source, and lay down more pheromones, and so forth. Over time, more and more ants will travel to the food source, taking it back to the colony, until it is exhausted at which point the pheromones will evaporate. Is this method the optimal way to gather food? Perhaps not, but it performs well enough for ants to have survived for millions of years!

Another example is the process of evolution through natural selection. The human body, and many other organisms, function remarkably well in their habitats, even if their biology is not exactly “optimal.”⁴ One of the most common heuristics in use today, and one described below, is based on applying principles of natural selection and mutation to a “population” of candidate solutions to an optimization problem, using an evolutionary process to identify better and better solutions over time. This volume describes two heuristics: *simulated annealing*, and *genetic algorithms*, both of which can be applied to many different optimization problems.

⁴For instance, in humans the retina is “backwards,” creating a blind spot; in giraffes, the laryngeal nerve takes an exceptionally long and roundabout path; and the descent of the testes makes men more vulnerable to hernias later in life.

4.5.1 Simulated annealing

Simulated annealing is a simple heuristic that makes an analogy to metallurgy, but this analogy is best understood after the heuristic is described.

As a better starting point for understanding simulated annealing, consider the following “local search” heuristic. (For now, descriptions will be a bit vague; more precise definitions will follow soon.) Local search proceeds as follows:

1. Choose some initial feasible solution $\mathbf{x} \in X$, and calculate the value of the objective function $f(\mathbf{x})$.
2. Generate a new feasible solution $\mathbf{x}' \in X$ which is close to the current solution \mathbf{x} .
3. Calculate $f(\mathbf{x}')$
4. If $f(\mathbf{x}') \leq f(\mathbf{x})$, the new solution is better than the old solution. So update the current solution by setting \mathbf{x} equal to \mathbf{x}' .
5. Otherwise, the old solution is better, so leave \mathbf{x} unchanged.
6. Return to step 2 and repeat until we are unable to make further progress.

One way to visualize local search is a hiker trying to find the lowest elevation point in a mountain range. In this analogy, the park boundaries are the feasible set, the elevation of any point is the objective function, and the location of the hiker is the decision variable. In local search, starting from his or her initial position, the hiker looks around and finds a nearby point which is lower than their current point. If such a point can be found, they move in that direction and repeat the process. If every neighboring point is higher, then they stop and conclude that they have found the lowest point in the park.

It is not hard to see why this strategy can fail; if there are multiple local optima, the hiker can easily get stuck in a point which is not the lowest (Figure 4.6). Simulated annealing attempts to overcome this deficiency of local search by allowing a provision for occasionally moving in an uphill direction, in hopes of finding an even lower point later on. Of course, we don’t always want to move in an uphill direction and have a preference for downhill directions, but this preference cannot be absolute if there is any hope of escaping local minima.

Simulated annealing accomplishes this by making the decision to move or not probabilistic, introducing a *temperature* parameter T which controls how likely you are to accept an “uphill” move. When the temperature is high, the probability of moving uphill is large, but when the temperature is low, the probability of moving uphill becomes small. In simulated annealing, the temperature is controlled with a *cooling schedule*. Initially, the temperature is kept high to encourage a broad exploration of the feasible region. As the temperature decreases slowly, the solution is drawn more and more to lower areas. The cooling schedule can be defined by an initial temperature T_0 , a final temperature T_f , the number of search iterations n at a given temperature level,

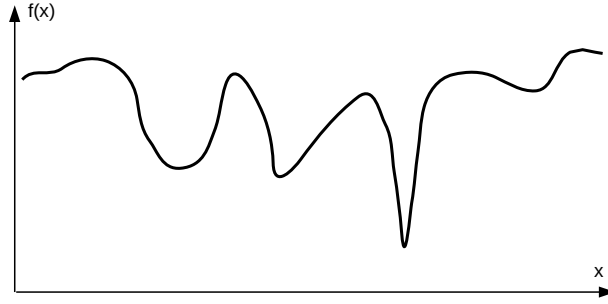


Figure 4.6: Moving downhill from the current location may not lead to the global optimum.

and a scaling factor $k \in (0, 1)$ which is applied every n iterations to reduce the temperature.

Finally, because the search moves both uphill and downhill, there is no guarantee that the final point of the search is the best point found so far. So, it is worthwhile to keep track of the best solution \mathbf{x}^* encountered during the algorithm. (This is analogous to the hiker keeping a record of the lowest point observed, and returning to that point when done searching.)

So, the simulated annealing algorithm can be stated as follows:

1. Choose some initial feasible solution $\mathbf{x} \in X$, and calculate the value of the objective function $f(\mathbf{x})$.
2. Initialize the best solution to the initial one $\mathbf{x}^* \leftarrow \mathbf{x}$.
3. Set the temperature to the initial temperature: $T \leftarrow T_0$
4. Repeat the following steps n times:
 - (a) Generate a new feasible solution \mathbf{x}' which neighbors \mathbf{x} .
 - (b) If $f(\mathbf{x}') < f(\mathbf{x}^*)$, it is the best solution found so far, so update $\mathbf{x}^* \leftarrow \mathbf{x}'$.
 - (c) If $f(\mathbf{x}') \leq f(\mathbf{x})$, it is a better solution than the current one, so update $\mathbf{x} \leftarrow \mathbf{x}'$.
 - (d) Otherwise, update $\mathbf{x} \leftarrow \mathbf{x}'$ with probability $\exp(-[f(\mathbf{x}') - f(\mathbf{x})]/T)$
5. If $T > T_f$, then reduce the temperature ($T \leftarrow kT$) and return to step 4.
6. Report the best solution found \mathbf{x}^*

The key step is step 4d. Notice how the probability of “moving uphill” depends on two factors: the temperature, and how much the objective function will increase. The algorithm is more likely to accept an uphill move if it is only slightly uphill, or if the temperature is high. The exponential function captures

these effects while keeping the probability between 0 and 1. A few points require explanation.

How should the cooling schedule be chosen? Unfortunately, it is hard to give general guidance here. Heuristics often have to be “tuned” for a particular problem: some problems do better with higher temperatures and slower cooling (k values closer to 1, n larger), others work fine with faster cooling. When you use simulated annealing, you should try different variations of the cooling schedule to identify one that works well for your specific problem.

How should an initial solution be chosen? It is often helpful if the initial solution is relatively close to the optimal solution. For instance, if the optimization problem concerns business operations, the current operational plan can be used as the initial solution for further optimization. However, it’s easy to go overboard with this. You don’t want to spend so long coming up with a good initial solution that it would have been faster to simply run simulated annealing for longer starting from a worse solution. The ideal is to think of a good, quick rule of thumb for generating a reasonable initial solution; failing that, you can always choose the initial solution randomly.

How do I define a neighboring solution for step 4a? Again, this is problem-specific, and one of the decisions that must be made when applying simulated annealing. A good neighborhood definition should involve points which are “close” to the current solution in some way, but ensure that feasible solutions are connected in the sense that any two feasible solutions can be reached by a chain of neighboring solutions. For the examples in Section 4.2, some possibilities (emphasis on *possibilities*, there are other choices) are:

Transit frequency setting problem: The decision variables \mathbf{n} are the number of buses assigned to each route. Given a current solution \mathbf{n} , a neighboring solution is one where exactly one bus has been assigned to a different route.

Scheduling maintenance: The decision variables are \mathbf{x} , indicating where and when maintenance is performed, and \mathbf{c} , indicating the condition of the facilities. In this problem, the constraints completely define \mathbf{c} in terms of \mathbf{x} , so for simulated annealing it is enough to simply choose a feasible solution \mathbf{x} , then calculate \mathbf{c} using the state evolution equations. In this problem, there is no reason to perform less maintenance than the budget allows each year. If the initial solution is chosen in this way, then a neighboring solution might be one where one maintenance activity on a facility is reassigned to another facility that same year. If the resulting solution is infeasible (because the resulting \mathbf{c} values fall below the minimum threshold), then another solution should be generated.

Facility location problem: The decision variables are the intersections that the three terminals are located at, L_1 , L_2 , and L_3 . Given current values for these, in a neighboring solution two of the three terminals are at the same location, but one of the three has been assigned to a different location.

Shortest path problem: The decision variables specify a path between the origin and destination. Given a current path, a neighboring path is one which differs in only intersection. (Can you think of a way to express this mathematically?)

How do I perform a step with a given probability? Most programming languages have a way to generate a uniform random variable between 0 and 1. If we want to perform a step with probability p , generate a random number from the continuous uniform $(0, 1)$ distribution. If this number is less than p , perform the step; otherwise, do not.

Example with facility location

This section demonstrates simulated annealing, on an instance of the facility location problem from Section 4.2. In this problem instance, the road grid consists of ten north-south and ten east-west streets. The cost of locating a terminal at each of the locations is shown in Figure 4.7 (these numbers were randomly generated). There are 30 customers, randomly located throughout the grid, as shown in Figure 4.8. In these figures, the coordinate system (x, y) is used to represent points, where x represents the number of blocks to the right of the upper-left corner, and y represents the number of blocks below the upper-left corner.

The cooling schedule is as follows: the initial temperature is $T_0 = 1000$, the final temperature is $T_f = 100$, $k = 0.75$, and the temperature is reduced every 8 iterations. (These values were chosen through trial-and-error, and work reasonably well for this problem.) The first several iterations are shown below.

The algorithm begins with an initial solution, generated randomly. Suppose this solution locates the three facilities at coordinates $(5, 1)$, $(8, 6)$, and $(6, 1)$, respectively; the total cost associated with this solution is 158.0 cost units. A neighboring solution is generated by picking one of the facilities (randomly) and assigning it to another location (randomly). Suppose that the third facility is reassigned to location $(1, 0)$, so the new candidate solution locates the facilities at $(5, 1)$, $(8, 6)$, and $(1, 0)$. This solution has a cost of 136.4 units, which is lower than the current solution. So, simulated annealing replaces the current solution with the new one.

In the next iteration, suppose that the first facility is randomly chosen to be reassigned to $(8, 2)$, so the candidate solution is $(8, 2)$, $(8, 6)$, and $(1, 0)$. This solution has a cost of 149.1, which is higher than the current cost of 136.4. The algorithm will not always move to such a solution, but will only do so with probability calculated as in Step 4c:

$$p = \exp(-[f(\mathbf{x}') - f(\mathbf{x})]/T) = \exp(-[149.1 - 136.4]/100) = 0.881$$

This probability is high, because (being one of the early iterations) the temperature is set high. If the temperature were lower, the probability of accepting this move would be lower as well.

8.7	13.5	11.0	10.2	6.3	6.4	12.4	13.5	9.7	6.8
5.8	12.4	8.8	15.1	14.7	13.9	9.0	5.2	8.7	12.7
11.8	9.8	12.4	12.6	6.3	13.3	7.7	10.2	13.6	13.5
5.4	6.9	11.1	8.4	14.2	10.9	10.7	11.7	8.1	8.9
4.9	12.5	10.8	6.6	12.0	6.8	11.9	9.2	9.5	10.0
14.4	9.7	8.6	11.6	8.0	6.7	12.7	5.9	7.6	14.1
7.7	9.6	6.4	9.9	9.2	13.3	12.3	14.7	15.0	5.1
7.6	14.7	7.1	5.5	5.5	9.7	9.7	14.4	7.9	15.1
10.9	12.5	9.2	12.1	14.8	6.4	6.1	10.5	12.5	10.8
12.0	5.9	13.1	10.0	11.9	10.0	8.6	8.4	10.7	5.3

Figure 4.7: Cost of locating facilities at each intersection.

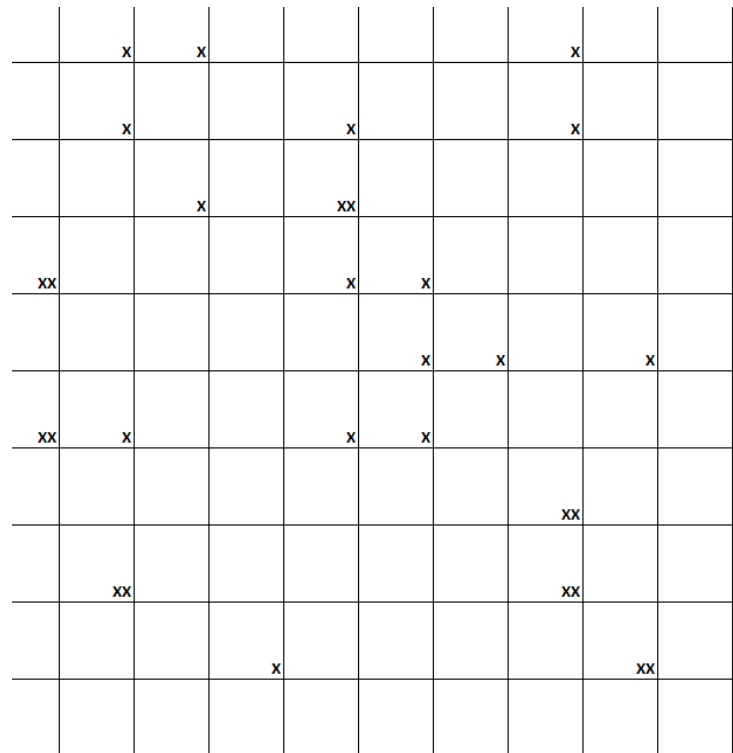


Figure 4.8: Locations of customers.

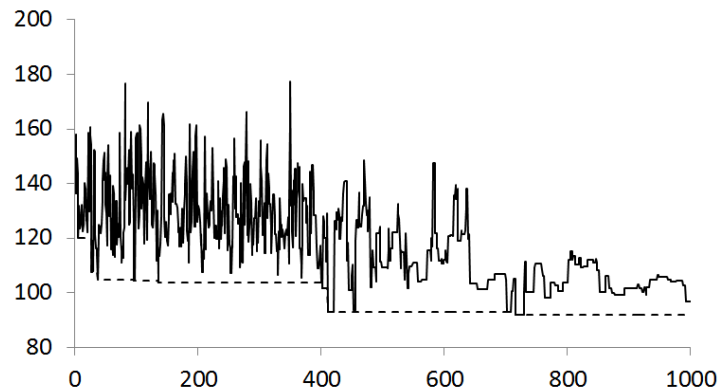


Figure 4.9: Progress of simulated annealing. Solid line shows current solution cost, dashed line best cost so far.

Supposing that the move is accepted, the algorithm replaces the current solution with the candidate and continues as before. If, on the other hand, the move is rejected, the algorithm generates another candidate solution based on the same current solution (5,1), (8,6), and (1,0). The algorithm continues in the same way, reducing the temperature by 25% every 8 iterations.

The progress of the algorithm until termination is shown in Figure 4.9. The solid line shows the cost of the current solution, while the dashed line tracks the cost of the best solution found so far. A few observations are worth making. First, in the early iterations the cost is highly variable, but towards termination the cost becomes more stable. This is due to the reduction in temperature which occurs over successive iterations. When the temperature is high, nearly any move will be accepted so one expects large fluctuations in the cost. When the temperature is low, the algorithm is less likely to accept cost-increasing moves, so fewer fluctuations are seen. Also notice that the best solution was found shortly after iteration 700. The final solution is not the best, although it is close.

The best facility locations found by simulated annealing are shown in Figure 4.10, with a total cost of 92.1.

4.5.2 Genetic algorithms

Another heuristic is the genetic algorithm. In contrast to simulated annealing, where there is a single search happening through the feasible region (recall the analogy with the hiker), genetic algorithms maintain a larger *population* of solutions which reflect a diverse set of points within the feasible region. Genetic algorithms work by attempting to improve the population from one iteration to the next.

The process of improvement is intended to mimic the processes of natural

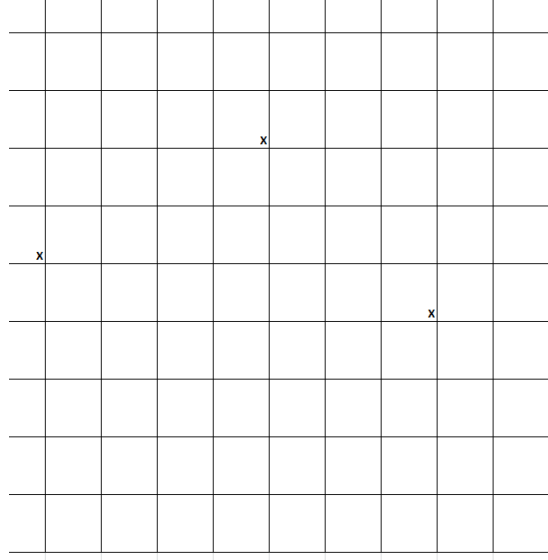


Figure 4.10: Locations of facilities found by simulated annealing (cost 92.1).

selection, reproduction, and mutation which are observed in biology. For the sake of our purposes, *natural selection* means identifying solutions in the population which have good (low) values of the objective function. The hope is that there are certain aspects or patterns in these solutions which make them good, which can be maintained in future generations. Given an initial population as the first generation, subsequent generations are created by choosing good solutions from the previous generation, and “breeding” them with each other in a process called *crossover* which mimics sexual reproduction: new “child” solutions are created by mixing characteristics from two “parents.” Lastly, there is a random *mutation* element, where solutions are changed externally with a small probability. If all goes well, after multiple generations the population will tend towards better and better solutions to the optimization problem.

At a high level, the algorithm is thus very straightforward and presented below. However, at a lower level, there are a number of careful choices which need to be made about how to implement each step for a particular problem. Thus, each of the steps is described in more detail below. The high-level version of genetic algorithms is as follows:

1. Generate an initial population of N feasible solutions (generation 0), set generation counter $g \leftarrow 0$.
2. Create generation $g + 1$ in the following way, repeating each step N times:
 - (a) Choose two “parent” solutions from generation g .
 - (b) Combine the two parent solutions to create a new solution.

(c) With probability p , mutate the new solution.

3. Increase g by 1 and return to step 2 unless done.

Although not listed explicitly, it is a good idea to keep track at every stage of the best solution \mathbf{x}^* found so far. Just like in simulated annealing, there is no guarantee that the best solution will be found in the last generation. So, whenever a new solution is formed, calculate its objective function value, and record the solution if it is better than any found so far.

Here are explanations of each step in more detail. It is very important to realize that *there are many ways to do each of these steps* and that what is presented below is one specific example intended to give the general flavor while sparing unnecessary complications at this stage.

Generate an initial population of feasible solutions: In contrast with simulated annealing, where we had to generate an initial feasible solution, with genetic algorithms we must generate a larger population of initial feasible solutions. While it is still desirable to have these initial feasible solutions be reasonably good, it is also very important to have some diversity in the population. Genetic algorithms work by combining characteristics of different solutions together. If there is little diversity in the initial population the difference between successive generations will be small and progress will be slow.

Selection of parent solutions: Parent solutions should be chosen in a way that better solutions (lower objective function values) are more likely to be chosen. This is intended to mimic natural selection, where organisms better adapted for a particular environment are more likely to reproduce. One way to do this is through *tournament selection*, where a number of solutions from the old generation are selected randomly, and the one with the best objective function value is chosen as the first parent. Repeating the “tournament” again, randomly select another subset of solutions from the old generation, and choose the best one as the second parent. The number of entrants in the tournament is a parameter that you must choose.

Combining parent solutions: This is perhaps the trickiest part of genetic algorithms: how can we combine two feasible solutions to generate a new feasible solution which retains aspects of both parents? The exact process will differ from problem to problem. Here are some ideas, based on the example problems from Section 4.2:

Transit frequency setting problem: The decision variables are the number of buses on each route. Another way of “encoding” this decision is to make a list assigning each bus in the fleet to a corresponding route. (Clearly, given such a list, we can construct the n_r values by counting how many times a route appears.) Then, to generate a new list from two parent lists, we can assign each bus to either its route in one parent, or its route in the other parent, choosing randomly for each bus.

Scheduling maintenance: As described above, optimal solutions must always exhaust the budget each year. So, along the lines of the transit

frequency setting problem, make a list of the number of maintenance actions which will be performed in each year, noting the facility assigned to each maintenance action in the budget. To create a new solution, for each maintenance action in the budget choose one of the facilities assigned to this action from the two parents, selecting randomly. Once this list is obtained, it is straightforward to calculate \mathbf{x} and \mathbf{c} .

Facility location problem: For each facility in the child solution, make its location the same as the location of that facility in one of the two parents (chosen randomly).

Mutating solutions: Mutation is intended to provide additional diversity to the populations, avoiding stagnation and local optimal solutions. Mutation can be accomplished in the same way that neighbors are generated in simulated annealing. The probability of mutation should not be too high — as in nature, most mutations are harmful — but enough to escape from local optima. The mutation probability p can be selected by trial and error, determining what works well for your problem.

Example with facility location

This section demonstrates genetic algorithms on the same facility location problem used to demonstrate simulated annealing. The genetic algorithm is implemented with a population size of 100, over ten generations. Tournaments of size 3 are used to identify parent solutions, and the mutation probability is 0.05. The initial population is generated by locating all the facilities are completely at random.

To form the next generation, 100 new solutions need to be created, each based on combining two solutions from the previous generation. These two parents are chosen using the tournament selection rule, as shown in Figure 4.11. The two winners of the tournament locate the three facilities at (2,3), (1,7), (8,7); and (1,4), (7,6), (5,0) respectively. These combine in the following way:

1. The first facility is located either at (2,3) or (1,4); randomly choose one of them, say, (2,3).
2. The second facility is located either at (1,7) or (7,6); randomly choose one of them, say, (1,7).
3. The third facility is located either at (8,7) or (5,0); randomly choose one of them, say, (5,0).

This gives a new solution (2,3), (1,7), (5,0) in the next generation, as shown in Figure 4.11. With 5% probability, this solution will be “mutated.” If this solution is selected for mutation, one of the three facility is randomly reassigned to another location. For instance, the facility at (1,7) may be reassigned to (0,4). This process is repeated until all 100 solutions in the next generation have been created.

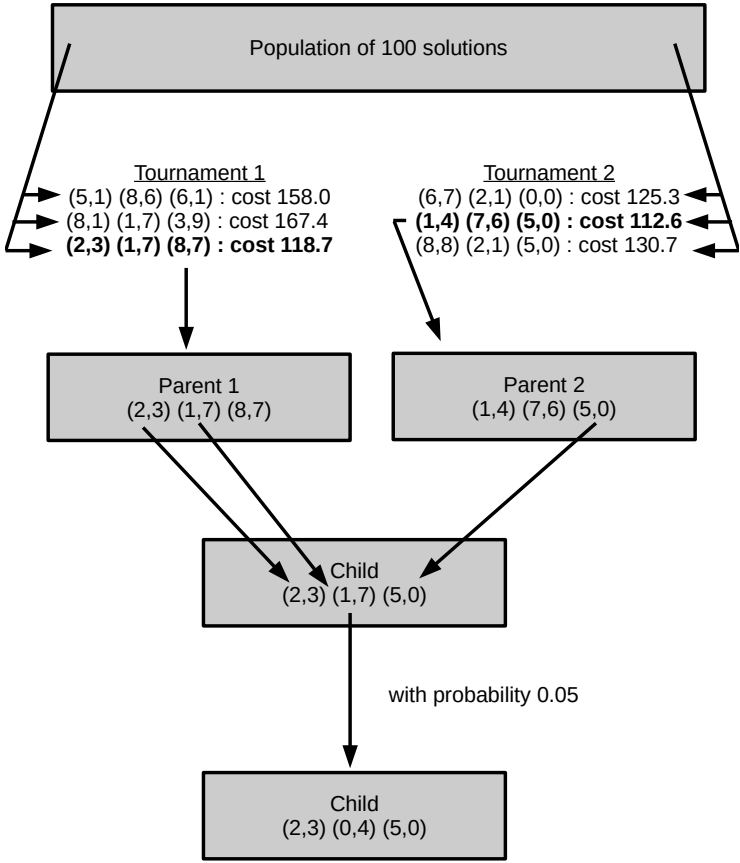


Figure 4.11: Generation of a new solution: reproduction and mutation.

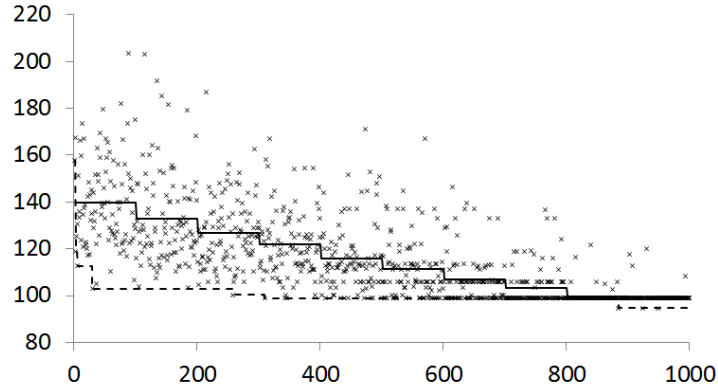


Figure 4.12: Progress of genetic algorithm. Solid line shows average generation cost, dashed line best cost so far, crosses cost of individual solutions.

Figure 4.12 shows the progress of the algorithm over ten generations. The solid line shows the average cost in each generation. The dashed line shows the cost of the best solution found so far, and the crosses show the cost of each of the solutions comprising the generations. Since lower-cost alternatives are more likely to win the tournaments and be selected as parents, the average cost of each generation decreases. Figure 4.13 shows the locations of the terminals in the best solution found.

4.6 Exercises

1. [11] Why is it generally a bad idea to use strict inequalities ($>$, $<$) in mathematical programs?
2. [12] Mathematically specify an optimization problem which has no optimal solution, and an optimization problem which has multiple optimal solutions.
3. [24] For the following functions, identify all stationary points and global optima.
 - (a) $f(x) = x^4 - 2x^3$
 - (b) $f(x_1, x_2) = 2x_1^2 + x_2^2 - x_1x_2 - 7x_2$
 - (c) $f(x_1, x_2) = x_1^6 + x_2^6 - 3x_1^2x_2^2$
 - (d) $f(x_1, x_2, x_3) = (x_1 - 4)^2 + (x_2 - 2)^2 + x_3^2 + x_1x_2 + x_1x_3 + x_2x_3$
4. [42] Express the constraint “ $x \in \mathbb{Z}$ ” (that is, x must be an integer) in the form $g(x) = 0$ for some continuous function g . (This shows that nonlinear programming can’t be any easier than integer programming!)

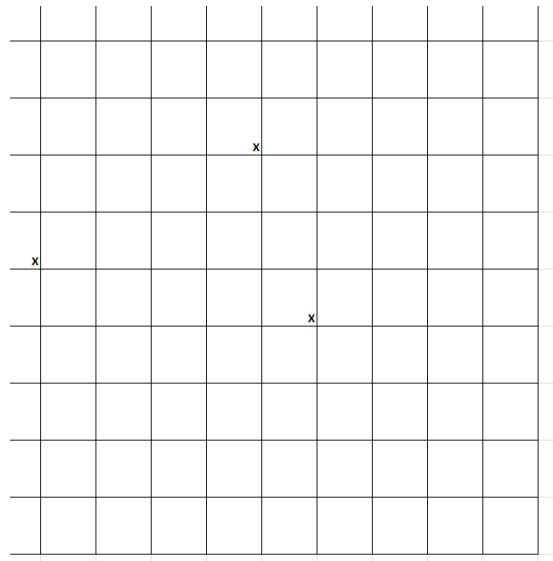


Figure 4.13: Locations of facilities found by genetic algorithm (cost 94.7).

Table 4.2: Data for Exercise 7

Bridge f	Initial condition c_f^0	Deterioration rate d_f	Repairability i_f	Maintenance cost k_f
1	85	5	3	2
2	95	10	1	4
3	75	1	1	3
4	80	6	6	1
5	100	3	3	1

5. [30] Can you ever improve the value of the objective function at the optimal solution by adding constraints to a problem? If yes, give an example. If no, explain why not.
6. [21] Prove Proposition 4.3.
7. [27] Write out the objective and all of the constraints for the maintenance scheduling problem of Example 4.3, using the bridge data shown in Table 4.2, where the maintenance cost is expressed in millions of dollars. Assume a two-year time horizon and an annual budget of \$5 million.
8. [44] Reformulate the transit frequency-setting problem of Example 4.2 so that the objective is to minimize the total amount of money spent, while achieving a pre-specified level of service (maximum delays).
9. [53] You are responsible for allocating bridge maintenance funding for a state, and must develop optimization models to assist with this for the

current year. Political realities and a strict budget will require you to develop multiple formulations for the purposes of comparison.

Let \mathcal{B} denote the set of bridges in the state. There is an “economic value” associated with the condition of each bridge: the higher the condition, the higher the value to the state, because bridges in worse condition require more routine maintenance, impose higher costs on drivers who must drive slower or put up with potholes, and carry a higher risk of unforeseen failures requiring emergency maintenance. To represent this, there is a value function $V_b(x_b)$ associated with each bridge $b \in \mathcal{B}$, giving the economic value of this bridge if x_b is spent on maintenance this year. You have a total budget of X to spend statewide. The state is also divided into n districts, and each bridge belongs to one district. Let \mathcal{D}_i denote the set of bridges which belong to district i ; different districts may contain different numbers of bridges.

In the following, be sure to define any additional notation you introduce. There is more than one possible formulation that achieves the stated goals, so feel free to explain any parts of your formulations which may not be self-evident.

- (a) **Maximize average benefits:** Formulate an optimization model (objective function, decision variables, and constraints) to maximize the average economic value of bridges across the state.
- (b) **Equitable allocation:** The previous model may recommend spending much more money in some districts than others, which is politically infeasible. Formulate an optimization model based on maximizing the total economic value to the state, but ensuring that each district receives a relatively fair share of the total budget X .
- (c) **Equitable benefits:** There is a difference between a fair allocation of money, and a fair allocation of *benefits* (say, if the functions V_b differ greatly between districts). Defining *benefit* as the difference between the economic value after investing in maintenance, and the “do-nothing” economic value, formulate an optimization model which aims to achieve a relatively fair distribution of benefits among districts.

10. [34] Consider the mathematical program

$$\begin{aligned}
 \max \quad & 3x_1 + 5x_2 + 6x_3 \\
 \text{s.t.} \quad & x_1 + x_2 + x_3 = 1 \\
 & 2x_1 + 3x_2 + 2x_3 \leq 5 \\
 & x_2 + x_3 \leq 1/2 \\
 & x_1, x_2, x_3 \geq 0
 \end{aligned}$$

- (a) Are all of the constraints needed?
- (b) Solve this problem graphically.

11. [57] In an effort to fight rising maintenance costs and congestion, a state transportation agency is considering a toll on a certain freeway segment during the peak period. Imposing a toll accomplishes two objectives at once: it raises money for the state, and also reduces congestion by making some people switch to less congested routes, travel earlier or later, carpool, take the bus, and so forth. Suppose that there are 10000 people who would want to drive on the freeway if there was no toll and no congestion, but the number of people who actually do is given by

$$x = 10000e^{(15-6\tau-t)/500}$$

where τ is the roadway toll (in dollars) and t is the travel time (in minutes). (That is, the higher the toll, or the higher the travel time, the fewer people drive.) The travel time, in turn, is given by

$$t = 15 \left[1 + 0.15 \left(\frac{x}{c} \right)^4 \right]$$

minutes, where $c = 8000$ veh/hr is the roadway capacity during rush hour. Regulations prohibit tolls exceeding \$10 at any time. Citizens are unhappy with both congestion and having to pay tolls. After conducting multiple surveys, the agency has determined that citizen satisfaction can be quantified as

$$s = 100 - t - \tau/5$$

- (a) Formulate two nonlinear programs based on this information, where the objectives are to either (a) maximize revenue or (b) maximize citizen satisfaction.
 - (b) Solve both of these problems, reporting the optimal values for the decision variables and the objective function.
 - (c) Comment on the two solutions. (e.g., do they give a similar toll value, or very different ones? How much revenue does the state give up in order to maximize satisfaction?)
 - (d) Name at least two assumptions that went into formulating this problem. Do you think they are realistic? Pick one of them, and explain how the problem might be changed to eliminate that assumption and make it more realistic.
12. [59] You are asked to design a traffic signal timing at the intersection of 8th & Grand. Assuming a simple two-phase cycle (where Grand Avenue moves in phase 1, and 8th Street in phase 2), no lost time when the signal changes, and ignoring turning movements, the total delay at the intersection can be written as

$$\frac{\lambda_1(c - g_1)^2}{2c \left(1 - \frac{\lambda_1}{\mu_1}\right)} + \frac{\lambda_2(c - g_2)^2}{2c \left(1 - \frac{\lambda_2}{\mu_2}\right)}$$

where g_1 and g_2 are the effective green time allotted to Grand Avenue and 8th Street, $c = g_1 + g_2$ is the cycle length, λ_1 and μ_1 are the arrival rate and saturation flow for Grand Avenue, and λ_2 and μ_2 are the arrival rate and saturation flow for 8th Street.

All signals downtown are given a sixty-second cycle length to foster good progression, and the arrival rate and saturation flow are 2200 veh/hr and 3600 veh/hr for Grand Avenue, respectively, and 300 veh/hr and 1900 veh/hr for 8th Street. Furthermore, no queues can remain at the end of the green interval; this means that $\mu_i g_i$ must be at least as large as $\lambda_i c$ for each approach i .

- (a) Why does the constraint $\mu_i g_i \geq \lambda_i c$ imply that no queues will remain after a green interval?
 - (b) Formulate a nonlinear program to minimize total delay.
 - (c) Simplify the nonlinear program so there is only one decision variable, and solve this nonlinear program using the bisection method of Section 4.4.2 (terminate when $b - a \leq 1$).
 - (d) Write code to automate this process, and perform a sensitivity analysis by plotting the effective green time on Grand Ave as λ_1 varies from 500 veh/hr to 3000 veh/hr, in increments of 500 veh/hr. Interpret your plot.
 - (e) Identify two assumptions in the above model. Pick one of them, and describe how you would change your model to relax that assumption.
13. [26] Find the global minima of the following functions in two ways: the bisection method, and using Newton's method to directly find a stationary point (if the Newton step leaves the feasible region, move to the boundary point closest to where Newton's method would go). Run each method for five iterations, and see which is closer to the optimum, making the comparison based on the value of the objective function at the final points.
- (a) $f(x) = -\arctan x$, $x \in [0, 10]$
 - (b) $f(x) = x \sin(1/(100x))$, $x \in [0.015, 0.04]$
 - (c) $f(x) = x^3$, $x \in [5, 15]$
14. [46] The golden section method, applicable for finding the minimum of a unimodal function f without needing derivatives.
- (a) The bisection method works by using three test points at each iteration k : the endpoints a_k and b_k , and the midpoint $c_k = (a_k + b_k)/2$ between them. If $f'(c_k)$ exists, we can use it to eliminate one half of the interval or the other. Show that if can only use the value of $f(c_k)$, and not the derivative there, it is not possible to know which half of the interval to eliminate.

- (b) The golden section method works by using *four* test points at each iteration k , dividing the interval $[a_k, b_k]$ into three subintervals: $[a_k, c_k]$, $[c_k, d_k]$, and $[d_k, b_k]$ where $a_k < c_k < d_k < b_k$. Let the widths of these three intervals be $\alpha_k = c_k - a_k$, $\beta_k = d_k - c_k$, and $\gamma_k = b_k - d_k$. Show that if we know the values $f(a_k)$, $f(b_k)$, $f(c_k)$, and $f(d_k)$, we can safely eliminate either $[a_k, c_k]$ or $[d_k, b_k]$ as not containing the minimum, and use the remainder as the interval $[a_{k+1}, b_{k+1}]$ for the next iteration.
- (c) In particular, the golden section algorithm always chooses c_k and d_k so that $\alpha = \gamma = (3 - \sqrt{5})(b_k - a_k)/2$ and $\beta = (\sqrt{5} - 2)(b_k - a_k)$. Show that, regardless of whether the upper or lower interval is eliminated, the ratio $(b_{k+1} - a_{k+1})/(b_k - a_k) = (\sqrt{5} - 1)/2$. (This is the *golden ratio*, from which the method gets its name.)
- (d) As a very useful byproduct of this choice of c_k and d_k , we only have to generate one new point (and one new objective function value) at each subsequent iteration; three of the four points needed can be reused from prior iterations. If the upper interval is eliminated in iteration k , show that $d_{k+1} = c_k$, and if the lower level is eliminated show that $c_{k+1} = d_k$.
15. [51] Give a network and a feasible solution \mathbf{x} to the mathematical formulation of the shortest path problem in Example 4.5 where the links with $x_{ij} = 1$ do not form a contiguous path between r and s , as alluded to at the end of the example.

Part II

Static Traffic Assignment

Chapter 5

Introduction to Static Assignment

5.1 Notation and Concepts

As discussed in Chapter 1, there are many possible measures of effectiveness for evaluating the impacts of a transportation project or policy. Recall that the traffic assignment problem assumes that the number of drivers traveling between each origin zone and destination zone is given, and we want to find the number of drivers using each roadway segment. This provides information on congestion, emissions, toll revenue, or other measures of interest. We are also given the underlying network, which has a set of links A representing the roadway infrastructure, and a set of nodes N representing junctions, and a set of centroids Z where trips are allowed to begin and end. Every centroid is represented by a node, so Z is a subset of N .

Let d^{rs} denote the number of drivers whose trips start at centroid r and end at centroid s . It is often convenient to write these values in matrix form, the *origin-destination matrix* (*OD matrix*, for short) \mathbf{D} where the number of trips between r and s is given by the value in the r -th row and s -th column. Together, an origin r and destination s form an *OD pair* (r, s) .

Associated with each link is its *flow*, denoted x_{ij} , representing the total number of vehicles wanting to use link (i, j) during the analysis period. The flow is also known as the *volume* or *demand*. For reasons described below, “demand” is actually the most accurate term, but “flow” and “volume” are the most common for reasons of tradition, and “flow” is used in this book. The travel time on link (i, j) is expressed as t_{ij} , and to represent congestion, we let this travel time be a function of x_{ij} and write $t_{ij}(x_{ij})$. Because of congestion effects, t_{ij} is typically increasing and convex, that is, its first two derivatives are typically positive. The function used to relate flow to travel time is called a *link performance function*. The most common used in practice is the Bureau of

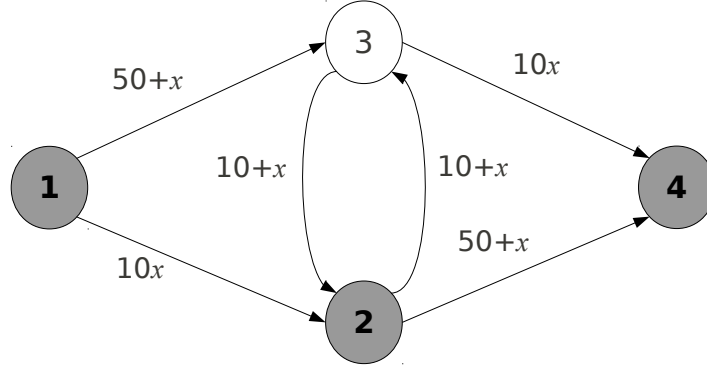


Figure 5.1: Example network for demonstration, with link performance functions shown.

Public Roads (BPR) function, named after the agency which developed it:

$$t_{ij}(x_{ij}) = t_{ij}^0 \left(1 + \alpha \left(\frac{x_{ij}}{u_{ij}} \right)^\beta \right), \quad (5.1)$$

where t_{ij}^0 is the “free-flow” travel time (the travel time with no congestion), u_{ij} is the practical capacity (typically the value of flow which results in a level of service of C or D), and α and β are shape parameters which can be calibrated to data. The values $\alpha = 0.15$ and $\beta = 4$ are commonly used if no calibration is done.

Notice that this function is well-defined for any value of x_{ij} , even when flows exceed the stated “capacity” u_{ij} . In the basic traffic assignment problem, there are no explicit upper bounds enforced on link flows. The interpretation of a “flow” greater than the capacity is actually that the *demand* for travel on the link exceeds the capacity, and queues will form. The delay induced by these queues is then reflected in the link performance function. Alternately, one can choose a link performance function which asymptotically grows to infinity as $x_{ij} \rightarrow u_{ij}$, implicitly enforcing the capacity constraint. However, this approach can introduce numerical issues in solution methods. More discussion on this issue follows in Section 5.1.1, but the short answer is that properly addressing capacity constraints in traffic assignment requires a *dynamic* traffic assignment model, which is the subject of Part III.

From the modeler’s perspective, the goal of traffic assignment is to determine the link flows in a network. But from the standpoint of the travelers themselves, it is easier to think of them each choosing a path connecting their origin to their destination. Let d^{rs} denote the number of vehicles which will be departing origin r to destination s , where r and s are both centroids. Using h^π to represent the

number of vehicles who choose path π , a *feasible assignment* is defined as a vector of path flows \mathbf{h} satisfying the following conditions:

1. $h^\pi \geq 0$ for all paths $\pi \in \Pi$. That is, path flows must be nonnegative.
2. $\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs}$ for all OD pairs (r, s) . Together with the nonnegativity condition, this requires that every vehicle traveling from r to s is assigned to exactly one of the paths connecting these zones.

Let H denote the set of feasible assignments.

Link flows x_{ij} and path flows h^π are closely linked, and we can obtain the former from the latter. Let δ_{ij}^π denote the number of times link (i, j) is used by path π , so $\delta_{ij}^\pi = 0$ if path π does not use link (i, j) , and $\delta_{ij}^\pi = 1$ if it does. With this notation, we have

$$x_{ij} = \sum_{r \in Z} \sum_{s \in Z} \sum_{\pi \in \Pi^{rs}} \delta_{ij}^\pi h^\pi. \quad (5.2)$$

This can be more compactly written using matrix notation as

$$\mathbf{x} = \mathbf{\Delta} \mathbf{h}, \quad (5.3)$$

where \mathbf{x} and \mathbf{h} are the vectors of link and path flows, respectively, and $\mathbf{\Delta}$ is the *link-path adjacency matrix*. The number of rows in this matrix is equal to the number of links, and the number of columns is equal to the number of paths in the network, and the value in the row corresponding to link (i, j) and the column corresponding to path π is δ_{ij}^π .

Given a feasible assignment, the corresponding link flows can be obtained by using equation (5.2) or (5.3). The set of *feasible link assignments* is the set X of vectors \mathbf{x} which satisfy (5.3) for some feasible assignment $\mathbf{h} \in H$.

Similarly, the path travel times c^π are directly related to the link travel times t_{ij} : the travel time of a path is simply the sum of the travel times of the links comprising that path. By the same logic, we have

$$c^\pi = \sum_{(i,j) \in A} \delta_{ij}^\pi t_{ij} \quad (5.4)$$

or, in matrix notation,

$$\mathbf{c} = \mathbf{\Delta}^T \mathbf{t}. \quad (5.5)$$

A small example illustrates these ideas. Consider the network in Figure 5.1 with four nodes and six links. The centroid nodes are shaded, and the link performance functions are as indicated. For travelers from node 2 to node 4, there are two paths: $\{(2, 3), (3, 4)\}$ and $\{(2, 4)\}$. Using the compact notation for paths, we could also write these as $[2, 3, 4]$ and $[2, 4]$. Therefore, the set of acyclic paths between these nodes is $\Pi^{24} = \{[2, 3, 4], [2, 4]\}$. You should verify for yourself that

$$\Pi^{14} = \{[1, 2, 4], [1, 2, 3, 4], [1, 3, 4], [1, 3, 2, 4]\}.$$

It will be useful to have specific indices for each path, i.e. $\Pi^{24} = \{\pi_1, \pi_2\}$ so $\pi_1 = [2, 3, 4]$ and $\pi_2 = [2, 4]$, and similarly $\Pi^{14} = \{\pi_3, \pi_4, \pi_5, \pi_6\}$ with $\pi_3 = [1, 2, 4]$, $\pi_4 = [1, 2, 3, 4]$, $\pi_5 = [1, 3, 4]$, and $\pi_6 = [1, 3, 2, 4]$.

Let's say that the demand for travel from centroid 1 to centroid 4 is 40 vehicles, and that the demand from 2 to 4 is 60 vehicles. Then $d^{14} = 40$ and $d^{24} = 60$, and we must have

$$\sum_{\pi \in \Pi^{14}} h^\pi = h^3 + h^4 + h^5 + h^6 = d^{14} = 40 \quad (5.6)$$

and

$$\sum_{\pi \in \Pi^{24}} h^\pi = h^1 + h^2 = d^{24} = 60. \quad (5.7)$$

Let's assume that the vehicles from each OD pair are divided evenly among all of the available paths, so $h^\pi = 10$ for each $\pi \in \Pi^{14}$ and $h^\pi = 30$ for each $\pi \in \Pi^{24}$. We can now use equation (5.2) to calculate the link flows. For instance the flow on link (1,2) is

$$\begin{aligned} \sum_{r \in Z} \left\{ \sum_{s \in Z} \left\{ \sum_{\pi \in \Pi^{rs}} \delta_{1,2}^\pi h^\pi \right\} \right\} &= \\ &= \left\{ \left\{ \delta_{(1,2)}^3 h^3 + \delta_{(1,2)}^4 h^4 + \delta_{(1,2)}^5 h^5 + \delta_{(1,2)}^6 h^6 \right\} \right\} + \left\{ \left\{ \delta_{(1,2)}^1 h^1 + \delta_{(1,2)}^2 h^2 \right\} \right\} = \\ &= \{ \{ 1 \times 10 + 1 \times 10 + 0 \times 10 + 0 \times 10 \} \} + \{ \{ 0 \times 30 + 0 \times 30 \} \} = 20. \end{aligned} \quad (5.8)$$

where the braces show how the summations “nest.” Remember, this is just a fancy way of picking the paths which use link (1,2), and adding their flows. The equation is for use in computer implementations or for large networks; when solving by hand, it's perfectly fine to just identify the paths using a particular link by inspection — in this case, only paths $\pi_1^{1,4}$ and $\pi_2^{1,4}$ use link (1,2). Repeating similar calculations, you should verify that $x_{13} = 20$, $x_{23} = 40$, $x_{32} = 10$, $x_{24} = 50$, and $x_{34} = 50$.

From these link flows we can get the link travel times by substituting the flows into the link performance functions, that is, $t_{12} = 10x_{12} = 200$, $t_{13} = 50 + x_{13} = 70$, $t_{23} = 10 + x_{23} = 50$, $t_{32} = 10 + x_{32} = 20$, $t_{24} = 50 + x_{24} = 100$, and $t_{34} = 10x_{34} = 500$. Finally, the path travel times can be obtained by either adding the travel times of their constituent links, or by applying equation (5.4). You should verify that $c^1 = 550$, $c^2 = 100$, $c^3 = 300$, $c^4 = 750$, $c^5 = 570$, and $c^6 = 190$.

The role of traffic assignment is to choose one path flow vector $\hat{\mathbf{h}}$ for purposes of forecasting and ranking alternatives, out of all of the feasible assignments in the network. An *assignment rule* is a principle used to determine this path flow vector. The most common assignment rule in practice is that the path flow vector should place all vehicles on a path with minimum travel time between their origins and destinations, although other rules are possible as well and will be discussed later in the book.

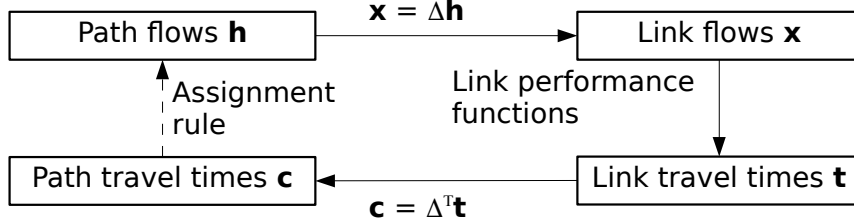


Figure 5.2: Traffic assignment as an iterative process.

5.1.1 Commentary

The equations and concepts mentioned in the previous subsection can be related as follows. Given a vector of path flows \mathbf{h} , we can obtain the vector of link flows \mathbf{x} from equation (5.3); from these, we can obtain the vector of link travel times \mathbf{t} by substituting each link's flow into its link performance function; from this, we can obtain the vector of path travel times \mathbf{c} from equation (5.5). This process is shown schematically in Figure 5.2.

The one component which does not have a simple representation is how to obtain path flows from path travel times using an assignment rule, “completing the loop” with the dashed line in the figure. This is actually the most complicated step, and answering it will require most of the remainder of the chapter. The main difficulty is that introducing some rule for relating path travel times to path flows creates a circular dependency: the path flows would depend on the path travel times, which depend on the link travel times, which depend on the link flows, which depend on the path flows, which depend on the path travel times and so on *ad infinitum*. We need to find a *consistent* solution to this process, which in this case means a set of path flows which remain unchanged when we go around the circuit: the path flows must be consistent with the travel times we obtain from those same path flows. Furthermore, the assignment rule must reflect the gaming behavior described in Section 1.3.

At this point, it is worthwhile to discuss the assumptions that go into the traffic assignment problem as stated here. The first assumption concerns the time scale over which the network modeling is occurring. This part of the book is focused entirely on what is called *static* network modeling, in which we assume that the network is close to a “steady state” during whatever length of time we choose to model (whether a peak hour, a multi-hour peak period, or a 24-hour model), and the link flows and capacities are measured with respect to the entire time period. That is, the capacity is the capacity *over the entire analysis period*, so if a facility has a capacity of 2200 veh/hr and we are modeling a three-hour peak period, we would use a capacity of 6600. Likewise, the link flows are the total flow over the three hours.

How long the analysis period should be is a matter of some balancing. Obviously, the longer the analysis period, the less accurate the steady state assumption is likely to be. In the extreme case of a 24-hour model, it is usually a (very big) stretch to assume that the congestion level will be the same over all 24 hours. On the other hand, choosing too short an analysis period can be problematic as well. In particular, if we are modeling a congested city or a large metropolitan area, trips from one end of the network to the other can easily take an hour or two, and it is good practice to have the analysis period be at least as long as most of the trips people are taking.

Properly resolving the issue of the “steady state” assumption requires moving to a *dynamic* traffic assignment (DTA) model. DTA models have the potential to more accurately model traffic, but are harder to calibrate, are more sensitive to having correct input data, and require more computer time. Further, DTA models end up requiring very different formulations and approaches. In short, while useful in some circumstances, they are not universally better than static models, and in any case they are surprisingly dissimilar. As a result, a full discussion of DTA models is deferred to the final part of this volume.

Another assumption we make is that *link and path flows can take any non-negative real value*. In particular, they are not required to be whole numbers, and there is no issue with saying that the flow on a link is 12.5 vehicles or that the number of vehicles choosing a path is, say, $\sqrt{2}$. This is often called the *continuum assumption*, because it treats vehicles as an infinitely-divisible fluid, rather than as a discrete number of “packets” which cannot be split. The reason for this assumption is largely practical — without the continuum assumption, traffic assignment problems become extremely difficult to solve, even on small networks. Further, from a practical perspective most links of concern have volumes in the hundreds or thousands of vehicles per hour, where the difference between fractional and integer values is negligible. Some also justify the continuum assumption by interpreting link and path flows to describe a long-term average of the flows, which may fluctuate to some degree from day to day.

5.2 Principle of User Equilibrium

The previous section introduced the main ideas and components of static traffic assignment, but did not provide much explanation of assignment rules other than claiming that a reasonable assignment rule results in all used routes connecting an origin and destination to have equal and minimal travel time, and that the “gaming” equilibrium idea described in Section 1.3 is relevant to route choice. This section explains this assignment rule in more detail.

Assignment rules are more complex than the other steps in traffic assignment shown in Figure 5.2, for two reasons. The first is that, being behavioral, it does not follow as mechanically as the other three steps, and has a more complicated answer. The other reason, alluding to the discussion in the previous section, is that adding this link creates a cycle of dependency. Linking path flows to path travel times, we now have four quantities which all depend on each other, and

Table 5.1: Potential criteria in choosing a route.

Low travel time
Reliable travel time
Low out-of-pocket cost (tolls, fuel, etc.)
Short distance
Bias toward (or away from) freeways
Low accident risk
Few potholes
Scenic view

untangling this cycle requires a little bit of thought.

To handle this, let's consider a simpler situation first. Rather than trying to find path flows (which require stating *every* driver's route choice), let's stick to a single driver. Why do they pick the route that they do? The number of potential paths in a network between two even slightly distant points in a network is enormous; in a typical city there are literally millions of possible routes one could theoretically take between an origin and destination. The vast majority of these are ridiculous, say, involving extraneous trips to the outskirts of the city and then back to the destination even though the origin is nearby.

Why are such paths ridiculous on their face? Table 5.1 gives a list of potential criteria which a desirable route would have, given in roughly decreasing order (at least according to my tastes for a trip to work or school). The most important times of day to properly model are the morning and evening peak periods, when most of the trips made are work-related commutes. For these types of trips, most people will choose the most direct route to their destination, where "direct" means some combination of travel time, cost, and distance (which are usually correlated). For simplicity, we'll use "low travel time" as our starting point, specified as Assumption 5.1.

Assumption 5.1. (Shortest path assumption.) *Each driver wants to choose the path between their origin and their destination with the least travel time.*

Notice that this principle does not include the impact a driver has on other drivers in the system, and a pithy characterization of Assumption 5.1 is that "people are greedy." If this seems unnecessarily pejorative, Section 5.3 shows a few ways that this principle can lead to suboptimal flow distributions on networks. Again, I emphasize that we adopt this assumption because we are modeling urban, peak period travel which is predominantly composed of work trips. If we were modeling, say, traffic flows around a national park during summer weekends, quality of scenery may be considerably more important than being able to drive at free-flow speed, and a different assumption would be needed. Further, the basic model developed in the first few weeks could really function just as well with cost or some other criterion, as long as it is *separable* and *additive* (that is, you can get the total value of the criterion by adding up its value on each link — the travel time of a route is the sum of the travel times

on its component links, the total monetary cost of a path is the sum of the monetary costs of each link, but the total scenic quality of a path may not be the sum of the scenic quality of each link), and can be expressed as a function of the link flows \mathbf{x} . So even though we will be speaking primarily of travel times, it is quite possible, and sometimes appropriate, to use other measures as well. A second assumption follows from modeling commute trips:

Assumption 5.2. *Drivers have perfect knowledge of link travel times.*

In reality, drivers' knowledge is not perfect — but commutes are typically habitual trips, so it is not unreasonable to assume that drivers are experienced and well-informed about congestion levels at different places in the network, at least along routes they might plausibly choose. We will later relax this assumption, but for now we'll take it as it greatly simplifies matters and is not too far from the truth for commutes. (Again, in a national park or other place with a lot of tourist traffic, this would be a poor assumption.)

Now, with a handle on individual behavior, we can try to scale up this assumption to large groups of travelers. What will be the resulting state if there are a large number of travelers who all want to take the fastest route to their destinations? For example, if you have to choose between two routes, one of which takes ten minutes and the other fifteen, you would always opt for the first one. If you are the only one traveling, this is all well and good. The situation becomes more complicated if there are others traveling. If there are ten thousand people making the same choice, and all ten thousand pick the first route, congestion will form and the travel time will rapidly increase. According to Assumption 5.2, drivers would become aware of this, and some people would switch from the first route to the second route, because the first would no longer be faster. This process would continue: as long as the first route is slower, people would switch away to the second route. If too many people switch to the second route, so the first becomes faster again, people would switch back.

With a little thought, it becomes clear that if there is *any* difference in the travel times between the two routes, people will switch from the slower route to the faster one. Note that Assumption 5.1 does not make any allowance for a driver being satisfied with a path which is a minute slower than the fastest path, or indeed a second slower, or even a nanosecond slower. Relaxing Assumption 5.1 to say that people may be indifferent as long as the travel time is “close enough” to the fastest path leads to an interesting, but more complicated line of research based on the concept of “bounded rationality,” which is discussed later, in Section 6.3.2. It is much simpler to assume that Assumption 5.1 holds strictly, in which case there are only three possible stable states:

1. Route 1 is faster, even when all of the travelers are using it.
2. Route 2 is faster, even when all of the travelers are using it.
3. Most commonly, neither route dominates the others. In this case, people use both Routes 1 and 2, and their travel times are *exactly equal*.

Because the third case is most common, this basic route choice model is called *user equilibrium*: the two routes are in equilibrium with each other. Why must the travel times be equal? If the first route was faster than the second, people would switch from the second to the first. This would decrease the travel time on the second route, and increase the travel time on the first, and people would continue switching until they were equal. The reverse is true as well: if the second route were faster, people would switch from the first route to the second, decreasing the travel time on the first route and increasing the travel time on the second. The *only* outcome where nobody has any reason to change their decision, is if the travel times are equal on both routes.

This is important enough to state again formally:

Corollary 5.1. (Principle of user equilibrium.) *Every used route connecting an origin and destination has equal and minimal travel time.*

Unused routes may of course have a higher travel time, and used routes connecting different origins and destinations may have different travel times, but any two used routes connecting the same origin and destination must have exactly the same travel times. Notice that we call this principle a corollary rather than an assumption: the real assumptions are the shortest path and full information assumptions. If you believe these are true, the principle of user equilibrium follows immediately and does not require you to assume anything more than you already have. The next section describes how to solve for equilibrium, along with a small example.

5.2.1 A trial-and-error solution method

We can develop a simple method for solving for path flows using the principle of user equilibrium, using nothing but the definition itself. For now, assume there is a single OD pair (r, s) . The method is as follows:

1. Select a set of paths $\hat{\Pi}^{rs}$ which you think will be used by travelers from this OD pair.
2. Write equations for the travel times of each path in $\hat{\Pi}^{rs}$ as a function of the path flows.
3. Solve the system of equations enforcing equal travel times on all of these paths, together with the requirement that the total path flows must equal the total demand d^{rs} .
4. Verify that this set of paths is correct; if not, refine $\hat{\Pi}^{rs}$ and return to step 2.

The first step serves to reduce a large potential set of paths to a smaller set of reasonable paths $\hat{\Pi}^{rs}$, which you believe to be the set of paths which will be used by travelers from r to s . Feel free to use engineering judgment here; if this is not the right set of paths, we'll discover this in step 4 and can adjust the set

accordingly. The second step involves applying equations (5.4) and (5.2). Write the equations for link flows as a function of the flow on paths in $\hat{\Pi}^{rs}$ (assuming all other paths have zero flow). Substitute these expressions for link flows into the link performance functions to get an expression for link travel times; then write the equations for path travel times as a function of link travel times.

If $\hat{\Pi}^{rs}$ truly is the set of used paths, all the travel times on its component paths will be equal. So, we solve a system of equations requiring just that. If there are n paths, there are only $n - 1$ independent equations specifying equal travel times. (If there are three paths, we have one equation stating paths one and two have equal travel time, and a second one stating paths two and three have equal travel time. An equation stating that paths one and three have equal travel time is redundant, because it is implied by the other two, and so it doesn't help us solve anything.) To solve for the n unknowns (the flow on each path), we need one more equation: the requirement that the sum of the path flows must equal the total flow from r to s (the “no vehicle left behind” equation requiring every vehicle to be assigned to a path). Solving this system of equations gives us path flows which provide equal travel times.

The last step is to verify that the set of paths is correct. What does this mean? There are two ways that $\hat{\Pi}^{rs}$ could be “incorrect”: either it contains paths that it shouldn't, or it omits a path that should be included. In the former case, you will end up with an infeasible solution (e.g., a negative or imaginary flow on one path), and should eliminate the paths with infeasible flows, and go back to the second step with a new set $\hat{\Pi}^{rs}$. In the latter case, you have a feasible solution and the travel times of paths in $\hat{\Pi}^{rs}$, but they are not minimal: you have missed a path which has a faster travel time than any of the ones in $\hat{\Pi}^{rs}$, so you need to include this path in $\hat{\Pi}^{rs}$ and again return to the second step.

Let's take a concrete example: Figure 5.3 shows 7000 travelers traveling from zone 1 to zone 2 during one hour, and choosing between the two routes mentioned above: route 1, with free-flow time 20 minutes and capacity 4400 veh/hr, and route 2, with free-flow time 10 minutes and capacity 2200 veh/hr. That means we have

$$t_1(x_1) = 20 \left(1 + 0.15 \left(\frac{x_1}{4400} \right)^4 \right), \quad (5.9)$$

$$t_2(x_2) = 10 \left(1 + 0.15 \left(\frac{x_2}{2200} \right)^4 \right). \quad (5.10)$$

This example is small enough that there is no real distinction between paths and links (because each path consists of a single link), so link flows are simply path flows ($x_1 = h_1$ and $x_2 = h_2$), and path travel times are simply link travel times ($c_1 = t_1$ and $c_2 = t_2$). As a starting assumption, we assume that both paths are used, so $\Pi^{12} = \{\pi_1, \pi_2\}$. We then need to choose the path flows h_1 and h_2 so that $t_1(h_1) = t_2(h_2)$ (equilibrium) and $h_1 + h_2 = 7000$ (no vehicle left behind). Substituting $h_2 = 7000 - h_1$ into the second delay function, the

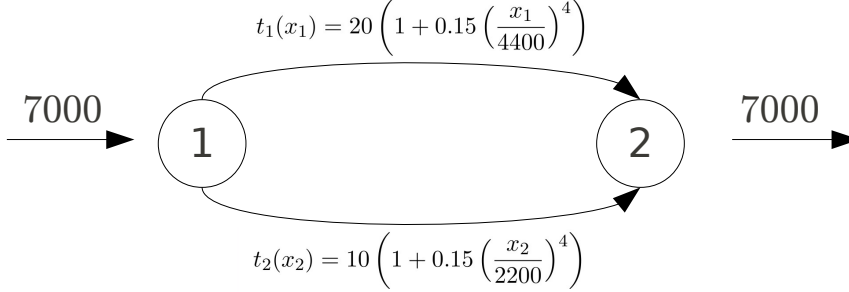


Figure 5.3: Small example using the two links.

equilibrium equation becomes

$$20 \left(1 + 0.15 \left(\frac{h_1}{4400} \right)^4 \right) = 10 \left(1 + 0.15 \left(\frac{7000 - h_1}{2200} \right)^4 \right). \quad (5.11)$$

Using a numerical equation solver, we find that equilibrium occurs for $h_1 = 3376$, so $h_2 = 7000 - 3376 = 3624$, and $t_1(x_1) = t_2(x_2) = 21.0$ minutes. Alternately, we can use a graphical approach. Figure 5.4 plots the travel time on *both* routes as a function of the flow on route 1 (because if we know the flow on route 1, we also know the flow on route 2). The point where they intersect is the equilibrium: $h_1 = 3376$, $t_1 = t_2 = 21.0$.

In the last step, we verify that the solution is reasonable (no paths have negative flow) and complete (there are no paths we missed which have a shorter travel time). Both conditions are satisfied, so we have found the equilibrium flow and stop.

Let's modify the problem slightly, so the travel time on link 1 is now

$$t_1(x_1) = 50 \left(1 + 0.15 \left(\frac{x_1}{4400} \right)^4 \right).$$

In this case, solving

$$50 \left(1 + 0.15 \left(\frac{h_1}{4400} \right)^4 \right) = 10 \left(1 + 0.15 \left(\frac{7000 - h_1}{2200} \right)^4 \right)$$

leads to a nonsensical solution: none of the answers involve real numbers without an imaginary part. The physical interpretation of this is that *there is no way to assign 7000 vehicles to these two paths so they have equal travel times*. Looking at a plot (Figure 5.5), we see that this happens because path 2 dominates path

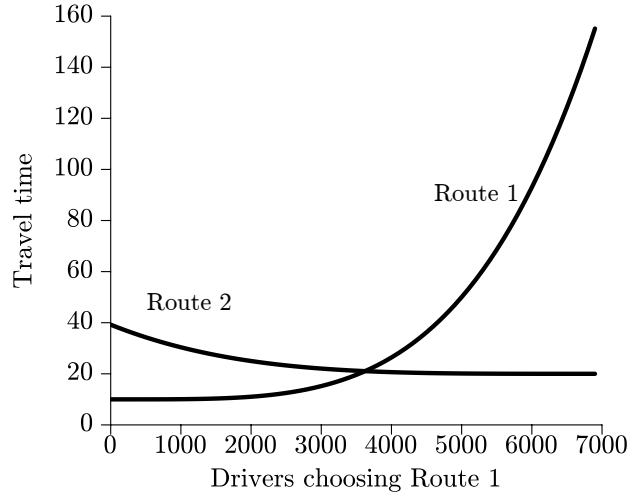


Figure 5.4: The equilibrium point lies at the intersection of the delay functions.

1: even with all 7000 vehicles on path 2, it has a smaller travel time. So, the solution to the modified problem is: $h_1 = 0$, $h_2 = 7000$, $t_1 = 50$, and $t_2 = 39.2$. This still satisfies the principle of user equilibrium: path 1 is not *used*, so it is fine for its travel time to not equal that of path 2.

To find this solution in the trial-and-error method, we would start by removing *both* paths (since imaginary path flows are not meaningful), and then adding in the shortest path at free-flow (since there must be at least one used path). When there is only one path in $\hat{\Pi}^{rs}$, the equilibrium state is easy to find, because the “system of equations” reduces trivially to placing all demand onto the single used path.

In more complicated networks, writing all of the equilibrium equations and solving them simultaneously is much too difficult. Instead, a more systematic approach is taken, and will be described in Chapter 7.

5.3 Three Motivating Examples

It is worth asking whether user equilibrium is the best possible condition. “Best” is an ambiguous term, but can be related towards our general goals as transportation engineers. User equilibrium probably does not lead to the flow pattern with, say, the best emissions profile, simply because there’s no reason to believe that the principle of user equilibrium has any connection whatsoever to emissions — it is based on people trying to choose fastest paths. But maybe it is related to congestion in some way, and you might find it plausible that each individual driver attempting to choose the best path for himself or herself would minimize congestion system-wide. This section presents three examples which

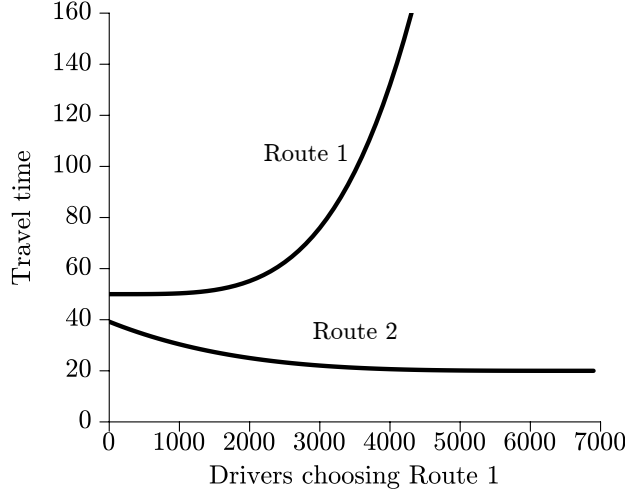


Figure 5.5: There is no intersection point: path 1 is dominated by path 2.

should shatter this innocent-seeming idea. In this section, we will not concern ourselves with how a feasible assignment satisfying this assignment rule is found (the trial-and-error method from the previous section would suffice), but will focus instead on how equilibrium can be used to evaluate the performance of potential alternatives.

In the first example, consider the network shown in Figure 5.6 where the demand between nodes 1 and 2 is $d^{12} = 30$ vehicles. Using h^\uparrow and h^\downarrow to represent the flows on the top and bottom paths in this network, the set of feasible assignments are the two-dimensional vectors $\mathbf{h} = [h^\uparrow \ h^\downarrow]$ which satisfy the conditions $h^\uparrow \geq 0$, $h^\downarrow \geq 0$, and $h^\uparrow + h^\downarrow = 30$. The figure shows the link performance functions for each of the two links in the network. Notice that the travel time on the top link is constant irrespective of the flow, while the travel time on the bottom link increases with its flow. However, at low values of flow, the bottom link is faster than the top one. This corresponds to a scenario where one link is shorter, but more subject to congestion, while the other link is longer but free of congestion.

With our stated assignment rule, the solution to the traffic assignment problem is $h^\uparrow = 25$, $h^\downarrow = 5$, because this will result in link flows of 25 and 5 on the top and bottom links, respectively, giving travel times of 50 on both paths. In this state, all vehicles in the network experience a travel time of 50.

Now, suppose that it is possible to improve one of the links in the network. The intuitive choice is to improve the link which is subject to congestion, say, changing its link performance function from $45 + x^\downarrow$ to $40 + \frac{1}{2}x^\downarrow$, as shown in Figure 5.7. (Can you see why this is called an improvement?) However, in this case the path flow solution which corresponds to the assignment rule is $h^\uparrow = 10$, $h^\downarrow = 20$, because this results in equal travel times on both the top and

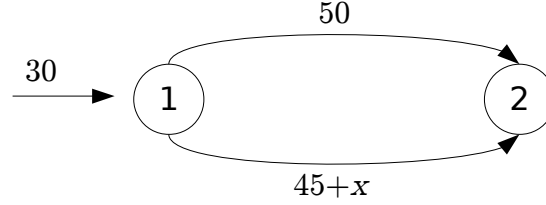


Figure 5.6: A two-link network.

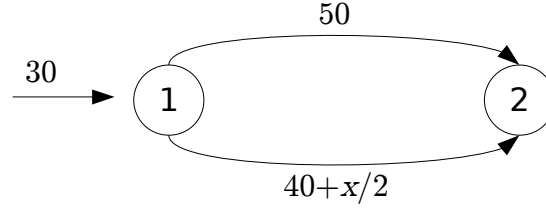


Figure 5.7: The two-link network with an improvement on the bottom link.

bottom paths. These travel times are still 50! Even though the bottom link was improved, the effect was completely offset by vehicles switching paths away from the top link and onto the bottom link.

In other words, improving a link (even the only congested link in a network) does not necessarily reduce travel times, because drivers can change their behavior in response to changes on the network. (If we could somehow force travelers to stay on the same paths they were using before, then certainly some vehicles would experience a lower travel time after the improvement is made.) This is called the *Knight-Pigou-Downs* paradox.

The second example was developed by Dietrich Braess, and shows a case where building a new roadway link can actually worsen travel times for all travelers in the network, after the new equilibrium is established. Assume we have the network shown in Figure 5.8a, with the link performance functions next to each link. The user equilibrium solution can be found by symmetry: since the top and bottom paths are exactly identical, the demand of six vehicles will evenly split between them. A flow of three vehicles on the two paths leads to flow of three vehicles on each of the four links; substituting into the link

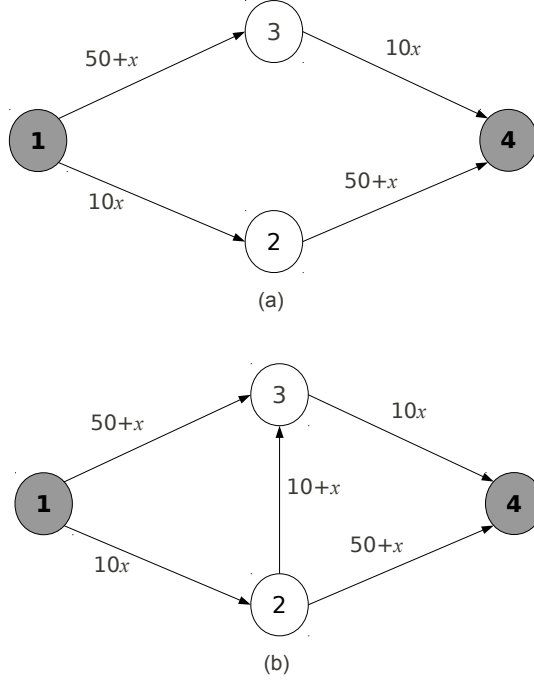


Figure 5.8: Braess network, before and after construction of a new link.

performance functions gives travel times of 53 on $(1, 3)$ and $(2, 4)$ and 30 on $(1, 2)$ and $(3, 4)$, so the travel time of both paths is 83, and the principle of user equilibrium is satisfied.

Now, let's modify the network by adding a new link from node 2 to node 3, with link performance function $10 + x_{23}$ (Figure 5.8b). We have added a new path to the network; let's label these as follows. Path 1 is the top route $[1, 3, 4]$; path 2 is the middle route $[1, 2, 3, 4]$, and path 3 is the bottom route $[1, 2, 4]$. Paths 1 and 3 each have a demand of three vehicles and a travel time of 83. Path 2, on the other hand, has a flow of zero vehicles and a travel time of 70, so the principle of user equilibrium is violated: the travel time on the used paths is equal, but not minimal.

Assumption 5.1 suggests that someone will switch their path to take advantage of this lower travel time; let's say someone from path 1 switches to path 2, so we have $h_1 = 2$, $h_2 = 1$, and $h_3 = 3$. From this we can predict new link flows: $x_{12} = 4$, $x_{13} = 2$, $x_{23} = 1$, $x_{24} = 3$, $x_{34} = 3$. Substituting into link performance functions gives new link travel times: $t_{12} = 40$, $t_{13} = 52$, $t_{23} = 11$, $t_{24} = 53$, $t_{34} = 30$, and finally we can recover the new path travel times: $c_1 = 82$, $c_2 = 81$, and $c_3 = 93$.

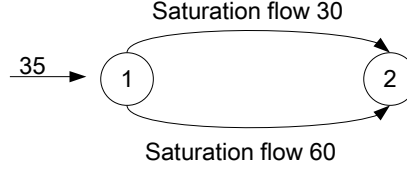


Figure 5.9: Network for Smith's paradox.

This is still not an equilibrium; perhaps someone from path 3 will switch to path 2 in an effort to save 12 minutes of time. So now $h_1 = h_2 = h_3 = 2$, the link flows are $x_{12} = x_{34} = 4$ and $x_{13} = x_{23} = x_{24} = 2$, the travel times are $t_{12} = t_{34} = 40$, $t_{13} = t_{24} = 52$ and $t_{23} = 12$, and the path travel times are $c_1 = c_2 = c_3 = 92$. We have found the new equilibrium, and it is unconditionally worse than the old one. Before adding the new link, each driver had a travel time of 83; now every driver has a travel time of 92. Nobody is better off. Everyone is worse off. This is the famous *Braess paradox*.

The third example, adapted from M. J. Smith, shows how the traffic assignment problem can provide more insight into other traffic engineering problems such as signal timing. Consider the network shown in Figure 5.9. Drivers choose one of two links, which join at a signalized intersection at the destination. Assume that each of these links has the same free-flow time ($t_0^\uparrow = t_0^\downarrow = 1$ min), but that the saturation flow of the bottom link is twice that of the top link ($s^\uparrow = 30$ veh/min and $s^\downarrow = 60$ veh/min). Vehicles enter the network at a demand level of $d = 35$ veh/min. The signal has a cycle length of $C = 1$ min and, for simplicity's sake, assume that there is no lost time so that the green times allocated to the top and bottom links equal the cycle length: $G^\uparrow + G^\downarrow = 1$.

In traditional traffic signal analysis, the *capacity* u of an approach is the saturation flow scaled by the proportion of green time given to that approach, so $u_i = s_i(G_i/C)$ where i can refer to either the top or bottom link. The *degree of saturation* X for an approach is the ratio of the link flow to the capacity, so $X_i = x_i/u_i = (x_i C)/(s_i G_i)$. With these quantities, the total travel time on each link (free-flow time plus signal delay) can be written as

$$t_i = 1 + \frac{9}{20} \left[\frac{C(1 - G_i/C)^2}{1 - x_i/s_i} + \frac{X_i^2}{x_i(1 - X_i)} \right]. \quad (5.12)$$

Assume that initially, the signal is timed such that $G^\uparrow = 40$ sec and $G^\downarrow = 20$ sec. Then the formulas for delay on the top and bottom links are functions of their flows alone, obtained from (5.12) by substituting the corresponding green times into the expressions for X_i , and the equilibrium solution can be obtained by solving the equations $t^\uparrow = t^\downarrow$ and $x^\uparrow + x^\downarrow = 35$ simultaneously. With these green times, the equilibrium solution occurs when $x^\uparrow = 23.6$ veh/min and

$x^\downarrow = 11.4$ veh/min, and the reader can confirm that the travel times on the two links are $t^\uparrow = t^\downarrow = 2.11$ min.

So far, so good. Now assume that it has been a long time since the signal was last retimed, and a traffic engineer decides to check on the signal and potentially change the timing. A traditional rule in traffic signal timing is that the green time given to an approach should be proportional to the degree of saturation. However, with the given solution, the degrees of saturation are $X^\uparrow = 0.982$ and $X^\downarrow = 0.953$, which are unequal — the top link is given slightly less green time than the equisaturation rule suggests, and the bottom link slightly more. Therefore, the engineer changes the green times to equalize these degrees of saturation, which occurs if $G^\uparrow = 48.3$ sec and $G^\downarrow = 11.7$ sec, a smallish adjustment. If drivers could be counted on to remain on their current routes, all would be well. However, changing the signal timing changes the delay formulas (5.12) on the two links. Under the assumption that drivers always seek the shortest path, the equilibrium solution will change as drivers swap from the (now longer) path to the (now shorter) path. Re-equating the travel time formulas, the flow rates on the top and bottom links are now $x^\uparrow = 23.8$ veh/min and $x^\downarrow = 11.2$ veh/min, with equal delays $t^\uparrow = t^\downarrow = 2.26$ min on each link. Delay has actually increased, because the signal re-timing (aimed at reducing delay) did not account for the changes in driver behavior after the fact.

A bit surprised by this result, our diligent traffic engineer notes that the degrees of saturation are still unequal, with $X^\uparrow = 0.984$ and $X^\downarrow = 0.959$, actually further apart than before the first adjustment. Undeterred, the engineer changes the signal timings again to $G^\uparrow = 48.5$ s and $G^\downarrow = 11.5$ s, which results in equal degrees of saturation under the new flows. But by changing the green times, the delay equations have changed, and so drivers re-adjust to move toward the shorter path, leading to $x^\uparrow = 23.9$ veh/min and $x^\downarrow = 11.1$ veh/min, and new delays of 2.43 minutes on each approach, even higher than before!

You can probably guess what happens from here, but Table 5.2 tells the rest of the story. As our valiant engineer stubbornly retimes the signals in a vain attempt to maintain equisaturation, flows always shift in response. Furthermore, the delays grow faster and faster, asymptotically growing to infinity as more and more adjustments are made. The moral of the story? *Changing the network will change the paths that drivers take*, and “optimizing” the network without accounting for how these paths will change is naive at best, and counterproductive at worst.

Table 5.2: Evolution of the network as signals are iteratively retimed.

Iteration	G^\uparrow (s)	G^\downarrow	x^\uparrow (v/min)	x^\downarrow	t^\uparrow (min)	t^\downarrow	X^\uparrow	X^\downarrow
0	48	12	23.6	11.4	2.11	2.11	0.982	0.953
1	48.3	11.7	23.8	11.2	2.26	2.26	0.984	0.959
2	48.5	11.5	23.9	11.1	2.43	2.43	0.986	0.965
3	48.7	11.3	24.1	10.9	2.63	2.63	0.988	0.969
4	48.9	11.1	24.2	10.8	2.86	2.86	0.99	0.974
5	49.1	10.9	24.3	10.7	3.12	3.12	0.991	0.977
10	49.5	10.5	24.7	10.3	5.11	5.11	0.996	0.989
20	49.88	10.12	24.91	10.09	16.58	16.58	0.9988	0.9971
50	49.998	10.002	24.998	10.002	855.92	855.93	0.99998	0.99995
∞	50	10	25	10	∞	∞	1	1

5.4 Exercises

- [15] Expand the diagram of Figure 1.4 to include a “government model” which reflects public policy and regulation regarding tolls. What would this “government” agent influence, and how would it be influenced?
- [23] How realistic do you think link performance functions are? Name at least two assumptions they make about traffic congestion, and comment on how reasonable you think they are.
- [23] How realistic do you think the principle of user equilibrium is? Name at least two assumptions it makes (either about travelers or congestion), and comment on how reasonable you think they are.
- [12] Develop analogies to the equilibrium principle to represent the following phenomena: (a) better restaurants tend to be more crowded; (b) homes in better school districts tend to be more expensive; (c) technologies to improve road safety (such as antilock brakes) can increase reckless driving. Specify what assumptions you are making in these analogies.
- [34] (*All-or-nothing solutions are corner points.*) The point \mathbf{x} is a *corner point* of the convex set X if it is impossible to write $\mathbf{x} = \lambda \mathbf{y} + (1 - \lambda)\mathbf{z}$ for $\mathbf{y} \in X$, $\mathbf{z} \in X$, and $\lambda \in (0, 1)$. Show that if H is the set of feasible path flows, \mathbf{h} is a corner point if and only if it is an *all-or-nothing solution* (a solution where all of the flow from each OD pair is on a single path). Repeat for the set of feasible link flows.
- [20] For each of the following games, list all equilibria (or state that none exist). In which of these games do equilibria exist; in which are the equilibria unique; and in which are there some equilibria which are inefficient? These games are all played by two players A and B: A chooses the row and B chooses the column, and each cell lists the payoffs to A and B, in that order.

(a)	L	R	(b)	L	R
U	(8, 13)	(5, 14)	U	(12, 12)	(2, 10)
D	(10, 10)	(7, 12)	D	(10, 2)	(4, 5)

(c)	L	R	(d)	L	R
U	(8, 6)	(10, 8)	U	(4, 5)	(5, 6)
D	(2, 3)	(8, 4)	D	(3, 4)	(6, 3)

- [14] There are three routes available to travelers between a single OD pair. For each case below, you are given the travel time t on each route, as well as the number of travelers x using each route. Indicate whether each case satisfies the principle of user equilibrium. If this principle is violated, explain why.

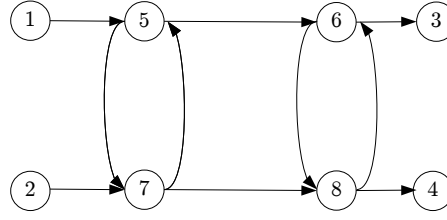


Figure 5.10: Network for Exercise 8.

- (a) $t_1 = 45$, $t_2 = 60$, and $t_3 = 30$ while $x_1 = 0$, $x_2 = 50$, and $x_3 = 0$.
 - (b) $t_1 = 45$, $t_2 = 30$, and $t_3 = 30$ while $x_1 = 0$, $x_2 = 50$, and $x_3 = 25$.
 - (c) $t_1 = 45$, $t_2 = 30$, and $t_3 = 30$ while $x_1 = 0$, $x_2 = 50$, and $x_3 = 0$.
 - (d) $t_1 = 15$, $t_2 = 30$, and $t_3 = 45$ while $x_1 = 30$, $x_2 = 0$, and $x_3 = 0$.
 - (e) $t_1 = 15$, $t_2 = 30$, and $t_3 = 45$ while $x_1 = 30$, $x_2 = 20$, and $x_3 = 10$.
 - (f) $t_1 = 15$, $t_2 = 30$, and $t_3 = 30$ while $x_1 = 0$, $x_2 = 50$, and $x_3 = 30$.
8. [26] The network in Figure 5.10 has 8 nodes, 12 links, and 4 zones. The travel demand is $d^{13} = d^{14} = 100$, $d^{23} = 150$, and $d^{24} = 50$. The dashed links have a constant travel time of 10 minutes regardless of the flow on those links; the solid links have the link performance function $15 + x/20$ where x is the flow on that link.
- (a) For each of the four OD pairs with positive demand, list all acyclic paths connecting that OD pair. In total, how many such paths are in the network?
 - (b) Assume that the demand for each OD pair is divided evenly among all of the acyclic paths you found in part (a) for that OD pair. What are the resulting link flow vector \mathbf{x} , travel time vector \mathbf{t} , and path travel time vector \mathbf{C} ?
 - (c) Does that solution satisfy the principle of user equilibrium?
 - (d) What is the total system travel time?
9. [32] Consider the network and OD matrix shown in Figure 5.11. The travel time on *every* link is $10 + x/100$, where x is the flow on that link. Find the link flows and link travel times which satisfy the principle of user equilibrium.
10. [36] Find the equilibrium path flows, path travel times, link flows, and link travel times on the Braess network (Figure 5.8b) when the travel demand from node 1 to node 4 is (a) 2; (b) 7; and (c) 10.
11. [52] In the Knight-Pigou-Downs network, the simple and seemingly reasonable heuristic to improve the congested link (i.e., changing its link

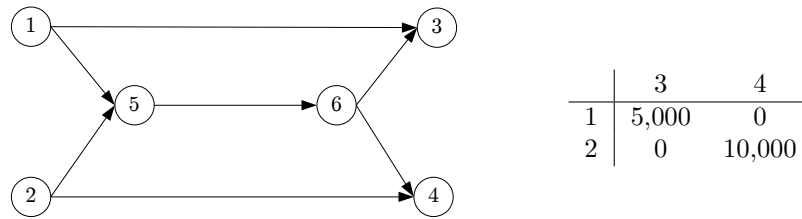


Figure 5.11: Network and OD matrix for Exercise 9.

performance function to lower the free-flow time) did not help. Can you identify an equally simple heuristic that would lead you to improve the other link? (Ideally, such a heuristic would be applicable in many kinds of networks, not just one with an uncongestible link.)

12. [20] Give nontechnical explanations of why uniqueness, efficiency, and existence of equilibrium solutions have practical implications, not just theoretical ones. Concrete examples may be helpful.
13. [77] In the trial-and-error method, identify several different strategies for choosing the initial set of paths. Test these strategies on networks of various size and complexity. What conclusions can you draw about the effectiveness of these different strategies, and the amount of effort they involve?
14. [55] The trial-and-error method generally involves the solution of a system of nonlinear equations. *Newton's method* for solving a system of n nonlinear equations is to move all quantities to one side of the equation, expressing the resulting equations in the form $\mathbf{F}(\mathbf{x}) = \mathbf{0}$ where \mathbf{F} is a vector-valued function mapping the n -dimensional vector of unknowns \mathbf{x} to another n -dimensional vector giving the value of each equation. An initial guess is made for \mathbf{x} , which is then updated using the rule $\mathbf{x} \leftarrow \mathbf{x} - (J\mathbf{F}(\mathbf{x}))^{-1}\mathbf{F}(\mathbf{x})$, where $(J\mathbf{F}(\mathbf{x}))^{-1}$ is the inverse of the Jacobian matrix of \mathbf{F} , evaluated at \mathbf{x} . This process continues until (hopefully) \mathbf{x} converges to a solution. A *quasi-Newton method* approximates the Jacobian with a diagonal matrix, which is equal to the Jacobian along the diagonal and zero elsewhere, which is much faster to calculate and convenient to work with. Extend your experiments from Exercise 13 to see whether Newton or quasi-Newton methods work better.
15. [84] What conditions on the link performance functions are needed for Newton's method (defined in the previous exercise) to converge to a solution of the system of equations? What about the quasi-Newton method? Assume that the network consists of n parallel links connecting a single origin to a single destination. Can you guarantee that the solution to the system of equations only involves real numbers, and not complex numbers?

Hint: it may be useful to redefine the link performance functions for negative x values. This should not have any effect on the ultimate equilibrium solution, since negative link flows are infeasible, but cleverly redefining the link performance functions in this region may help show convergence.

16. [96] Repeat the previous exercise, but for a general network with any number of links and nodes, and where paths may overlap.

Chapter 6

The Traffic Assignment Problem

This chapter formalizes the user equilibrium traffic assignment problem defined in Chapter 5. Using the mathematical language from Chapter 3, we are prepared to model and solve equilibrium problems even on networks of realistic size, with tens of thousands of links and nodes. Section 6.1 begins with fixed point, variational inequality, and optimization formulations of the user equilibrium problem. Section 6.2 then introduces important existence and uniqueness properties of the user equilibrium assignment, as was first introduced with the small two-player games of Section 1.3. Section 6.3 names several alternatives to the user equilibrium rule for assigning traffic to a network, including the system optimal principle, and the idea of bounded rationality. This chapter concludes with Section 6.4, exploring the inefficiency of the user equilibrium rule and unraveling the mystery of the Braess paradox.

6.1 Mathematical Formulations

The previous section introduced a trial-and-error method for solving the user equilibrium problem. To solve equilibrium on larger networks, we need something more sophisticated than trial-and-error, but doing so requires a more convenient representation of the principle of user equilibrium in mathematical notation, using the tools defined in the preceding chapters. This section does so by first formulating the equilibrium problem as the solution to a variational inequality; as the solution to a fixed point problem based on the variational inequality; and as the solution to a convex optimization problem. The reason for presenting all three of these formulations (rather than just one) is that each is useful in different ways. Just as true fluency in a language requires being able to explain the same concept in different ways, fluency in traffic assignment problems requires familiarity with all of these ways of understanding the principle of user equilibrium. (It is also possible to formulate the principle of user

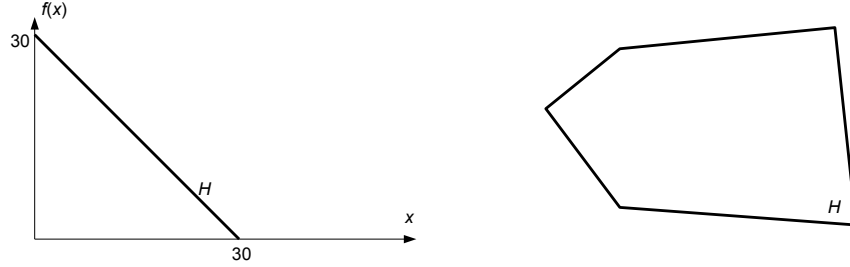


Figure 6.1: Feasible path-flow sets H for the simple two-link network (left) and a schematic representing H for more complex problems (right).

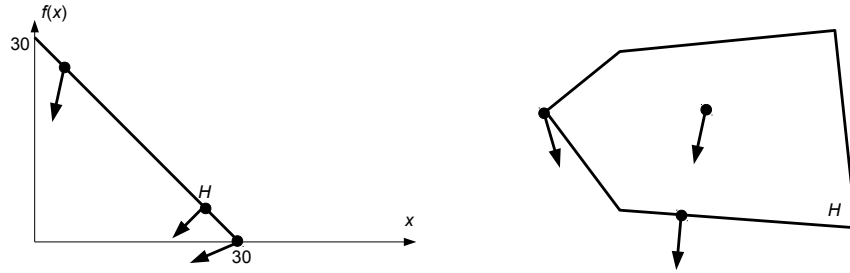


Figure 6.2: “Force” vectors $-\mathbf{c}$ for the simple two-link network (left) and a schematic representing H for more complex problems (right).

equilibrium as a fixed point problem without referring to a variational inequality through the use of multifunctions, as described in the optional subsection 6.1.1.)

The variational inequality formulation proceeds as follows. For each feasible path flow vector, we can associate a direction representing travelers’ desire for lower travel time paths. Figure 6.1 shows the feasible set for the two-link equilibrium problem on the left, and a conceptual schematic of the feasible set for the general equilibrium problem on the right. Notice that in all cases the feasible set is convex and has no “gaps.” For each feasible vector \mathbf{h} , we can calculate the path travel time vector \mathbf{c} associated with these path flows. Associate with each point a vector pointing in the *opposite* direction as \mathbf{c} , that is, each point \mathbf{h} is associated with the direction $-\mathbf{c}(\mathbf{h})$. You can think about this as a force acting on the point \mathbf{h} . The interpretation is that the path flow vector is “pulled” in the direction of decreasing travel times, and we will shortly show that an equilibrium is a point which is unmoved by this force. Of course, this intuitive interpretation must be established mathematically as well, and we will show that a vector \mathbf{h} satisfies the principle of user equilibrium if, and only if, the force criterion is satisfied.

A few points are shown in Figure 6.2 as examples. Recall that a non-corner point is unmoved by such a force only if the force is perpendicular to the bound-

ary of the feasible set at that point. At a corner point, the force must make a right or obtuse angle with all boundary directions. So we can characterize stable path flow vectors $\hat{\mathbf{h}}$ as points where the force $-\mathbf{c}(\hat{\mathbf{h}})$ makes a right or obtuse angle with any feasible direction $\mathbf{h} - \hat{\mathbf{h}}$, where \mathbf{h} is any other feasible path flow vector. Recalling vector operations, saying that two vectors make a right or obtuse angle with each other is equivalent to saying that their dot product is nonpositive. Thus, $\hat{\mathbf{h}}$ is a stable point if it satisfies

$$-\mathbf{c}(\hat{\mathbf{h}}) \cdot (\mathbf{h} - \hat{\mathbf{h}}) \leq 0 \quad \forall \mathbf{h} \in H \quad (6.1)$$

or, equivalently,

$$\mathbf{c}(\hat{\mathbf{h}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) \leq 0 \quad \forall \mathbf{h} \in H. \quad (6.2)$$

This is a variational inequality (VI) in the form shown in Section 3.2.2. We now prove that solutions of this VI correspond to equilibria, as shown by the next result:

Theorem 6.1. *A path flow vector $\hat{\mathbf{h}}$ solves the variational inequality (6.2) if and only if it satisfies the principle of user equilibrium.*

Proof. The theorem can equivalently be written “a path flow vector $\hat{\mathbf{h}}$ does not solve (6.2) if and only if it does not satisfy the principle of user equilibrium,” which is easier to prove. Assume $\hat{\mathbf{h}}$ does not solve (6.2). Then there exists some $\mathbf{h} \in H$ such that $\mathbf{c}(\hat{\mathbf{h}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) > 0$, or equivalently $\mathbf{c}(\hat{\mathbf{h}}) \cdot \hat{\mathbf{h}} > \mathbf{c}(\hat{\mathbf{h}}) \cdot \mathbf{h}$. Now, $\mathbf{c}(\hat{\mathbf{h}}) \cdot \hat{\mathbf{h}}$ is the total system travel time when the path flows are $\hat{\mathbf{h}}$, and $\mathbf{c}(\hat{\mathbf{h}}) \cdot \mathbf{h}$ is the total system travel time if the travel times were held constant at $\mathbf{c}(\hat{\mathbf{h}})$ even when the flows changed to \mathbf{h} . For the latter to be strictly less than the former, switching from $\hat{\mathbf{h}}$ to \mathbf{h} must have reduced at least one vehicle’s travel time even though the path travel times did not change. This can only happen if that vehicle was not on a minimum travel time path to begin with, meaning that $\hat{\mathbf{h}}$ does not satisfy user equilibrium.

Conversely, assume that $\hat{\mathbf{h}}$ is not a user equilibrium. Then there is some OD pair (r, s) and path π such that $\hat{h}_\pi > 0$ even though $c_\pi(\hat{\mathbf{h}}) > \min_{\pi' \in \Pi^{rs}} c_{\pi'}(\hat{\mathbf{h}})$. Let $\hat{\pi}$ be a minimum travel time path for this OD pair. Create a new path flow vector \mathbf{h} which is the same as $\hat{\mathbf{h}}$ except that \hat{h}_π is reduced by some small positive amount ϵ and $\hat{h}_{\hat{\pi}}$ is increased by ϵ . As long as $0 < \epsilon < \hat{h}_\pi$ the new point \mathbf{h} remains feasible. By definition all the components of $\hat{\mathbf{h}} - \mathbf{h}$ are equal to zero, except the component for π is ϵ and the component for $\hat{\pi}$ is $-\epsilon$. So $\mathbf{c}(\hat{\mathbf{h}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) = \epsilon(c_\pi(\hat{\mathbf{h}}) - c_{\hat{\pi}}(\hat{\mathbf{h}})) > 0$, so $\hat{\mathbf{h}}$ does not solve (6.2). \square

This variational inequality formulation leads directly to a fixed point formulation, as was discussed in Section 3.2.2. Theorem 6.1 shows that the equilibrium path flow vectors are exactly the solutions of the VI (6.2), that is, the stable points with respect to the fictitious force $-\mathbf{c}$. Therefore, an equilibrium vector $\hat{\mathbf{h}}$ is a fixed point of the function $\text{proj}_H(\mathbf{h} - \mathbf{c}(\mathbf{h}))$, that is, for all equilibrium solutions $\hat{\mathbf{h}}$ we have

$$\hat{\mathbf{h}} = \text{proj}_H(\hat{\mathbf{h}} - \mathbf{c}(\hat{\mathbf{h}})). \quad (6.3)$$

Finally, user equilibrium can also be shown as the solution to a convex optimization problem. The main decision variables are the path flows \mathbf{h} , and we need to choose one from the set of feasible assignments H ; the trick is identifying an appropriate function $f(\mathbf{h})$ in terms of the path flows, which is minimized when the path flows satisfy the principle of user equilibrium. So, the constraint set of our optimization problem consists of the equations defining H :

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (6.4)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (6.5)$$

Different assignment rules will lead to different objective functions. It turns out that the objective function corresponding to the principle of user equilibrium has a somewhat unintuitive form. Therefore, we will derive this function in the same way that it was originally derived, by working backwards. *Rather than writing down the optimization problem, and then deriving the optimality conditions, we start by writing down the optimality conditions we want, then determining what kind of optimization problem has that form.*

Writing the (unknown!) objective function as $f(\mathbf{h})$, using the procedures described in Chapter 4 we can Lagrangianize the constraints (6.4) introducing multipliers κ^{rs} for each OD pair, providing the following optimality conditions:

$$\frac{\partial f}{\partial h^\pi} - \kappa^{rs} \geq 0 \quad \forall (r, s) \in Z^2, \pi \in \Pi^{rs} \quad (6.6)$$

$$h^\pi \left(\frac{\partial f}{\partial h^\pi} - \kappa^{rs} \right) = 0 \quad \forall (r, s) \in Z^2, \pi \in \Pi^{rs} \quad (6.7)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (6.8)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (6.9)$$

The last two of these are simply (6.4) and (6.5), requiring that the solution be feasible. The condition (6.7) is the most interesting of these, requiring that the product of each path's flow and another term involving f must be zero. For this product to be zero, either h^π must be zero, or $\frac{\partial f}{\partial h^\pi} = \kappa^{rs}$ must be true; and by (6.6), for all paths $\frac{\partial f}{\partial h^\pi} \geq \kappa^{rs}$. That is, in a solution to this problem, whenever h^π is positive we have $\frac{\partial f}{\partial h^\pi} = \kappa^{rs}$; if a path is unused, then $\frac{\partial f}{\partial h^\pi} \geq \kappa^{rs}$ — in other words, at optimality all used paths connecting OD pair (r, s) must have equal and minimal $\frac{\partial f}{\partial h^\pi}$, which is equal to κ^{rs} . According to the principle of user equilibrium, all used paths have equal and minimal travel time... so if we can choose a function $f(\mathbf{h})$ such that $\frac{\partial f}{\partial h^\pi} = c^\pi$, we are done!

So, the objective function must involve some integral of the travel times. A first guess might look something like

$$f(\mathbf{h}) = \sum_{\pi \in \Pi} \int_{\text{?}}^{\text{?}} c^\pi dh^\pi, \quad (6.10)$$

where the bounds of integration and other details are yet to be determined. The trouble is that c^π is not a function of h^π alone: the travel time on a path depends on the travel times on other paths as well, so the partial derivative of this function with respect to h^π will not simply be c^π , but contain other terms as well. However, these interactions are not arbitrary, but instead occur where paths overlap, that is, where they share common links. In fact, if we try writing a similar function to our guess but in terms of *link flows*, instead of *path flows*, we will be done.

To be more precise, let $\mathbf{x}(\mathbf{h})$ be the link flows as a function of the path flows \mathbf{h} , as determined by equation (5.2). Then the function

$$f(\mathbf{h}) = \sum_{(i,j) \in A} \int_0^{x_{ij}(\mathbf{h})} t_{ij}(x) dx \quad (6.11)$$

satisfies our purposes. To show this, calculate the partial derivative of f with respect to the flow on an arbitrary path π , using the fundamental theorem of calculus and the chain rule:

$$\frac{\partial f}{\partial h^\pi} = \sum_{(i,j) \in A} t_{ij}(x_{ij}(\mathbf{h})) \frac{\partial x_{ij}}{\partial h^\pi} = \sum_{(i,j) \in A} \delta_{ij}^\pi t_{ij}(x_{ij}(\mathbf{h})) = c^\pi \quad (6.12)$$

where the last two equalities respectively follow from differentiating (5.2) and from (5.4).

Finally, we can clean up the notation a bit by simply introducing the link flows \mathbf{x} as a new set of decision variables, adding equations (5.2) as constraints to ensure they are consistent with the path flows. This gives the following optimization problem, first formulated by Martin Beckmann:

$$\min_{\mathbf{x}, \mathbf{h}} \quad \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (6.13)$$

$$\text{s.t.} \quad x_{ij} = \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \quad \forall (i,j) \in A \quad (6.14)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r,s) \in Z^2 \quad (6.15)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (6.16)$$

Section 6.2 shows that the objective function is convex, and the feasible region is a convex set, so this problem will be relatively easy to solve.

6.1.1 Multifunctions and application to network equilibrium (*)

(This section is optional and can be skipped. However if you are curious about the “other road” to using fixed points to show equilibrium existence, read this section to learn about multifunctions and Kakutani’s theorem.)

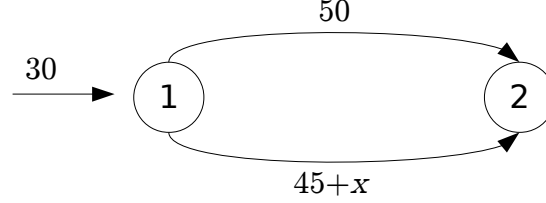


Figure 6.3: A simple two-link network for demonstration.

In the transit ridership example of Section 3.2.1, we could formulate the solution to the problem as a fixed point problem directly, without first going through a variational inequality. The traffic assignment problem is slightly more complex, because at equilibrium all used routes will have equal travel time. So, if we try to apply the same technique as in the transit ridership problem, we will run into a difficulty with “breaking ties.” To see this concretely, consider the two-route network in Figure 6.3, and let \mathbf{h} and \mathbf{c} be the vectors of route flows and route travel times. As Figure 1.3 suggests, we can try to define two functions: $\mathbf{H}(\mathbf{c})$ gives the vector of path flows representing route choices *if the travel times were fixed at \mathbf{c}* , and $\mathbf{C}(\mathbf{h})$ gives the vector of path travel times when the path flows are \mathbf{h} . The function \mathbf{C} is clearly defined: simply calculate the cost on each route using the link performance functions. However, if both paths are used at equilibrium, then $c_1 = c_2$, which means that *any* path flow vector would be a valid choice for $\mathbf{H}(\mathbf{c})$. Be sure you understand this difficulty: because \mathbf{H} assumes that the travel times are fixed, if they are equal any path flow vector is consistent with our route choice assumptions. If we relax the assumption that the link performance functions are fixed, then \mathbf{H} is no simpler than solving the equilibrium problem in the first place.

To resolve this, we introduce the concept of a *multifunction*.¹ Recall that a regular function from X to Y associates each value in its domain X with exactly one value in the set Y . A multifunction, on the other hand, associates each value in the domain X with some subset of Y . If F is such a multifunction the notation $F : X \rightrightarrows Y$ can be used. For example, consider the multifunction $R(\mathbf{c})$. If $c_1 < c_2$ (and these travel times were fixed), then the only consistent path flows are to have everybody on path 1. Likewise, if $c_1 > c_2$, then everyone would have to be on path 2. Finally, if $c_1 = c_2$ then people could split in any proportion while satisfying the rule that drivers are using least-time paths. That

¹Also known by many other names, including *correspondence*, *point-to-set map*, *set-valued map*, or *set-valued function*.

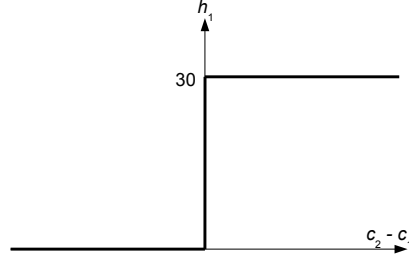


Figure 6.4: Multifunction corresponding to the two-link network; notice multiple values when $c_1 = c_2$.

is,

$$R(c_1, c_2) = \begin{cases} \left\{ \begin{bmatrix} 30 & 0 \end{bmatrix} \right\} & c_1 < c_2 \\ \left\{ \begin{bmatrix} 0 & 30 \end{bmatrix} \right\} & c_1 > c_2 \\ \left\{ \begin{bmatrix} h & 30 - h \end{bmatrix} : h \in [0, 30] \right\} & c_1 = c_2 \end{cases} \quad (6.17)$$

The notation in equation (6.17) is chosen very carefully to reflect the fact that H is a multifunction, which means that its values are *sets*, not a specific number. In the first two cases, the set only consists of a single element, so this distinction may seem a bit pedantic; but in the latter case, the set contains a whole range of possible values. In this way, multifunctions generalize the concept of a function by allowing R to take multiple “output” values for a single input. This is graphically represented in Figure 6.4, using the fact that the function can be parameterized in terms of a single variable $c_2 - c_1$.

Fixed-point problems can be formulated for multifunctions as well as for functions. If F is an arbitrary multifunction defined on the set K , then a fixed point of F is a point x such that $x \in F(x)$. Note the use of set inclusion \in rather than equality $=$ because F can associate multiple values with x . Just as Brouwer’s theorem guarantees existence of fixed points for functions, Kakutani’s theorem guarantees existence of fixed points for multifunctions under general conditions:

Theorem 6.2. (Kakutani). *Let $F : K \rightrightarrows K$ be a multifunction (from the set K to itself), where K is convex and compact. If F has a closed graph and $F(x)$ is nonempty and convex for all $x \in K$, then there is at least one point $x \in K$ such that $x \in F(x)$.*

The new terminology here is a *closed graph*; we say that the multifunction F has a closed graph if the set $\{(\mathbf{x}, \mathbf{y}) : \mathbf{x} \in X, \mathbf{y} \in F(\mathbf{x})\}$ is closed. Again, each of these conditions is necessary; you might find it helpful to visualize these conditions geometrically similar to Figure 3.9.

As a result of Kakutani’s theorem, we know that an equilibrium solution must always exist in the traffic assignment problem. Let H denote the set of

all feasible path flow vectors, that is, vectors \mathbf{h} such that $\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs}$ for all OD pairs (r, s) and $h^\pi \geq 0$ for all $\pi \in \Pi$. Let the function $\mathbf{C}(\mathbf{h})$ represent the path travel times when the path flows are \mathbf{h} .² Let the multifunction $R(\mathbf{c})$ represent the set of path flows which could possibly occur if all travelers chose least cost paths given travel times \mathbf{c} . Mathematically

$$R(\mathbf{c}) = \left\{ \mathbf{h} \in H : h^\pi > 0 \text{ only if } c^\pi = \min_{\pi' \in \Pi^{rs}} c^{\pi'} \text{ for all } (r, s) \in Z^2, \pi \in \Pi^{rs} \right\} \quad (6.18)$$

This is the generalization of equation (6.17) when there are multiple OD pairs and multiple paths connecting each OD pair — be sure that you can see how (6.17) is a special case of (6.18) when there is just one OD pair connected by two paths.

Theorem 6.3. *If the link performance functions t_{ij} are continuous, then at least one equilibrium solution exists to the traffic assignment problem.*

Proof. Let the multifunction F be the composition of the multifunction R and the function \mathbf{C} defined above, so $F(\mathbf{h}) = R(\mathbf{C}(\mathbf{h}))$ is the set of path flow choices which are consistent with drivers choosing fastest paths when the travel times correspond to path flows \mathbf{h} . An equilibrium path flow vector is a fixed point of F : if $\mathbf{h} \in F(\mathbf{h})$ then the path flow choices and travel times are consistent with each other. The set H of feasible \mathbf{h} is convex and compact by Lemma 6.1 in the next section. Examining the definition of R in (6.18), we see that for any vector \mathbf{c} , $R(\mathbf{c})$ is nonempty (since there is at least one path of minimum cost), and compact by the same argument used in the proof of Lemma 6.1. Finally, the graph of F is closed, as you are asked to show in the exercises. \square

6.2 Properties of User Equilibrium

This section uses the mathematical formulations from the previous section to explore basic properties of user equilibrium solutions. In Section 6.2.1 we give fairly general conditions under which there is one (and only one) user equilibrium link flow solution in a network. Although this section is mostly mathematical, these properties are actually of great practical importance. We ultimately want to use traffic assignment as a tool to help evaluate and rank transportation projects, using the principle of user equilibrium to predict network conditions. If it is possible that there is no feasible assignment which satisfies the principle of user equilibrium, then we would be at a loss as to how projects should be ranked, and we would need to find another assignment rule. Or, if there could be multiple feasible assignments which satisfy the principle of user equilibrium, again it is unclear how projects should be ranked.

Interestingly, while the conditions for a unique *link flow* user equilibrium are fairly mild, there will almost certainly be multiple *path flow* solutions which satisfy the principle of equilibrium — infinitely many, in fact. As explained in

²If you look back to Chapter 5, you will see that $\mathbf{C}(\mathbf{h}) = \mathbf{\Delta}^T \mathbf{t}(\mathbf{\Delta} \mathbf{h})$ in matrix notation.

Section 6.2.2, this suggests that the principle of user equilibrium is not strong enough to identify path flows, simply link flows. If path flows are needed to evaluate a project, an alternative approach is needed, and this section explains the concepts of maximum entropy and proportionality which provide this alternative.

Lastly, when solving equilibrium problems on large networks both the link flow and path flow representations have significant limitations — algorithms only using link flows tend to be slow, while algorithms only using path flows can require a large amount of memory and tend to return low-entropy path flow solutions. A compromise is to use the link flows, but distinguish the flow on each link by its origin or destination. While explored more in Chapter 7, Section 6.2.3 lays the groundwork by showing how flow can be decomposed this way, and that at equilibrium the links with positive flow from each origin or destination form an acyclic network.

6.2.1 Existence and link flow uniqueness

We would like to use the mathematical formulations from the previous section — variational inequality, fixed point, and convex optimization — to identify conditions under which one (and only one) user equilibrium solution exists. This section relies heavily on mathematical formulations and results proved previously; interested readers may refer to those sections for more detail, and others can still read the proofs here to obtain the general idea. Recall the following results from Chapters 3 and 4:

- **Brouwer's Theorem 3.1:** Let X be a compact and convex set, and let $f : X \rightarrow X$ be a continuous function. Then there is at least one fixed point $x \in X$.
- **Proposition 4.6:** Let X be a convex set, and let $f : X \rightarrow \mathbb{R}$ be a strictly convex function. Then there is a unique global minimum \hat{x} of f on X .

Since all of these results rely on some properties of H and X , we establish them first:

Lemma 6.1. *The set of feasible path assignments H and the set of feasible link assignments X are both compact and convex.*

Proof. To show that H is compact, we must show that it is closed and bounded. The set H is defined by a combination of linear weak inequalities ($h^\pi \geq 0$) and linear equalities ($\sum h^\pi = d^{rs}$), so together Propositions 3.3b and 3.4a–b show that it is closed. For boundedness, consider any $\mathbf{h} \in H$, and OD pair (r, s) , and any path $\pi \in \Pi^{rs}$. We must have $h^\pi \leq d^{rs}$, since each h^π is nonnegative and $\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs}$. Therefore, if D is the largest entry in the OD matrix, we have $h^\pi \leq D$ for any path π . Therefore $|\mathbf{h}| \leq \sqrt{|\Pi|}D$, so H is contained in the ball $B_{\sqrt{|\Pi|}D}(\mathbf{0})$ and H is bounded.

Finally, for convexity, consider any $\mathbf{h}_1 \in H$, any $\mathbf{h}_2 \in H$, any $\lambda \in [0, 1]$, and the resulting vector $\mathbf{h} = \lambda\mathbf{h}_1 + (1 - \lambda)\mathbf{h}_2$. For any path π , $h_1^\pi \geq 0$ and $h_2^\pi \geq 0$,

so $h^\pi = \lambda h_1^\pi + (1 - \lambda)h_2^\pi \geq 0$ as well. Furthermore, for any OD pair (r, s) we have

$$\begin{aligned} \sum_{(r,s) \in \Pi^{rs}} h^\pi &= \sum_{(r,s) \in \Pi^{rs}} [\lambda h_1^\pi + (1 - \lambda)h_2^\pi] \\ &= \lambda \sum_{(r,s) \in \Pi^{rs}} h_1^\pi + (1 - \lambda) \sum_{(r,s) \in \Pi^{rs}} h_2^\pi \\ &= \lambda d^{rs} + (1 - \lambda)d^{rs} \\ &= d^{rs} \end{aligned}$$

so $\mathbf{h} \in H$ as well.

Every feasible link assignment $\mathbf{x} \in X$ is obtained from a linear transformation of some $\mathbf{h} \in H$ by (6.14) so X is also closed, bounded, and convex. \square

To show existence of solutions, we use Brouwer's Theorem based on the formulation of user equilibrium as a fixed point of the function $f(\mathbf{h}) = \text{proj}_H(\mathbf{h} - \mathbf{c}(\mathbf{h}))$.

Proposition 6.1. *If the link performance functions t_{ij} are continuous for each link $(i, j) \in A$, then there is a feasible assignment satisfying the principle of user equilibrium.*

Proof. Taking H as the set of feasible assignments, the range of the function $f(\mathbf{h}) = \text{proj}_H(\mathbf{h} - \mathbf{c}(\mathbf{h}))$ clearly lies in H because of the projection. By Lemma 6.1, H is compact and convex, so if we show that f is continuous, then Brouwer's Theorem guarantees existence of a fixed point. The discussion in Section 6.1 showed that each fixed point is a user equilibrium solution, which will be enough to prove the result.

We use the result that the composition of continuous functions is continuous (Proposition 3.6). The function f is the composition of two other functions (call them f_1 and f_2), where $f_1(\mathbf{h}) = \text{proj}_H(\mathbf{h})$ and $f_2(\mathbf{h}) = \mathbf{h} - \mathbf{c}(\mathbf{h})$. By Proposition 3.7, f_1 is continuous because H is a convex set. Furthermore, f_2 is continuous if \mathbf{c} is a continuous function of \mathbf{h} , which is true by hypothesis. Therefore, the conditions of Brouwer's Theorem are satisfied, and f has at least one fixed point (which satisfies the principle of user equilibrium). \square

The easiest way to show uniqueness of solutions is to make use of the convex optimization formulation, minimizing the function $f(\mathbf{x}) = \sum_{ij} \int_0^{x_{ij}} t_{ij}(x) dx$ over the set X .

Proposition 6.2. *If the link performance functions t_{ij} are differentiable, and $t'_{ij}(x) > 0$ for all x and links (i, j) , there is exactly one feasible link assignment satisfying the principle of user equilibrium.*

Proof. The previous section showed that user equilibrium link flow solutions correspond to minimum points of f on X . If we can show that f is strictly convex, then there can only be one minimum point. (Differentiability implies

continuity, so we can assume that at least one minimum point exists by Proposition 6.1). Since f is a function of multiple variables (each link's flow), to show that f is convex we can write its Hessian matrix of second partial derivatives. The first partial derivatives take the form $\frac{\partial f}{\partial x_{ij}} = t_{ij}(x_{ij})$. So, the diagonal entries of the Hessian take the form $\frac{\partial^2 f}{\partial x_{ij}^2} = t'_{ij}(x_{ij})$, while the off-diagonal entries take the form $\frac{\partial^2 f}{\partial x_{ij} \partial x_{kl}} = 0$. Since the Hessian is a diagonal matrix, and its diagonal entries are strictly positive by assumption, f is a strictly convex function. Therefore there is only one feasible link assignment satisfying the principle of user equilibrium. \square

It is also possible to prove the same result under the slightly weaker condition that the link performance functions are continuous and strictly increasing; this is undertaken in the exercises.

To summarize the above discussion, if the link performance functions are continuous, then at least one user equilibrium solution exists; if in addition the link performance functions are strictly increasing, then *exactly* one user equilibrium solution exists. For many problems representing automobile traffic, these conditions seem fairly reasonable, if not universally so: adding one more vehicle to the road is likely to increase the delay slightly, but not dramatically.

6.2.2 Path flow nonuniqueness, entropy, and proportionality

The previous subsection gave a relatively mild condition for the equilibrium solution to be unique in terms of link flows. However, a more subtle point is that the equilibrium solution need not be unique in terms of path flows. Consider the example shown in Figure 6.5. With the paths as numbered in the figure, the reader can verify that both $\mathbf{h}_1 = [20 \ 10 \ 20 \ 10]$ and $\mathbf{h}_2 = [30 \ 0 \ 10 \ 20]$ produce the equilibrium link flows, and therefore satisfy the principle of user equilibrium (all paths having equal travel time). Furthermore, any weighted average of \mathbf{h}_1 and \mathbf{h}_2 also produces the same link flows, so there are an infinite number of feasible path flow solutions which satisfy the principle of user equilibrium.

At first glance this may appear to contradict the proof of Proposition 6.2, which showed that the function f was strictly convex. Although f is strictly convex *as a function of \mathbf{x}* , it is *not* strictly convex as a function of \mathbf{h} . This is easy to see mathematically — if there are two distinct path flow solutions \mathbf{h}_1 and \mathbf{h}_2 which satisfy the principle of user equilibrium, then $\lambda \mathbf{h}_1 + (1 - \lambda) \mathbf{h}_2$ is an equilibrium as well for $\lambda \in (0, 1)$. Since all three are equilibria, $f(\mathbf{h}_1) = f(\mathbf{h}_2) = f(\lambda \mathbf{h}_1 + (1 - \lambda) \mathbf{h}_2)$ even though strict convexity would require $f(\lambda \mathbf{h}_1 + (1 - \lambda) \mathbf{h}_2) < \lambda f(\mathbf{h}_1) + (1 - \lambda) f(\mathbf{h}_2) = f(\mathbf{h}_1)$. From a more intuitive standpoint, in typical networks there are many more paths than links; therefore, it is very likely that multiple path flow solutions correspond to the same link flow solution. Since path travel times only depend on the values of the link flows (because link flows determine link travel times, which are added up to get

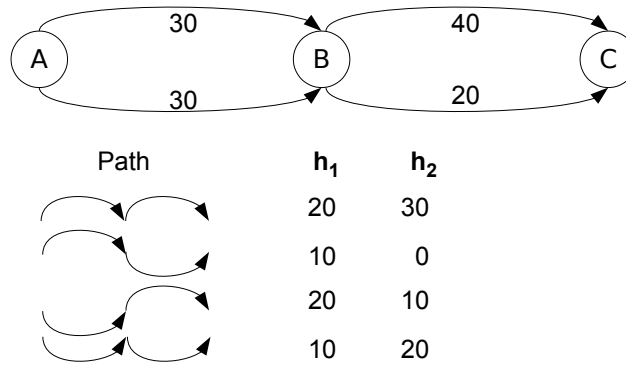


Figure 6.5: Nonuniqueness of equilibrium in terms of path flows

path travel times), and since the principle of user equilibrium is defined in terms of path travel times, it is reasonable to expect multiple path flow solutions to satisfy the equilibrium principle.

The existence of multiple path flow equilibria (despite a unique link flow equilibrium) is not simply a technical curiosity. Rather, it plays an important role in using network models to evaluate and rank alternatives. Thus far, we've been content to ask ourselves what the equilibrium link flows are, but in practice these numbers are usually used to generate other, more interesting measures of effectiveness such as the total system travel time and total vehicle-miles traveled. If a neighborhood group is worried about increased traffic, link flows can support this type of analysis. If we are concerned about safety, we can use link flows as one input in estimating crash frequency, and so forth.

Other important metrics, however, require more information than just the total number of vehicles on a link. Instead, we must know the entire *paths* used by drivers. Examples include

Select link analysis: In transportation planning parlance, a “select link analysis” goes beyond asking “how many vehicles are on a link,” to “exactly which vehicles are using this link?” This is used to produce nice visualizations, as well as to identify which origins and destinations are using a particular link. This is needed to identify which neighborhoods and regions are affected by transportation improvements, both positively and negatively.

Equity and environmental justice: As a special case of the above, planners are focusing more and more on the notion of equity. Most transportation projects involve both “winners” and “losers.” For instance, constructing a new freeway through an existing neighborhood may improve regional mobility, but cause significant harm to those whose neighborhood is disrupted in the process. Historically, disadvantaged populations tend

to bear the brunt of such projects. To know who is benefiting or suffering from a project, we need to know the origins and destinations of the travelers on different links.

Emissions: Vehicles emit more pollutants toward the start of trips, when the engine is cold, than later on once the catalytic converter has warmed up. Therefore, to use the output of a network model to predict emissions from vehicles on a link, ideally we need to distinguish between vehicles which have just started their trips from vehicles which have already traveled some distance. This requires knowing path flows.

In fact, software vendors often discuss such abilities as key features of their applications. In this light, the fact that equilibrium path flows are nonunique should trouble you somewhat. If there are multiple path flow equilibria (and therefore multiple values for emissions forecasts, equity analyses, etc.), which one is the “right” one? Simply picking one possible path flow solution and hoping that it’s the right one is not particularly rigorous, and inappropriate for serious planning applications.

So, the principle of user equilibrium is only strong enough to determine link flows, not path flows, even though having a path flow solution would be very helpful in analyzing engineering alternatives. This has motivated the search for the “most likely path flow” solution which satisfies equilibrium. Even though user equilibrium is not a strong enough concept to determine path flows uniquely, we can still ask “of all the path flow vectors which satisfy the principle of user equilibrium, which do we think is most likely to occur?” This section describes the concepts most commonly used to this end. To my knowledge, there has been little to no research directly validating this principle, but nobody has come up with a better, generally agreed-upon solution to this issue.

A physical analogy is presented to motivate the idea. Figure 6.6 shows the locations of n gas molecules within a box. In the left panel, all n molecules happen to be located in the top half of the box, while in the right panel $n/2$ molecules are located in the top half and $n/2$ in the bottom half. Both of these situations are physically possible (that is, they satisfy the laws of mechanics, etc.) although you would certainly think that the scenario on the right is more likely or plausible than the scenario on the left. In exactly the same way, although there are many path flow solutions which satisfy the “law” of user equilibrium, we are not forced to conclude that all such solutions are equally likely to occur.

Why does the scenario on the left seem so unlikely? If the molecules form an ideal gas, the location of each molecule is independent of the location of every other molecule. Therefore, the probability that any given molecule is in the top half of the box is $1/2$, and the probability that *all* n molecules are as in the left figure is

$$p^L = \left(\frac{1}{2}\right)^n. \quad (6.19)$$

To find the probability that the distribution of molecules is as in the right figure

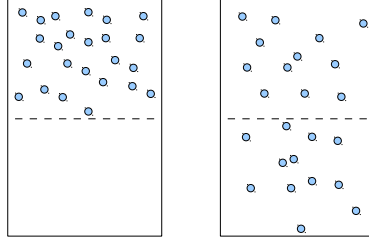


Figure 6.6: Molecules of an ideal gas in a box. Which situation is more likely?

(p^R) , we use the binomial distribution:

$$p^R = \binom{n}{n/2} \left(\frac{1}{2}\right)^{n/2} \left(\frac{1}{2}\right)^{n/2} = \frac{n!}{(n/2)!(n/2)!} \left(\frac{1}{2}\right)^n \quad (6.20)$$

Since $n! \gg [(n/2)!]^2$, we have $p^R \gg p^L$, that is, the situation on the right is much, much likelier than that on the left even though both are physically possible.

We want to use the same principle for path flows. The principle of user equilibrium assumes that each user is identical (in the sense that they choose routes only according to minimal travel time) and chooses routes independent of other drivers (except insofar as other drivers' choices affect travel times, drivers do not coordinate and strategize about their route choices with each other). Therefore, we can apply the same logic as with the gas molecules. For now, assume there is a single OD pair, with integer demand d , and let $\hat{\Pi} = \{\pi_1, \dots, \pi_k\}$ denote the set of minimum travel-time paths at the equilibrium solution. For simplicity assume that these paths have constant travel times independent of flow. Since drivers only care about travel time, each path is essentially identical, and the probability any particular driver chooses any particular path is $1/k$. The probability that the path flow vector \mathbf{h} takes a particular value $[h_1 \dots h_k]$ is then given by the multinomial distribution:

$$p = \frac{d!}{h_1!h_2! \dots h_k!} \left(\frac{1}{k}\right)^d \quad (6.21)$$

and the most likely path flow vector is the one which maximizes this product. Since $(1/k)^d$ is a constant, the most likely path flow vector simply maximizes

$$p = \frac{d!}{h_1!h_2! \dots h_k!} \quad (6.22)$$

Now we introduce a common trick: since the logarithm function is strictly increasing, the path flows which maximize p also maximize $\log p$, which is

$$\log p = \log d! - \sum_{\pi=1}^k \log h_{\pi}! \quad (6.23)$$

To simplify further, we use Stirling's approximation, which states that $n! \approx n^n e^{-n} \sqrt{2\pi n}$, or equivalently $\log n! \approx n \log n - n + (1/2) \log(2\pi n)$. This approximation is asymptotically exact in the sense that the ratio between $n!$ and $n^n e^{-n} \sqrt{2\pi n}$ approaches 1 as n grows large. Further, when n is large, n is much larger than $\log n$, so the last term can be safely ignored and $\log n! \approx n \log n - n$. Substituting into (6.23) we obtain

$$\log p \approx (d \log d - d) - \sum_{\pi} (h_{\pi} \log h_{\pi} - h_{\pi}) \quad (6.24)$$

Since $d = \sum_{\pi} h_{\pi}$, we can manipulate (6.24) to obtain

$$\log p \approx - \sum_{\pi} h_{\pi} \log(h_{\pi}/d) \quad (6.25)$$

and the path flows h_1, \dots, h_k maximizing this quantity (subject to the constraint $d = \sum_{\pi} h_{\pi}$ and nonnegativity) approximately maximize the probability that this particular path flow vector will occur in the field if travelers are choosing routes independent of each other.

To move towards the traffic assignment problem we're used to, we need to make the following changes:

1. In the traffic assignment problem, the demand d and path flows h_{π} do not need to be integers, but can instead take on real values. The corresponding interpretation is that the demand is "divided" into smaller and smaller units, each of which is assumed to act independently of every other unit. (This is how user equilibrium works with continuous flows.) This doesn't cause any problems, and in fact helps us out: as we take the limit towards infinite divisibility, Stirling's approximation becomes exact — so we can replace the \approx in our formulas with exact equality.
2. There are multiple origins and destinations. This doesn't change the basic idea, the formulas just become incrementally more complex as we sum equations of the form (6.25) for each OD pair.
3. Travel times are not constant, but are instead flow dependent. Again, this doesn't change any basic ideas; we just have to add a constraint stating that the path flows correspond to an equilibrium solution. Since the equilibrium link flows are unique, we simply have to state that the path flows correspond to the equilibrium link flows $\hat{\mathbf{x}}$.

Putting all this together, we seek the path flows \mathbf{h} which solve the optimization

problem

$$\max_{\mathbf{h}} \sum_{(r,s) \in Z^2} \sum_{\pi \in \hat{\Pi}^{rs}} h_{\pi} \log(h_{\pi}/d^{rs}) \quad (6.26)$$

$$\text{s.t.} \quad \sum_{\pi \in \Pi} \delta_{ij}^{\pi} h_{\pi} = \hat{x}_{ij} \quad \forall (i,j) \in A \quad (6.27)$$

$$\sum_{\pi \in \hat{\Pi}^{rs}} h_{\pi} = d^{rs} \quad \forall (r,s) \in Z^2 \quad (6.28)$$

$$h_{\pi} \geq 0 \quad \forall \pi \in \Pi \quad (6.29)$$

The objective function (6.26) is called the *entropy* of a particular path flow solution, and the constraints (6.28), (6.27), and (6.29) respectively require that the path flows are consistent with the OD matrix, equilibrium link flows, and nonnegativity. The set $\hat{\Pi}^{rs}$ is the set of paths which are used by OD pair (r,s) .

The word *entropy* is meant to be suggestive. The connection with the thermodynamic concept of entropy may be apparent from the physical analogy at the start of this section.³ In physics, entropy can be interpreted as a measure of disorder in a system. Both the left and right panels in Figure 6.6 are “allowable” in the sense that they obey the laws of physics. However, the scenario in the right has much less structure and much higher entropy, and is therefore more likely to occur.

It is not trivial to solve the optimization problem (6.26)–(6.29). It turns out that entropy maximization implies a much simpler condition, proportionality among pairs of alternate segments. Any two paths with a common origin and destination imply one or more pairs of alternate segments where the paths diverge, separated by links common to both paths. In Figure 6.5, there are two pairs of alternate segments: the top and bottom links between nodes A and B, and the top and bottom links between B and C. It turns out that path flows \mathbf{h}_1 are the solution to the entropy maximization problem.

You might notice some regularity or patterns in this solution. For instance, looking only at the top and bottom links between A and B, at equilibrium these links have equal flow (30 vehicles each). Paths 1 and 3 are the same except for between nodes A and B, and these paths also have equal flow (20 vehicles each). The same is true for paths 2 and 4 (10 vehicles each). Or, more subtly, between B and C, the ratio of flows between the top and bottom link is 2:1, the same as the ratio of flows on paths 1 and 2 (which only differ between B and C), and the ratio of flows on paths 3 and 4. This is no coincidence; in fact, we can show that *entropy maximization implies proportionality*.

Before proving this fact, we show that it holds even for different OD pairs. The network in Figure 6.7 has 40 vehicles traveling from origin A to destination F, and 120 vehicles from B to E, and the equilibrium link flows are shown in the figure. Each OD pair has two paths available to it, one using the top link between C and D, and the other using the bottom link between C and D. Using

³Actually, the term here is more directly drawn from the fascinating field of information theory, which is unfortunately beyond the scope of this text.

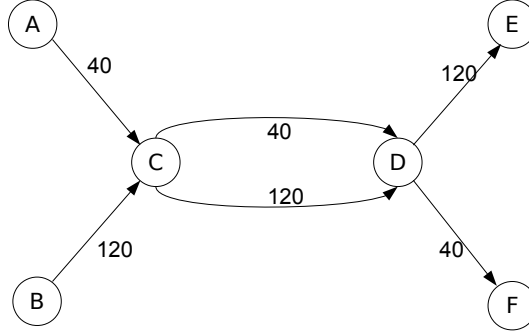


Figure 6.7: Multiple OD pair example of proportionality.

an upward-pointing arrow to denote the first type of path, and a downward-pointing arrow to describe the second, the four paths are h_{AF}^{\uparrow} , h_{AF}^{\downarrow} , h_{BE}^{\uparrow} , and h_{BE}^{\downarrow} . Solving the optimization problem (6.26)–(6.29), we find the most likely path flows are $h_{AF}^{\uparrow} = 10$, $h_{AF}^{\downarrow} = 30$, $h_{BE}^{\uparrow} = 30$, and $h_{BE}^{\downarrow} = 90$. The equilibrium flows for the top and bottom links between C and D have a ratio of 1:3; you can see that this ratio also holds between h_{AF}^{\uparrow} and h_{AF}^{\downarrow} , as well as between h_{BE}^{\uparrow} and h_{BE}^{\downarrow} .

In particular, the obvious-looking solution $h_{AF}^{\uparrow} = 40$, $h_{AF}^{\downarrow} = 0$, $h_{BE}^{\uparrow} = 0$, $h_{BE}^{\downarrow} = 120$ has extremely low entropy, because it implies that for some reason all travelers from one OD pair are taking one path, and all travelers from the other are taking the other path, even though both paths are perceived identically by all travelers and even though travelers are making choices independent of each other. This is exceptionally unlikely.

We now derive the proportionality condition by defining it a slightly more precise way:

Theorem 6.4. *Let π_1 and π_2 be any two paths connecting the same OD pair. If the path flows \mathbf{h} solve the entropy-maximizing problem (6.26)–(6.29), then the ratio of flows h_1/h_2 for paths π_1 and π_2 is identical regardless of the OD pair these paths connect, and only depends on the pairs of alternate segments distinguishing these two paths.*

The proof is a bit lengthy, and is deferred to the end of the section. But even though the derivation is somewhat involved, the proportionality condition itself is fairly intuitive: when confronted with the same set of choices, travelers from different OD pairs should behave in the same way. The proportionality condition implies nothing more than that the share of travelers choosing one alternative over another is the same across OD pairs. Proportionality is also a relatively easy condition to track and enforce.

It would be especially nice if proportionality can be shown fully equivalent to entropy maximization. Theorem 6.4 shows that entropy maximization im-

plies proportionality, but is the converse true? Unfortunately, the result is no, and examples can be created which satisfy proportionality without maximizing entropy. Therefore, proportionality is a weaker condition than entropy maximization. The good news is that proportionality “gets us most of the way there.” In the Chicago regional network, there are over 93 million equal travel-time paths at equilibrium; after accounting for the equilibrium and “no vehicle left behind” constraints (6.27) and (6.28), there are still over 90 million degrees of freedom in how the path flows are chosen. Accounting for proportionality reduces the number of degrees of freedom to 91 (a reduction of 99.9999%!) So, in practical terms, enforcing proportionality seems essentially equivalent to maximizing entropy.

6.2.3 Aggregation by origins or destinations

Thus far, we have looked at flow representations in one of two ways: path flows \mathbf{h} or link flows \mathbf{x} . In a way, these can be thought of as two extremes. The path flow solution contains the most information, showing the exact path chosen by every traveler in the network. The price paid for this level of detail is the amount of memory needed. The number of paths in a typical network is very, very large and storing the entire vector of path flows is impractical. The link flows, on the other hand, are much more compact and provide enough information for many common measures of effectiveness (such as vehicle-hours or vehicle-miles traveled, and volumes or delays on individual links). The link flow vector can be thought of as an “aggregated” form of the path flow vector, where we lump together all the travelers using the same link (regardless of which specific path, they are using, or which specific OD pair they are from). Naturally, some information is lost in this process of aggregation. So, there is a tradeoff: path flow solutions contain a great deal of information, at the expense of requiring a large amount of memory; link flow solutions are more compact but contain less information.

So, it is natural to ask if there is an intermediate “Goldilocks” flow representation which contains more detail than link flow solutions, without invoking the combinatorial explosion in storage space associated with a path flow solution. One such representation involves aggregating all travelers departing the same origin together.⁴ On a link (i, j) , let x_{ij}^r be the flow on this link who left origin r , that is,

$$x_{ij}^r = \sum_{s \in Z} \sum_{\pi \in \Pi^{rs}} \delta_{ij}^{\pi} h^{\pi} \quad (6.30)$$

and let \mathbf{x}^r be the vector of link flows associated with origin r .

Many of the most efficient algorithms for solving the traffic assignment problem use such a representation, often called an *origin-based* or *bush-based* representation. The term “origin-based” describes how the aggregation occurs. The term “bush-based” will make more sense once we get to Chapter 7. Let $\hat{\mathbf{h}}$ be

⁴Everything in this section would apply equally well to an aggregation by destination; the presentation from here on will be in terms of origins only to avoid repetition.

a feasible assignment satisfying the principle of user equilibrium, and let \mathbf{x}^r be the link flows associated with origin r at this solution.

It turns out that, at equilibrium, the links with positive x_{ij}^r values form an acyclic subnetwork. As a result, they have a topological order, which allows network calculations to be done much more rapidly than in general networks. This property is exploited heavily by bush-based traffic assignment algorithms, which are discussed in Section 7.4.

The following result shows that this must be true when each link's travel time is strictly positive. With slightly more effort, a similar result can be shown even when there are zero-time links.

Proposition 6.3. *Let \mathbf{x}^r be the link flows associated with origin r at some feasible assignment satisfying the principle of user equilibrium. If $t_{ij} > 0$ for all links (i, j) , then the subset of arcs $A^r = \{(i, j) \in A : x_{ij}^r > 0\}$ with positive flow from origin r contains no cycle.*

Proof. By contradiction, assume that there is a cycle $[i_1, i_2, \dots, i_k, i_1]$, where $(i_1, i_2), (i_2, i_3), \dots, (i_k, i_1)$ are all in A^r . Let L_i^r represent the travel time on the shortest path from origin r to node i . Consider any link $(i, j) \in A^r$; this implies that $x_{ij}^r > 0$ and that some used path starting from origin r includes (i, j) . Because the principle of user equilibrium holds, this path must be a shortest path to some destination. Clearly $L_i^r + t_{ij} \geq L_j^r$ by the definition of the labels \mathbf{L} ; but if $L_i^r + t_{ij} > L_j^r$, then there is a shorter path to node j than any of them passing through node i , and such a path cannot be used at equilibrium. Therefore we must have $L_i^r + t_{ij} = L_j^r$ for all links in A^r ; since all travel times are strictly positive this implies $L_i^r < L_j^r$ for all $(i, j) \in A^r$. So, for the cycle under consideration, we have $L_1 < L_2 < \dots < L_k < L_1$, which is impossible. Therefore the links in A^r cannot contain a cycle. \square

Proof of Theorem 6.4:

Begin by forming the Lagrangian for the entropy-maximizing problem:

$$\begin{aligned} \mathcal{L}(\mathbf{h}, \beta, \gamma) = & - \sum_{(r,s) \in Z^2} \sum_{\pi \in \hat{\Pi}^{rs}} h_\pi \log \left(\frac{h_\pi}{d^{rs}} \right) + \sum_{(i,j) \in A} \beta_{ij} \left(\hat{x}_{ij} - \sum_{\pi \in \Pi} \delta_{ij}^\pi h_\pi \right) \\ & + \sum_{(r,s) \in Z^2} \gamma_{rs} \left(d^{rs} - \sum_{\pi \in \hat{\Pi}^{rs}} h_\pi \right) \end{aligned} \quad (6.31)$$

using β and γ to denote the Lagrange multipliers. Note that the nonnegativity constraint can be effectively disregarded since the objective function is only defined for strictly positive \mathbf{h} . Therefore at the optimum solution the partial derivative of \mathcal{L} with respect to any path flow must vanish:

$$\frac{\partial \mathcal{L}}{\partial h_\pi} = -1 - \log \left(\frac{h_\pi}{d^{rs}} \right) - \sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij} - \gamma_{rs} = 0 \quad (6.32)$$

where (r, s) is the OD pair connected by path π . Solving for h^π we obtain

$$h^\pi = d^{rs} \exp \left(-1 - \sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij} - \gamma_{rs} \right) \quad (6.33)$$

Likewise we have

$$\frac{\partial \mathcal{L}}{\partial \gamma_{rs}} = d^{rs} - \sum_{\pi' \in \hat{\Pi}^{rs}} h_{\pi'} = 0 \quad (6.34)$$

so

$$d^{rs} = \sum_{\pi' \in \hat{\Pi}^{rs}} h_{\pi'} = d^{rs} \exp(-1 - \gamma_{rs}) \sum_{\pi \in \hat{\Pi}^{rs}} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij} \right) \quad (6.35)$$

substituting the result from (6.33). Therefore we can solve for γ_{rs} :

$$\gamma_{rs} = -1 + \log \left(\sum_{\pi' \in \hat{\Pi}^{rs}} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^{\pi'} \beta_{ij} \right) \right) \quad (6.36)$$

Finally, substituting (6.36) into (6.33) and simplifying, we obtain

$$h^\pi = \frac{d^{rs}}{\sum_{\pi' \in \hat{\Pi}^{rs}} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^{\pi'} \beta_{ij} \right)} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij} \right) \quad (6.37)$$

Noting that the fraction in (6.37) only depends on the OD pair (r, s) , we can simply write

$$h^\pi = K_{rs} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij} \right) \quad (6.38)$$

where K_{rs} is a constant associated with OD pair (r, s) .

Let A_1 be the set of links not in either π_1 or π_2 , A_2 the set of links common to both paths, A_3 the links in π_1 but not π_2 , and A_4 the links in π_2 but not π_1 .

Then the ratio h_1/h_2 can be written

$$\frac{h_1}{h_2} = \frac{K_{rs} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^{\pi_1} \beta_{ij} \right)}{K_{rs} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^{\pi_2} \beta_{ij} \right)} \quad (6.39)$$

$$= \frac{K_{rs} \exp \left(\begin{array}{c} - \sum_{(i,j) \in A_1} \delta_{ij}^{\pi_1} \beta_{ij} - \sum_{(i,j) \in A_2} \delta_{ij}^{\pi_1} \beta_{ij} \\ - \sum_{(i,j) \in A_3} \delta_{ij}^{\pi_1} \beta_{ij} - \sum_{(i,j) \in A_4} \delta_{ij}^{\pi_1} \beta_{ij} \end{array} \right)}{K_{rs} \exp \left(\begin{array}{c} - \sum_{(i,j) \in A_1} \delta_{ij}^{\pi_2} \beta_{ij} - \sum_{(i,j) \in A_2} \delta_{ij}^{\pi_2} \beta_{ij} \\ - \sum_{(i,j) \in A_3} \delta_{ij}^{\pi_2} \beta_{ij} - \sum_{(i,j) \in A_4} \delta_{ij}^{\pi_2} \beta_{ij} \end{array} \right)} \quad (6.40)$$

$$= \frac{K_{rs} \exp \left(- \sum_{(i,j) \in A_2} \delta_{ij}^{\pi_1} \beta_{ij} - \sum_{(i,j) \in A_3} \delta_{ij}^{\pi_1} \beta_{ij} \right)}{K_{rs} \exp \left(- \sum_{(i,j) \in A_2} \delta_{ij}^{\pi_2} \beta_{ij} - \sum_{(i,j) \in A_4} \delta_{ij}^{\pi_2} \beta_{ij} \right)} \quad (6.41)$$

$$= \frac{K_{rs} \exp \left(- \sum_{(i,j) \in A_2} \delta_{ij}^{\pi_1} \beta_{ij} \right) \exp \left(- \sum_{(i,j) \in A_3} \delta_{ij}^{\pi_1} \beta_{ij} \right)}{K_{rs} \exp \left(- \sum_{(i,j) \in A_2} \delta_{ij}^{\pi_2} \beta_{ij} \right) \exp \left(- \sum_{(i,j) \in A_4} \delta_{ij}^{\pi_2} \beta_{ij} \right)} \quad (6.42)$$

$$= \frac{\exp \left(- \sum_{(i,j) \in A_3} \delta_{ij}^{\pi_1} \beta_{ij} \right)}{\exp \left(- \sum_{(i,j) \in A_4} \delta_{ij}^{\pi_2} \beta_{ij} \right)} \quad (6.43)$$

where the steps of the derivation respectively involve substituting (6.37), expanding A into $A_1 \cup A_2 \cup A_3 \cup A_4$, using the definitions of the A_i sets to identify sums where $\delta_{ij}^{\pi_i} = 0$, splitting exponential terms, and canceling common factors.

Thus in the end we have

$$\frac{h_1}{h_2} = \frac{\exp \left(- \sum_{(i,j) \in A_3} \delta_{ij}^{\pi_1} \beta_{ij} \right)}{\exp \left(- \sum_{(i,j) \in A_4} \delta_{ij}^{\pi_2} \beta_{ij} \right)} \quad (6.44)$$

regardless of the OD pair h_1 and h_2 connect, and this ratio only depends on A_3 and A_4 , that is, the pairs of alternate segments distinguishing π_1 and π_2 .

6.3 Alternative Assignment Rules

The previous section described the most important assignment rule, the principle of user equilibrium. However, it is not the only possible assignment rule, and it is not hard to see that other factors may enter route choice beyond travel time, or that other assumptions of the user equilibrium model are not perfectly realistic. Over time researchers have introduced a number of alternative assignment rules that can be used. This text by necessity can only cover a few of these alternative assignment rules, and three representative rules were chosen: assignment with perception errors, system optimal assignment, and bounded rationality.

Assignment with perception errors relaxes the assumption that drivers have perfect knowledge of travel times. The interpretation is that every driver *perceives* a particular travel time on each path (which may or may not equal its actual travel time), and chooses a path that they believe to be shortest based on these perceptions. The net effect is that travelers are distributed across paths with different travel times, but are concentrated more on paths with lower travel times. (If a path is very much longer than the shortest one, it is unlikely that someone's perception would be so far off as to mistakenly think it is shortest.)

The mathematical tool most commonly used to model perception errors is the theory of *discrete choice* from economics. The resulting model is termed *stochastic user equilibrium*, and this important model is the focus of Section 9.3.

The other two alternative assignment rules are discussed next.

6.3.1 System optimal assignment

Imagine for a moment that route choice was taken out of the hands of individual travelers, and instead placed in the hands of a dictator who could assign each traveler a path that they would be required to follow. Suppose also that this dictator was benevolent, and wanted to act in a way to minimize average travel delay in the network. The resulting assignment rule results in the *system optimal* state.

While this scenario is a bit fanciful, the system optimal state is important for several reasons. First, it provides a theoretical lower bound on the delay in a network, a benchmark for the best performance that could conceivably be achieved that more realistic assignment rules can be compared to. Second, there are ways to actually achieve the system optimal state even without taking route choice out of the hands of travelers, if one has the ability to freely charge tolls or provide incentives on network links. Third, there are some network problems where a single agent can exert control over the routes chosen by travelers, as in certain logistics problems where a dispatcher can assign specific routes to vehicles.

At first, it may not be obvious that this assignment rule is meaningfully different from the user equilibrium assignment rule. After all, in user equilibrium each traveler is individually choosing routes, while in system optimum a single agent is choosing routes for everyone, but both are doing so to minimize travel times. To see why they might be different, consider the network used for the Knight-Pigou-Downs paradox in Section 1.5: a two-link network (Figure 5.6) with a constant travel time of 50 minutes on the top link, a link performance function $45 + x^\downarrow$ on the bottom link, and a total demand of 30 vehicles. The user equilibrium solution was $x^\uparrow = 25$, $x^\downarrow = 5$, with equal travel times of 50 minutes on both top and bottom routes.

So, in the user equilibrium state the average travel time is 50 minutes for all

vehicles. Is it possible to do better? We can write the average travel time as

$$\begin{aligned} \frac{1}{50} (50x^\uparrow + (45 + x^\downarrow)x^\downarrow) &= \frac{1}{50} (50(30 - x^\downarrow) + (45 + x^\downarrow)x^\downarrow) \\ &= \frac{1}{50} ((x^\downarrow)^2 - 5x^\downarrow + 1500) . \end{aligned} \quad (6.45)$$

This is a quadratic function which obtains its minimum at $x^\downarrow = 2.5$. At the solution $x^\uparrow = 27.5$, $x^\downarrow = 2.5$, the travel times are unequal ($t^\uparrow = 50$, $t^\downarrow = 47.5$), but the average travel time of 49.8 minutes is slightly less than the average travel time of 50 minutes at the user equilibrium solution.

Therefore, travelers individually choosing routes to minimize their own travel times may create more delay than would be obtained with central control. This reveals an important, but subtle point about user equilibrium assignment, a point important enough that Section 6.4 is devoted entirely to explaining this issue. So, we will not belabor the point here, but it will be instructive to start thinking about why the user equilibrium and system optimal states need not coincide.

To be more precise mathematically, the system optimal assignment rule chooses the feasible assignment minimizing the average travel time. Since the total number of travelers is a constant, we can just as well minimize the *total system travel time*, defined as

$$TSTT = \sum_{\pi \in \Pi} h^\pi c^\pi = \sum_{(i,j) \in A} x_{ij} t_{ij} \quad (6.46)$$

where it is not hard to show that calculating $TSTT$ through path flows and link flows produces the same value.

6.3.2 Bounded rationality (*)

(This optional section discusses an assignment rule which relaxes the assumption that all travelers are on the shortest path.)

An alternative assignment rule is based on another variation of shortest path-seeking behavior. Under *bounded rationality*, drivers will choose paths that are within some threshold of the true shortest path travel time. The behavioral rationale is that a driver will switch paths from a longer path to a shorter one, but only if the new path is sufficiently faster than their current choice. Perhaps a driver is unlikely to switch paths to save, say, a tenth of a second in travel time: the burden of learning a new path outweighs any travel time savings which are possible.

There is some evidence from behavioral economics that humans often make choices in a boundedly rational way. Specifically in network assignment, bounded rationality can explain observed reactions to changes in a network in a way that the principle of user equilibrium cannot. Of course, these advantages must be balanced against other concerns: there may be many distinct bounded rational

assignments, and efficient methods for finding bounded rational assignments in large networks have not yet been developed.

The boundedly rational user equilibrium (BRUE) problem can be stated as follows. Let ϵ represent some tolerance value in terms of travel time, and let κ^{rs} be the travel time on the shortest path connecting origin r to destination s .

Definition 6.1. (Principle of boundedly rational user equilibrium [BRUE].) *Every used path has travel time no more than ϵ in excess of the travel time on the shortest path connecting these nodes. That is, the feasible assignment $\mathbf{h} \in H$ is a BRUE if, for each $(r, s) \in Z^2$ and $\pi \in \Pi^{rs}$, $h^\pi > 0$ implies $c^\pi \leq \kappa^{rs} + \epsilon$.*

As an example, consider once again the network in Figure 6.3, but now assume we are seeking a BRUE where the tolerance is $\epsilon = 2$ minutes. That is, travelers are willing to use any path as long as it is within one minute of being the shortest path. The user equilibrium solution $x^\uparrow = 25$, $x^\downarrow = 5$ is clearly a BRUE as well. However, there are others: if $x^\uparrow = 26$, $x^\downarrow = 4$, then the travel times on the top and bottom paths are 50 and 49 minutes, respectively. Even though the majority of travelers (those on the top link) are not on the shortest path, the travel time on their path is close enough to the shortest path that the situation is acceptable, and this solution is BRUE as well. However, the flows $x^\uparrow = 28$, $x^\downarrow = 2$ are not BRUE: the travel times are 50 and 47 minutes, and the difference between the (used) top path and the shortest path (3 minutes) exceeds the tolerance value (2 minutes). You should convince yourself that $x^\uparrow \in [23, 27]$ and $x^\downarrow = 30 - x^\uparrow$ characterize all of the BRUE solutions in this simple network.

To come up with a mathematical formulation of BRUE assignment, it is helpful to introduce an auxiliary variable ρ^π indicating the amount by which the travel time on path π falls short of the maximum acceptable travel time on a path. Specifically, define

$$\rho^\pi = [\kappa^{rs} + \epsilon - c^\pi]^+ \quad \forall (r, s) \in Z^2, \pi \in \Pi^{rs} \quad (6.47)$$

So, in the example of the previous paragraph, if $\epsilon = 1$, $x^\uparrow = 26$ and $x^\downarrow = 4$, then $t^\uparrow = 50$, $t^\downarrow = 49$, and thus $\rho^\uparrow = [49 + 2 - 50]^+ = 1$ and $\rho^\downarrow = [49 + 2 - 49]^+ = 2$. If the path π is unacceptable ($c^\pi > \kappa^{rs} + \epsilon$), then $\rho^\pi = 0$ by (6.47). We can then define the *effective travel time* \hat{c}^π as

$$\hat{c}^\pi = c^\pi + \rho^\pi. \quad (6.48)$$

These new quantities are helpful because the principle of BRUE can now be formulated in a more familiar way:

Proposition 6.4. *A feasible assignment $\mathbf{h} \in H$ is a BRUE if and only if every used path has equal and minimal effective travel time, that is, for each $(r, s) \in Z^2$ and $\pi \in \Pi^{rs}$, $h^\pi > 0$ implies $\hat{c}^\pi = \min_{\pi' \in \Pi^{rs}} \hat{c}^{\pi'}$.*

Proof. Let $\mathbf{h} \in H$ satisfy the principle of BRUE, and consider any OD pair (r, s) and path $\pi \in \Pi^{rs}$ with $h^\pi > 0$. Since \mathbf{h} is BRUE, $c^\pi \leq \kappa^{rs} + \epsilon$, so $\rho^\pi = \kappa^{rs} + \epsilon - c^\pi$ and $\hat{c}^\pi = c^\pi + \kappa^{rs} + \epsilon - c^\pi = \kappa^{rs} + \epsilon$. But for any $\pi' \in \Pi^{rs}$, $\hat{c}^{\pi'} = c^{\pi'} + \rho^{\pi'} \geq c^{\pi'} + \kappa^{rs} + \epsilon - c^{\pi'} = \kappa^{rs} + \epsilon$, so $\hat{c}^\pi = \min_{\pi' \in \Pi^{rs}} \hat{c}^{\pi'}$.

Contrariwise, let $\mathbf{h} \in H$ be such that $h^\pi > 0$ implies $\hat{c}^\pi = \min_{\pi' \in \Pi^{rs}} \hat{c}^{\pi'}$, and consider any OD pair (r, s) . Let $\hat{\pi}$ be a shortest path connecting r to s , so $\hat{c}^{\hat{\pi}} = \kappa^{rs} + [\kappa^{rs} + \epsilon - \kappa^{rs}]^+ = \kappa^{rs} + \epsilon$ and this must be minimal among all effective travel times on paths connecting r to s . For any path $\pi \in \Pi^{rs}$ with $h^\pi > 0$, we thus have $\hat{c}^\pi = \kappa^{rs} + \epsilon$, so $c^\pi = \hat{c}^\pi - \rho^\pi = \kappa^{rs} + \epsilon - \rho^\pi \leq \kappa^{rs} + \epsilon$ since $\rho^\pi \geq 0$. Thus \mathbf{h} satisfies the principle of BRUE as well. \square

Proposition 6.4 may give you hope that we can write a convex optimization problem whose solutions correspond to BRUE. Unfortunately, this is not the case, for reasons that will be shown shortly. However, we can write a variational inequality where both the path flows \mathbf{h} and the auxiliary variables $\boldsymbol{\rho}$ are decision variables. In what follows, let $\kappa^{rs}(\mathbf{h}) = \min_{\pi \in \Pi^{rs}} c^\pi(\mathbf{h})$ and, with slight abuse of notation, let κ^π refer to the κ^{rs} value corresponding to the path π , and $\boldsymbol{\kappa}$ be the vector of κ^π values.

Proposition 6.5. *A vector of path flows $\hat{\mathbf{h}} \in H$ is a BRUE if there exists a vector of nonnegative auxiliary costs $\hat{\boldsymbol{\rho}} \in \mathbb{R}_+^{|\Pi|}$ such that*

$$\begin{bmatrix} \hat{\mathbf{c}}(\hat{\mathbf{h}}, \hat{\boldsymbol{\rho}}) \\ \mathbf{c}(\hat{\mathbf{h}}) + \hat{\boldsymbol{\rho}} - \boldsymbol{\kappa}(\hat{\mathbf{h}}) - \epsilon \end{bmatrix} \cdot \begin{bmatrix} \hat{\mathbf{h}} - \mathbf{h} \\ \hat{\boldsymbol{\rho}} - \boldsymbol{\rho} \end{bmatrix} \leq 0 \quad (6.49)$$

for all $\mathbf{h} \in H$, $\boldsymbol{\rho} \in \mathbb{R}_+^{|\Pi|}$.

Proof. Assume that $(\hat{\mathbf{h}}, \hat{\boldsymbol{\rho}})$ solve the variational inequality (6.49). The first component of the VI $\hat{\mathbf{c}}(\hat{\mathbf{h}}, \hat{\boldsymbol{\rho}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) \leq 0$ shows that all used paths have equal and minimal effective cost, and by Proposition 6.4 satisfy the principle of BRUE, assuming that $\boldsymbol{\rho}$ is consistent with (6.47). To show this, consider the second component of the VI: $(\boldsymbol{\kappa}(\hat{\mathbf{h}}) + \epsilon - \hat{\boldsymbol{\rho}}) \cdot \hat{\boldsymbol{\rho}} - \boldsymbol{\rho} \leq 0$, or equivalently

$$\sum_{\pi \in \Pi} (c^\pi(\hat{\mathbf{h}}) + \hat{\rho}_\pi - \kappa^\pi(\hat{\mathbf{h}}) - \epsilon) \hat{\rho}_\pi \leq \sum_{\pi \in \Pi} (c^\pi(\hat{\mathbf{h}}) + \hat{\rho}_\pi - \kappa^\pi(\hat{\mathbf{h}}) - \epsilon) \rho_\pi \quad (6.50)$$

for all $\boldsymbol{\rho} \in \mathbb{R}_+^{|\Pi|}$. If $\hat{\rho}_\pi > 0$ for any path, inequality (6.50) can be true only if $\rho^\pi = \kappa^\pi(\hat{\mathbf{h}}) + \epsilon - c^\pi(\hat{\mathbf{h}})$; or, if $\hat{\rho}_\pi = 0$, then inequality (6.50) can be true only if $\rho^\pi = 0 \leq \kappa^\pi(\hat{\mathbf{h}}) + \epsilon - c^\pi(\hat{\mathbf{h}})$. In either case (6.47) is satisfied.

In the reverse direction, assume that $\hat{\mathbf{h}}$ is BRUE. Choose $\hat{\boldsymbol{\rho}} \in \mathbb{R}_+^{|\Pi|}$ according to (6.47). By Proposition 6.4 surely $\hat{\mathbf{c}}(\hat{\mathbf{h}}, \hat{\boldsymbol{\rho}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) \leq 0$ for all $\mathbf{h} \in H$. For any path where $\rho^\pi = \kappa^\pi(\hat{\mathbf{h}}) + \epsilon - c^\pi(\hat{\mathbf{h}})$, clearly $\sum_{\pi \in \Pi} (c^\pi(\hat{\mathbf{h}}) + \hat{\rho}_\pi - \kappa^\pi(\hat{\mathbf{h}}) - \epsilon)(\hat{\rho}_\pi - \rho^\pi) = 0$; and for any path where $\rho^\pi = 0$, $\kappa^\pi(\hat{\mathbf{h}}) + \epsilon - c^\pi(\hat{\mathbf{h}}) \geq 0$ and $\sum_{\pi \in \Pi} (c^\pi(\hat{\mathbf{h}}) + \hat{\rho}_\pi - \kappa^\pi(\hat{\mathbf{h}}) - \epsilon)(\hat{\rho}_\pi - \rho^\pi) \geq 0$. In either case the second component of the VI is satisfied as well. \square

This BRUE formulation has been helpful in explaining observed changes in network flows after a disruption. For instance, assume that a link is removed from the network due to a disaster of some sort, and that flows adjust towards a new equilibrium in the network without the affected link. When the link is

restored, flows will adjust again. If the principle of user equilibrium is true, the flows will move back to exactly the same values as before. However, in practice there has been some “stickiness” observed, and not all drivers will return to the same routes they were initially on. The BRUE framework provides a logical explanation for this: when the network is disrupted, certain drivers were forced to choose new paths. When the network is restored, they will only switch back to their original paths if the travel time savings are sufficiently large. Otherwise, they will remain on their new paths.

However, an implication of this finding is that the BRUE solution need not be unique, even in link flows. (This was also seen in the small example at the start of this subsection.) In general, the set of BRUE solutions is not convex either — which immediately shows that there is no convex optimization problem whose solutions correspond to BRUE assignments, since the sets of minima to convex optimization problems are always convex. However, there is always at least one BRUE solution by Proposition 6.1, since any path flows satisfying the principle of user equilibrium also satisfy BRUE.

Still, the lack of uniqueness means that BRUE should be applied carefully. While it is certainly more realistic than the principle of user equilibrium (since user equilibrium can always be obtained as a special case with $\epsilon = 0$), it is considerably more difficult to use BRUE to rank projects (which BRUE assignment should be used?) and at present it is not clear what value of ϵ best represents travel choices. Whether the added realism offsets these disadvantages is application and network dependent.

6.4 User Equilibrium and System Optimal

The relationship between the user equilibrium and system optimal assignment rules is worth studying further, for a few reasons. First, the fact that these rules lead to distinct solutions means that travelers cannot be relied upon to independently behave in a way that minimizes congestion⁵, even though each driver is minimizing his or her own travel time. This means that there is a role for transportation professionals to reduce congestion through planning and operational control. Exploring why this happens leads to a discussion of economic externalities, which are illustrated in the Braess paradox discussed in Section 5.3.

Second, the mathematical structure of the system optimal assignment problem is actually very similar to that of the user equilibrium problem. In fact, if one has a “black box” procedure for solving user equilibrium assignment, one can just as easily solve system optimal assignment using that same procedure, making some easy modifications to the network. Likewise, if one has a procedure for solving system optimal assignment, one can just as easily solve user equilibrium by modifying the network. So, when we discuss algorithms for solving the user equilibrium traffic assignment problem in Chapter 7, they all apply

⁵Perhaps not a very surprising observation to those used to driving in crowded cities.

equally well to the system optimal problem. This relationship is explored in Section 6.4.2.

Finally, there is a more recent and interesting set of results which attempts to quantify how far apart the user equilibrium and system optimal assignments could possibly be, in terms of total system travel time. At first glance it seems that these assignment rules would tend to give relatively similar total travel times, since both are based on travel time minimization (one individual and the other collective). However, the Braess paradox shows that these can diverge in surprising ways. Nevertheless, in many cases one can in fact bound how much inefficiency is caused by allowing individual drivers to pick their own routes, a set of results called the *price of anarchy* and discussed in Section 6.4.3.

6.4.1 Externalities

Recall the three examples in Section 5.3, in which equilibrium traffic assignment exhibited counterintuitive results. Perhaps the most striking of these was the Braess paradox, in which adding a new roadway link actually worsened travel times for all travelers. This is counterintuitive, because if the travelers simply agreed to stay off of the newly-built road, their travel times would have been no worse than before. This section describes exactly why this happens. But before delving into the roots of the Braess paradox, let's discuss some of the implications:

User equilibrium does *not* minimize congestion. In principle, everyone could have stayed on their original paths and retained the lower travel time of 83; the only catch was that this was no longer an equilibrium solution with the new link, and the equilibrium solution with the link is worse. Therefore, there may be traffic assignments with lower travel times (for everybody!) than the user equilibrium one.

The “invisible hand” may not work in traffic networks. In economics, Adam Smith's celebrated “invisible hand” suggests that each individual acting in his or her own self-interest also tends to maximize the interests of society as well. The Braess paradox shows that this is not necessarily true in transportation networks. The two drivers switched paths out of self-interest, to reduce their own travel times; but the end effect made things worse for everyone, including the drivers who switched paths.

There may be room to “improve” route choices. A corollary of the previous point, when the invisible hand does not function well, a case can be made for regulation or another coordination mechanism to make everyone better off. In this case, this mechanism might take the form of certain transportation policies. What might they be?

If you have studied economics, you might already have some idea of why the Braess paradox occurs. The “invisible hand” requires certain things to be true if individual self-interest is to align with societal interest, and fails otherwise.

A common case where the invisible hand fails is in the case of *externalities*: an externality is a cost or benefit resulting from a decision made by one person, and felt by another person who had no say in the first person's decision. Examples of externalities include industrial pollution (we cannot rely on industries to self-regulate pollution, because all citizens near a factory suffer the effects of poor air quality, even though they had no say in the matter), driving without a seatbelt (if you get into an accident without a seatbelt, it is likely to be more serious, costing others who help pay your health care costs and delaying traffic longer while the ambulance arrives; those who help pay your costs and sit in traffic had no control over your decision to not wear a seatbelt), and education (educated people tend to commit less crime, which benefits all of society, even those who do not receive an education). The first two examples involve costs, and so are called negative externalities; the latter involves a benefit, and is called a positive externality.

Another example of an externality is seen in the prisoner's dilemma (the Erica-Fred game of Section 1.3). When Erica chooses to testify against Fred, she does so because it reduces her jail time by one year. The fact that her choice also increases Fred's jail time by fourteen years was irrelevant (such is the nature of greedy decision-making). When *both* Erica and Fred behaved in this way, the net effect was to dramatically increase the jail time they experienced, even though each of them made the choices they did in an attempt to minimize their jail time.

The relevance of this economics tangent is that *congestion* can be thought of as an externality. Let's say I'm driving during the peak period, and I can choose an uncongested back road which is out of my way, or a congested freeway which is a direct route (and thus faster, even with the congestion). If I choose the freeway, I end up delaying everyone with the bad luck of being behind me. Perhaps not by much; maybe my presence increases each person's commute by a few seconds. However, multiply these few seconds by a large number of drivers sitting in traffic, and we find that my presence has cost society as a whole a substantial amount of wasted time. This is an externality, because those other drivers had no say over my decision to save a minute or two to take the freeway. (This is why the user equilibrium assumption is said to be "greedy." It assumes a driver will happily switch routes to save a minute even if it increases the total time people spend waiting in traffic by an hour.) These concepts are made more precise in the next subsection.

As a society, we tend to adopt regulation or other mechanisms for minimizing the impact of negative externalities, such as emissions regulation or seatbelt laws (although you can probably think of many negative externalities which still exist). Transfer payments are another option: if a factory had to directly compensate each citizen for the harm they suffer as a result of pollution, the externality no longer exists because the factory now has to account for the costs it imposed on society, and will only continue to operate if its profits exceed the cost of the pollution it causes. What types of regulation or transfer payments might exist in the transportation world?

6.4.2 Mathematical equivalence

Recall that the user equilibrium traffic assignment was the solution to the following convex optimization problem:

$$\min_{\mathbf{x}, \mathbf{h}} \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (6.51)$$

$$\text{s.t. } x_{ij} = \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \quad \forall (i, j) \in A \quad (6.52)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (6.53)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (6.54)$$

while the system optimal traffic assignment solves this convex optimization problem:

$$\min_{\mathbf{x}, \mathbf{h}} \sum_{(i,j) \in A} t_{ij}(x_{ij}) x_{ij} \quad (6.55)$$

$$\text{s.t. } x_{ij} = \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \quad \forall (i, j) \in A \quad (6.56)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (6.57)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (6.58)$$

Inspecting these, the only difference is in their objective functions, and even these are quite similar: both are sums of terms involving each link in the network. The problems share essentially the same structure, the only difference is the term corresponding to each link. For the user equilibrium problem, the term for link (i, j) is $\int_0^{x_{ij}} t_{ij}(x) dx$, while for the system optimum problem it is $t_{ij}(x_{ij})x_{ij}$.

Suppose for a moment that we had some algorithm which can solve the user equilibrium traffic assignment problem for any input network and OD matrix. (A number of these algorithms will be discussed in Chapter 7.) Now suppose that we had to solve the system optimal problem. If we were to replace the link performance functions $t_{ij}(x_{ij})$ with modified functions $\hat{t}_{ij}(x_{ij})$ defined by

$$\hat{t}_{ij}(x_{ij}) = t_{ij}(x_{ij}) + t'_{ij}(x_{ij})x_{ij} \quad (6.59)$$

where t'_{ij} is the derivative of the link performance functions, then

$$\int_0^{x_{ij}} \hat{t}_{ij}(x) dx = \int_0^{x_{ij}} (t_{ij}(x) + t'_{ij}(x)x) dx \quad (6.60)$$

Integrating the second term by parts, this simplifies to

$$\int_0^{x_{ij}} \hat{t}_{ij}(x) dx = t_{ij}(x_{ij})x_{ij} \quad (6.61)$$

which is the same term for link (i, j) used in the system optimal function.

That is, solving user equilibrium with the modified functions \hat{t}_{ij} produces the same link and path flow solution as solving system optimum with the original performance functions t_{ij} . The interpretation of the formula (6.59) is closely linked with the concept of externalities introduced in the previous subsection. Equation (6.59) consists of two parts: the actual travel time $t_{ij}(x_{ij})$, and an additional term $t'_{ij}(x_{ij})x_{ij}$. If an additional vehicle were to travel on link (i, j) , it would experience a travel time of $t_{ij}(x_{ij})$. At the margin, the travel time on link (i, j) would increase by approximately $t'_{ij}(x_{ij})$, and this marginal increase in the travel time would be felt by the other x_{ij} vehicles on the link. So, the second term in (6.59) expresses the additional, external increase in travel time caused by a vehicle using link (i, j) . For this reason, the modified functions \hat{t}_{ij} can be said to represent the *marginal cost* on a link.

Since the system optimal solution is a “user equilibrium” with respect to the modified costs \hat{t}_{ij} , we can formulate a “principle of system optimum” similar to the principle of user equilibrium:

Definition 6.2. (Principle of system optimum.) *Every used route connecting an origin and destination has equal and minimal marginal cost.*

So, system optimum can be seen as a special case of user equilibrium, with modified link performance functions. The converse is true as well. Suppose we had an algorithm which would solve the system optimal problem for any input network and OD matrix. If we replace the link performance functions $t_{ij}(x_{ij})$ with modified functions $\tilde{t}_{ij}(x_{ij})$ defined by

$$\tilde{t}_{ij}(x_{ij}) = \frac{\int_0^{x_{ij}} t_{ij}(x) dx}{x_{ij}} \quad (6.62)$$

when $x_{ij} > 0$ and $\tilde{t}_{ij}(x_{ij}) = 0$ otherwise, a similar argument shows that the objective function for the system optimal problem with link performance functions \tilde{t}_{ij} is identical to that for the user equilibrium problem with link performance functions t_{ij} .

As a result, despite very different physical interpretations, the user equilibrium and system optimal problems are essentially the same mathematically. If you can solve one, you can solve the other by making some simple changes to the link performance functions. So, there is no need to develop separate algorithms for the user equilibrium and system optimal problems.

6.4.3 Price of anarchy

The discussion of the user equilibrium and system optimal problems so far has gone in several directions. On the one hand, it seems like the problems are fairly similar: both involve routes chosen according to some travel time minimization rule, and they are mathematically equivalent. On the other hand, the real-world interpretations are very different (individual control vs. centralized control), and the Braess paradox shows that user equilibrium can behave in a far more

counterintuitive manner than system optimal. This subsection concludes the discussion by showing that, while the user equilibrium solution may be different than the system optimal solution, and worse in terms of higher total system travel time, in many cases there is a limit to how much worse it can be.

In particular, if the link performance functions are linear (that is, of the form $t_{ij}(x_{ij}) = a_{ij} + b_{ij}x_{ij}$), the ratio between the total system travel time at user equilibrium, and the total system travel time at system optimum, can be no greater than $\frac{4}{3}$. This is often called the *price of anarchy*, since it represents the amount of additional delay which can potentially be caused by allowing individual drivers to choose their own routes, compared to a centrally-controlled solution. This bound holds no matter what the network structure or the level of demand, as long as the link performance functions are affine. This is an elegant and perhaps surprising result, which is actually not difficult to show using the variational inequality formulation of equilibrium.

Theorem 6.5. *Given a network with link performance functions of the form $t_{ij}(x_{ij}) = a_{ij} + b_{ij}x_{ij}$ with $a_{ij} \geq 0$ and $b_{ij} \geq 0$ for all links (i, j) , let feasible link flows $\hat{\mathbf{x}} \in X$ satisfy the principle of user equilibrium and $\hat{\mathbf{y}} \in X$ the system optimum principle. Denoting the total system travel times corresponding to these solutions as $TSTT(\hat{\mathbf{x}})$ and $TSTT(\hat{\mathbf{y}})$, we have*

$$\frac{TSTT(\hat{\mathbf{x}})}{TSTT(\hat{\mathbf{y}})} \leq \frac{4}{3} \quad (6.63)$$

Proof. Since $\hat{\mathbf{x}}$ satisfies the principle of user equilibrium, we have

$$\mathbf{t}(\hat{\mathbf{x}}) \cdot (\hat{\mathbf{x}} - \hat{\mathbf{y}}) \leq 0 \quad (6.64)$$

or, equivalently,

$$\mathbf{t}(\hat{\mathbf{x}}) \cdot \hat{\mathbf{x}} \leq \mathbf{t}(\hat{\mathbf{x}}) \cdot \hat{\mathbf{y}} \quad (6.65)$$

$$= (\mathbf{t}(\hat{\mathbf{x}}) - \mathbf{t}(\hat{\mathbf{y}})) \cdot \hat{\mathbf{y}} + \mathbf{t}(\hat{\mathbf{y}}) \cdot \hat{\mathbf{y}} \quad (6.66)$$

which is the same as

$$TSTT(\hat{\mathbf{x}}) \leq (\mathbf{t}(\hat{\mathbf{x}}) - \mathbf{t}(\hat{\mathbf{y}})) \cdot \hat{\mathbf{y}} + TSTT(\hat{\mathbf{y}}). \quad (6.67)$$

If we can show that

$$(\mathbf{t}(\hat{\mathbf{x}}) - \mathbf{t}(\hat{\mathbf{y}})) \cdot \hat{\mathbf{y}} \leq \frac{1}{4} TSTT(\hat{\mathbf{x}}) \quad (6.68)$$

then we have established the result, by substitution into (6.67). We do this by showing an even stronger result, namely

$$(t_{ij}(\hat{x}_{ij}) - t_{ij}(\hat{y}_{ij}))\hat{y}_{ij} \leq \frac{1}{4} t_{ij}(\hat{x}_{ij})\hat{x}_{ij} \quad (6.69)$$

for all links (i, j) which is clearly sufficient.

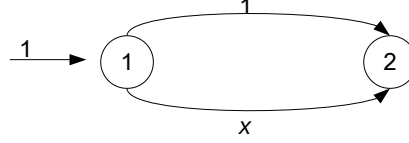


Figure 6.8: The price of anarchy bound is tight.

If $\hat{x}_{ij} \leq \hat{y}_{ij}$, then the left-hand side of (6.69) is nonpositive (the link performance functions are nondecreasing since $b_{ij} \geq 0$), while the right-hand side is nonnegative (since $a_{ij} \geq 0$), so the inequality is certainly true. On the other hand, if $\hat{x}_{ij} > \hat{y}_{ij}$, then

$$(t_{ij}(\hat{x}_{ij}) - t_{ij}(\hat{y}_{ij}))\hat{y}_{ij} = b_{ij}(\hat{x}_{ij} - \hat{y}_{ij})\hat{y}_{ij}. \quad (6.70)$$

Now, the function $b_{ij}(\hat{x}_{ij} - z)z$ is a concave quadratic function in z , which obtains its maximum when $z = \hat{x}_{ij}/2$. Therefore

$$(t_{ij}(\hat{x}_{ij}) - t_{ij}(\hat{y}_{ij}))\hat{y}_{ij} \leq \frac{b_{ij}(\hat{x}_{ij})^2}{4} \leq \frac{1}{4}(a_{ij} + b_{ij}\hat{x}_{ij})\hat{x}_{ij} = t_{ij}(\hat{x}_{ij})\hat{x}_{ij}, \quad (6.71)$$

proving the result. \square

Furthermore, this bound is tight. Consider the more extreme version of the Knight-Pigou-Downs network shown in Figure 6.8 where the demand is 1 unit, the travel time on the upper link is 1 minute, and the travel time on the bottom link is $t^\downarrow = x^\downarrow$. The user equilibrium solution is $x^\uparrow = 0$, $x^\downarrow = 1$, when both links have equal travel times of 1 minute, and the total system travel time is 1. You can verify that the system optimal solution is $x^\uparrow = x^\downarrow = \frac{1}{2}$, when the total system travel time is $\frac{3}{4}$. Thus the ratio between the user equilibrium and system optimal total system travel times is $\frac{4}{3}$.

You may be wondering if a price of anarchy can be found when we relax the assumption that the link performance functions are affine. In many cases, yes; for instance, if the link performance functions are quadratic, then the price of anarchy is $\frac{3\sqrt{3}}{3\sqrt{3}-2}$, if cubic, then the price of anarchy is $\frac{4\sqrt[3]{4}}{4\sqrt[3]{4}-3}$, and so on. In all of these cases the modified Knight-Pigou-Downs network can show that this bound is tight. On the other hand, if the link performance functions have a vertical asymptote (e.g., $t_{ij} = (u_{ij} - x_{ij})^{-1}$), then the ratio between the total system travel times at user equilibrium and system optimum may be made arbitrarily large.

6.5 Exercises

1. [34] Consider a network with two parallel links connecting a single origin to a single destination; the link performance function on each link is $2 + x$ and the total demand is $d = 10$.
 - (a) Write the equations and inequalities defining the set of feasible path flows H , and draw a sketch.
 - (b) What are the vectors $-\mathbf{C}(\mathbf{h})$ for the following path flow vectors? (1) $\mathbf{h} = [0, 10]$ (2) $\mathbf{h} = [5, 5]$ (3) $\mathbf{h} = [10, 0]$ Draw these vectors at these three points on your sketch.
 - (c) For each of the vectors from part (b), identify the point $\text{proj}_H(\mathbf{h} - \mathbf{C}(\mathbf{h}))$ and include these on your sketch.
2. [45] Consider a network with two parallel links connecting a single origin to a single destination; the link performance function on the first link is $10 + x$, and the link performance function on the second link is $x^2/20$. The total demand is $d = 30$.
 - (a) Demonstrate that the solution $\hat{\mathbf{h}} = [20 \ 10]$ is not an equilibrium, and then provide a vector \mathbf{h} showing that the variational inequality (6.2) is not satisfied.
 - (b) Demonstrate that the solution $\hat{\mathbf{h}} = [10 \ 20]$ is an equilibrium, and prove that the variational inequality (6.2) is satisfied no matter what $\mathbf{h} \in H$ is.
3. [43] For the Braess network of Figure 5.8(b), write down the optimality conditions corresponding to Beckmann's formulation (6.13)–(6.16).
4. [27] Which of the following multifunctions F are closed? Each of these multifunctions maps $[0, 1]$ to subsets of $[0, 1]$. Draw sketches of each of these multifunctions.
 - (a) $F(x) = \{y : 0 \leq y \leq x\}$
 - (b) $F(x) = \{y : 0 < y < 1\}$
 - (c) $F(x) = \{0\} \cup \{1 - x\} \cup \{1\}$
 - (d) $F(x) = \{x^2\}$
5. [62] Show that if $f(x)$ is a continuous function with a compact domain, the single-valued “multifunction” $F(x) = \{f(x)\}$ is closed.
6. [23] For each of the multifunctions in Exercise 4, identify all of the fixed points.
7. [62] Specify the multifunction $R(\mathbf{c})$ for the Braess network (Figure 5.8b), and identify all of its fixed points.

8. [42] Complete the proof of Theorem 6.3 by showing that the graph of F is indeed closed.
9. [51] Create a simple network where one or more link performance functions are *not* continuous, and where no user equilibrium solution exists. (Don't worry about creating "realistic" functions for this problem.)
10. [51] Create a simple network where one or more link performance functions are *not* strictly increasing, and where the user equilibrium link flows are not unique.
11. [11] Show that the BPR link delay function (5.1) does not satisfy the condition $t'_{ij}(x_{ij}) > 0$ for all x_{ij} if $\beta > 1$.
12. [46] In spite of Exercise 11, the BPR link delay functions are strictly increasing, and the resulting link flow equilibrium solution is still unique. Generalize the proof of Proposition 6.2 to handle the case when the link delay functions are differentiable and strictly increasing, by showing that the Beckmann function is still strictly convex even if $t'_{ij}(x_{ij})$ is not always strictly positive.
13. [51] Show that the user-equilibrium and system-optimal link flows are the same if there is no congestion, that is, if $t_a(x_a) = \tau_a$ for some constant τ_a .
14. [10] Consider the network of Figure 6.9, ignoring the labels next to the link in the figure. The demand in this network is given by $d^{19} = 400$. Does the following path-flow solution satisfy the principle of user equilibrium? 100 vehicles on path [1,2,3,6,9], 100 vehicles on path [1,2,3,6,9], 100 vehicles on path [1,4,5,6,9], and 100 vehicles on path [1,4,5,2,3,6,9]. You can answer this without any calculation.
15. [45] In the network in Figure 6.9, the demand is given by $d^{13} = d^{19} = d^{17} = d^{91} = d^{93} = d^{97} = 100$. The flow on each link is shown in the figure.
 - (a) Find a path flow vector \mathbf{h} which corresponds to these link flows.
 - (b) Show the resulting origin-based link flows \mathbf{x}^1 and \mathbf{x}^9 from the two origins in the network.
 - (c) Show that the links used by these two origins form an acyclic subnetwork by finding topological orders for each subnetwork.
 - (d) Determine whether these link flows satisfy the principle of user equilibrium.
16. [35] Consider the network in Figure 6.10, along with the given (equilibrium) link flows. There is only one OD pair, from node 1 to node 3. Identify three values of path flows which are consistent with these link flows, *in addition* to the most likely (entropy-maximizing) path flows.

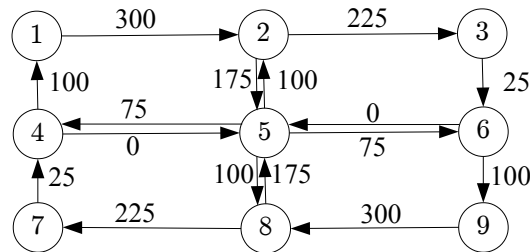


Figure 6.9: Network for Exercises 14 and 15. Each link has link performance function $10 + x_{ij}$

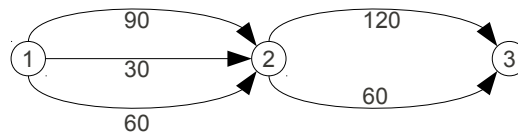


Figure 6.10: Network for Exercise 16.

17. [37] In the network shown in Figure 6.11, 320 vehicles travel from A to C, 640 vehicles travel from A to D, 160 vehicles travel from B to C, and 320 vehicles travel from B to D. The equilibrium link flows are shown.
- Give a path flow solution which satisfies proportionality and produces the equilibrium link flows.
 - At the proportional solution, what fraction of flow on the top link connecting 3 and 4 is from origin A?
 - Is there a path flow solution which produces the equilibrium link flows, yet has *no* vehicles from origin A on the top link connecting 3 and 4? If so, list it. If not, explain why.

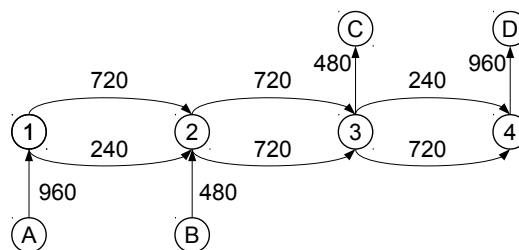


Figure 6.11: Network for Exercise 17.

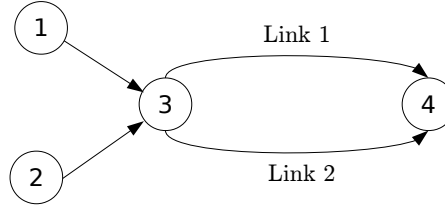


Figure 6.12: Network for Exercise 18.

18. [38] Consider the network in Figure 6.12, where 2 vehicles travel from 1 to 4 (with a value of time of \$20/hr), and 4 vehicles travel from 2 to 4 (with a value of time of \$8/hr). The equilibrium volumes are 3 vehicles on Link 1, and 3 vehicles on Link 2.
 - (a) Assuming that the vehicles in this network are discrete and cannot be split into fractions, identify every combination of path flows which give the equilibrium link volumes (there should be 20). Assuming each combination is equally likely, show that the proportional division of flows has the highest probability of being realized.
 - (b) What is the average value of travel time on Link 1 at the most likely path flows? What are the upper and lower limits on the average value of travel time on this link?
19. [65] Derive the optimality conditions for the system-optimal assignment, and provide an interpretation of these conditions which intuitively relates them to the concept of system optimality.
20. [33] Find the system optimal assignment in the Braess network (Figure 5.8b), assuming a demand of 6 vehicles from node 1 to node 4.
21. [51] Consider a network where every link performance function is linear, of the form $t_{ij} = a_{ij}x_{ij}$. Show that the user equilibrium and system optimum solutions are the same.
22. [34] Calculate the ratio of the total system travel time between the user equilibrium and system optimal solutions in the Braess network (Figure 5.8b).
23. [46] Find the set of boundedly rational assignments in the Braess network (Figure 5.8b).

Chapter 7

Algorithms for Traffic Assignment

This chapter presents algorithms for solving the basic traffic assignment problem (TAP), which was defined in Chapter 6 as the solution $\hat{\mathbf{x}}$ to the variational inequality

$$\mathbf{t}(\hat{\mathbf{x}}) \cdot (\hat{\mathbf{x}} - \mathbf{x}) \leq 0 \quad \forall \mathbf{x} \in X, \quad (7.1)$$

which can also be expressed in terms of path flows $\hat{\mathbf{h}}$ as

$$\mathbf{c}(\hat{\mathbf{h}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) \leq 0 \quad \forall \mathbf{h} \in H, \quad (7.2)$$

or equivalently as the solution $\hat{\mathbf{x}}$ to the optimization problem

$$\min_{\mathbf{x}, \mathbf{h}} \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (7.3)$$

$$\text{s.t. } x_{ij} = \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \quad \forall (i, j) \in A \quad (7.4)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (7.5)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (7.6)$$

While there are general algorithms for variational inequalities or nonlinear optimization problems, TAP involves tens of thousands or even millions of variables and constraints for practical problems. So, these general algorithms are outperformed by specialized algorithms which are designed to exploit some specific features of TAP. The most significant features to exploit are the network structure embedded in the constraints, and the fact that the objective function is separable by link.

This chapter presents four types of algorithms for TAP. The first three are aimed at finding an equilibrium solution (either link flows $\hat{\mathbf{x}}$ or path flows $\hat{\mathbf{h}}$). Broadly speaking, these algorithms can be divided into link-based, path-based,

and bush-based algorithms, according to the way the solution is represented. The chapter focuses entirely on the basic TAP, and not the alternative assignment rules from Section 6.3. System optimal assignment can be transformed mathematically into a user equilibrium problem with modified cost functions, so all of these algorithms can be easily adapted for the system optimal problem, but most of these algorithms cannot be directly applied to the variations with bounded rationality or perception errors.

Section 7.1 provides an introduction, justifying the need for efficient algorithms for TAP and discussing issues such as convergence criteria which are relevant to all algorithms. The next three sections respectively present link-based, path-based, and bush-based algorithms for finding user equilibrium solutions.

7.1 Introduction to Assignment Algorithms

This introductory section touches on three topics: first, the advantages and disadvantages of link-based, path-based, and bush-based algorithms; second, the general framework for solving equilibrium problems; and third, the question of convergence criteria, which is pertinent to all of these.

7.1.1 Why do different algorithms matter?

It may not be clear why we need to present so many different algorithms for the same problem. Why not simply present the one “best” algorithm for solving TAP? Link-based, path-based, and bush-based algorithms all exhibit advantages and disadvantages relative to each other. In this introductory section, these algorithms are compared qualitatively. As you read through the following sections, which include details and specifications of each algorithm, keep these concepts in mind.

Link-based algorithms only keep track of the aggregate link flows \mathbf{x} , not the path flows \mathbf{h} which lead to these link flows. Since the number of links in a large network is much less than the number of paths, link-based algorithms are very economical in terms of computer memory. For this reason, link-based algorithms were the first to be used in practice decades ago, when computer memory was very expensive. Link-based algorithms also tend to be easier to parallelize. With modern desktop machines containing multiple cores, this parallelization can reduce computation time significantly. Finally, link-based algorithms tend to be easy to code and implement (even in a spreadsheet). A significant drawback is that they tend to be slow, particularly when high precision is demanded. The first few iterations of link-based algorithms make good progress, but they stall quickly, and final convergence to the equilibrium solution can be extremely slow.

Path-based algorithms, by contrast, keep track of the path flow vector \mathbf{h} , which can be used to generate the link flows \mathbf{x} whenever needed. Although this representation requires more memory than a link-based solution, it also retains a great deal of information which is lost when one aggregates path flows to

link flows. This information can be exploited to converge much faster than link-based algorithms, especially when a very precise solution is needed. Because the number of paths in a network is so large, much of the effort involved in coding path-based algorithms is associated with clever data structures and algorithm schemes which are aimed at minimizing the number of paths which need to be stored. This increases the complexity of the coding of these algorithms, even if the algorithmic concepts are not difficult.

Bush-based algorithms are the most recent developed, and aim to offer speed comparable to path-based algorithms, while requiring less memory (although still significantly more than link-based algorithms). Bush-based algorithms do this by selectively aggregating the path flows: as shown in Section 6.2.3, at equilibrium the set of paths used by all travelers from the same origin (or traveling to the same destination) forms an acyclic subgraph (a bush). Bush-based algorithms aim to identify these bushes for each origin and destination, exploiting the fact that calculations in acyclic networks are very fast. Bush-based algorithms also tend to return higher-entropy solutions than path-based algorithms, which is important when interpreting the path flow solution. The downside to these algorithms is that their design involves more complexity, and the success of these algorithms depends highly on implementation and programming skill.

7.1.2 Framework

It turns out that none of these solution methods get to the right answer immediately, or even after a finite number of steps. There is no “step one, step two, step three, and then we’re done” recipe for solving large-scale equilibrium problems. Instead, an iterative approach is used where we start with some feasible assignment (link or path), and move closer and closer to the equilibrium solution as you repeat a certain set of steps over and over, until you’re “close enough” to quit and call it good. One iterative algorithm you probably saw in calculus was Newton’s method for finding zeros of a function. In this method, one repeats the same step over and over until the function is sufficiently close to zero.

Broadly speaking, all equilibrium solution algorithms repeat the following three steps:

1. Find the shortest (least travel time) path between each origin and each destination.
2. Shift travelers from slower paths to faster ones.
3. Recalculate link flows and travel times after the shift, and return to step one unless we’re close enough to equilibrium.

The shortest path computation can be done quickly and efficiently even in large networks, as was described in Section 2.4. The third step is even more straightforward, and is nothing more than re-evaluating the link performance

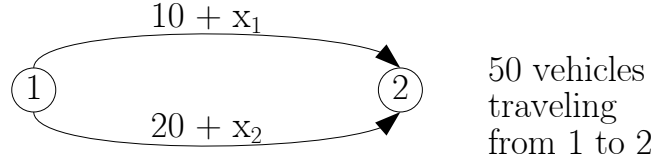


Figure 7.1: Two-link example for demonstration.

functions on each link with the new volumes. The second step requires the most care; the danger here is shifting either too few travelers onto faster paths, or shifting too many. If we shift too few, then it will take a long time to get to the equilibrium solution. On the other hand, systematically shifting too many can be even more dangerous, because it creates the possibility of “infinite cycling” and never finding the true equilibrium.

In the simple example in Figure 7.1, by inspection the equilibrium is for thirty travelers to choose the top route, and twenty to choose the bottom route, with an equal travel time of 40 minutes on both paths. Solving this example using the above process, initially (i.e., with nobody on the network) the fastest path is the top one (step one), so let’s assign all 50 travelers onto the top path (step two). Performing the third step, we recalculate the travel times as 60 minutes on the top link, and 20 on the bottom. This is not at all an equilibrium, so we go back to the first step, and see that the bottom path is now faster, so we have to shift some people from the top to the bottom. If we wanted, we could shift travelers one at a time, that is, assigning 49 to the top route and 1 to the bottom, seeing that we still haven’t found equilibrium, so trying 48 and 2, then 47 and 3, and so forth, until finally reaching the equilibrium with 30 and 20. Clearly this is not efficient, and is an example of shifting too few travelers at a time.

At the other extreme, let’s say we shift *everybody* onto the fastest path in the second step. That is, we go from assigning 50 to the top route and 0 to the bottom, to assigning 0 to the top and 50 to the bottom. Recalculating link travel times, the top route now has a travel time of 10 minutes, and the bottom a travel time of 70. This is even worse!¹ Repeating the process, we try to fix this by shifting everybody back (50 on top, 0 on bottom), but now we’re just back in the original situation. If we kept up this process, we’d keep bouncing back and forth between these solutions. This is clearly worse than shifting too few, because we never reach the equilibrium no matter how long we work! You might think it’s obvious to detect if something like this is happening. With this small example, it might be. Trying to train a computer to detect this, or trying to detect cycles with over millions OD pairs (as is common in practice), is much harder. The key step in all of the algorithms for finding equilibria is determining how much flow to shift.

¹By “worse” we mean farther from equilibrium.

7.1.3 Convergence criteria

A general issue is how one chooses to stop the iterative process, that is, how one knows when a solution is “good enough” or close enough to equilibrium. This is called a *convergence criterion*. Many convergence criteria have been proposed over the years; perhaps the most common in practice is the *relative gap*, which is defined first. Unfortunately, the relative gap has been defined in several different ways, so it is important to be familiar with all of the definitions.

Remembering that the multiplier κ^{rs} represents the time spent on the fastest path between origin r and destination s , one definition of the relative gap γ_1 is:

$$\gamma_1 = \frac{\sum_{(i,j) \in A} t_{ij} x_{ij}}{\sum_{(r,s) \in Z^2} \kappa^{rs} d^{rs}} - 1 = \frac{\mathbf{t} \cdot \mathbf{x}}{\boldsymbol{\kappa} \cdot \mathbf{d}} - 1 \quad (7.7)$$

The numerator of the fraction is the total system travel time (TSTT). The relative gap is always nonnegative, and it is equal to zero if and only if the flows \mathbf{x} satisfy the principle of user equilibrium. It is these properties which make the relative gap a useful convergence criterion: once it is close enough to zero, our solution is “close enough” to equilibrium. For most practical purposes, a relative gap of 10^{-4} – 10^{-6} is small enough.

A second definition of the relative gap γ_2 is based on the Beckmann function itself. Let f denote the Beckmann function, and \hat{f} its value at equilibrium (which is a global minimum). In many algorithms, given a current solution \mathbf{x} , it is not difficult to generate upper and lower bounds on \hat{f} based on the current solution, respectively denoted $\bar{f}(\mathbf{x})$ and $\underline{f}(\mathbf{x})$. A trivial upper bound is its value at the current solution: $\bar{f} = f(\mathbf{x})$, since clearly $\hat{f} \leq f(\mathbf{x})$ for any feasible link assignment \mathbf{x} . Sometimes, a corresponding lower bound can be identified as well. Assuming that these bounds can become tighter over time, and that in the limit both $\bar{f} \rightarrow \hat{f}$ and $\underline{f} \rightarrow \hat{f}$, the difference or gap $\bar{f}(\mathbf{x}) - \underline{f}(\mathbf{x})$ can be used as a convergence criterion. These values are typically normalized, leading to one definition of the relative gap:

$$\gamma_2 = \frac{\bar{f} - \underline{f}}{\underline{f}}, \quad (7.8)$$

or a slightly modified version

$$\gamma_3 = \frac{\bar{f} - \max \underline{f}}{\max \underline{f}}, \quad (7.9)$$

where $\max \underline{f}$ is the greatest lower bound found to date, in case the sequence of \underline{f} values is not monotone over iterations. A disadvantage of these definitions of the relative gap is that different algorithms calculate these upper and lower bounds differently. While they are suitable as termination criteria in an algorithm, it is not possible (or at least not easy) to directly compare the relative gap calculated by one algorithm to that produced by another to assess which of two solutions is closer to equilibrium.

One drawback of the relative gap (in all of its forms) is that it is unitless and does not have an intuitive meaning. Furthermore, it can be somewhat confusing to have several slightly different definitions of the relative gap, even though they all have the same flavor. A more recently proposed metric is the *average excess cost*, defined as

$$AEC = \frac{\sum_{(i,j) \in A} t_{ij} x_{ij} - \sum_{(r,s) \in Z^2} \kappa^{rs} d^{rs}}{\sum_{(r,s) \in Z^2} d^{rs}} = \frac{\mathbf{t} \cdot \mathbf{x} - \boldsymbol{\kappa} \cdot \mathbf{d}}{\mathbf{d} \cdot \mathbf{1}} \quad (7.10)$$

This quantity represents the average difference between the travel time on each traveler's *actual* path, and the travel time on the *shortest* path available to him or her. Unlike the relative gap, *AEC* has units of time, and is thus easier to interpret.

Another convergence measure with time units is the *maximum excess cost*, which relates directly to the principle of user equilibrium. The maximum excess cost is defined as the largest amount by which a used path's travel time exceeds the shortest path travel time available to that traveler:

$$MEC = \max_{(r,s) \in Z^2} \left\{ \max_{\pi \in \Pi^{rs}: h^\pi > 0} \{c^\pi - \kappa^{rs}\} \right\} \quad (7.11)$$

This is often a few orders of magnitude higher than the average excess cost. One disadvantage of the maximum excess cost is that it is only applicable when the path flow solution is known. This is easy in path-based or bush-based algorithms. However, since many path-flow solutions correspond to the same link-flow solution (cf. Section 6.2.2), *MEC* is not well suited for link-based algorithms.

Finally, this section concludes with two convergence criteria which are inferior to those discussed thus far. The first is to simply use the Beckmann function itself; when it is sufficiently close to the global optimal value \hat{f} , terminate. A moment's thought should convince you that this criterion is not practical: there is no way to know the value of \hat{f} until the problem has already been solved. (Upper and lower bounds are possible to calculate, though, as with γ_2 .) A more subtle version is to terminate when the Beckmann function stops decreasing, or (in a more common form) to terminate the algorithm when the link or path flows stabilize from one iteration to the next. The trouble with these convergence criteria is that they cannot distinguish between a situation when the flows stabilize because they are close to the equilibrium solution, and when they stabilize because the algorithm “gets stuck” and cannot improve further due to a flaw in its design or a bug in the programming. For this reason, it is always preferable to base the termination criteria on the equilibrium principle itself.

7.2 Link-Based Algorithms

Link-based algorithms for traffic assignment are the simplest to understand and implement, and require the least amount of computer memory. Given a current

set of link flows \mathbf{x} , a link-based algorithm attempts to move closer to equilibrium by performing two steps. First, a target point \mathbf{x}^* is identified; moving in the direction of \mathbf{x}^* from \mathbf{x} should lead towards an equilibrium solution. Then, a step of size $\lambda \in [0, 1]$ is taken in this direction, updating the link flows to $\lambda\mathbf{x}^* + (1 - \lambda)\mathbf{x}$. This is a specific way of implementing the first two steps of the framework specified in the previous section: λ represents the fraction of flow which is shifted from the paths at the current solution \mathbf{x} to the paths at the target solution \mathbf{x}^* . If $\lambda = 1$, all of the flow has shifted to the target solution; if $\lambda = 0$, no flow has shifted at all. Since the set of feasible link flow solutions X is convex, as long as \mathbf{x} and \mathbf{x}^* are feasible, we can be assured that the convex combination $\lambda\mathbf{x}^* + (1 - \lambda)\mathbf{x}$ is feasible as well.

Link-based algorithms differ in two primary ways: first, how the target \mathbf{x}^* is chosen; and second, how the step size λ is chosen. If λ is too small, convergence to equilibrium will be very slow, but if λ is too large, the solution may never converge at all — the example from Figure 7.1 in the previous section is an example of what can happen if $\lambda = 1$ for all iterations. This section presents three link-based algorithms, in increasing order of sophistication (but in increasing order of convergence speed.) The first is the method of successive averages (MSA), which is perhaps the simplest equilibrium algorithm. The second is the Frank-Wolfe algorithm, which can be thought of as a version of MSA with a more intelligent choice of step size λ . (Frank-Wolfe was the most common used in practice for several decades.) The third is the conjugate Frank-Wolfe algorithm, which can be thought of as a version of Frank-Wolfe with a more intelligent choice of target \mathbf{x}^* .

All of these algorithms use the following framework; the only difference is how \mathbf{x}^* and λ are calculated.

1. Generate an initial solution $\mathbf{x} \in X$.
2. Generate a target solution $\mathbf{x}^* \in X$.
3. Update the current solution: $\mathbf{x} \leftarrow \lambda\mathbf{x}^* + (1 - \lambda)\mathbf{x}$ for some $\lambda \in [0, 1]$.
4. Calculate the new link travel times.
5. If the convergence criterion is satisfied, stop; otherwise return to step 2.

7.2.1 Method of successive averages

Although the method of successive averages (MSA) is not competitive with other equilibrium solution algorithms, its simplicity and clarity in applying the three-step iterative process make it an ideal starting place. To specify MSA, we need to specify how the target solution \mathbf{x}^* is chosen, and how the step size λ is chosen.

The target \mathbf{x}^* is an all-or-nothing assignment. That is, *assuming that the current travel times are fixed*, identify the shortest path between each origin and destination, and load all of the demand for that OD pair onto that path. Thus, \mathbf{x}^* is the state which would occur if literally every driver was to switch

paths onto what is currently the shortest path. Of course, if we were to switch everybody onto these paths, which would occur if we choose $\lambda = 1$, those paths would almost certainly not be “shortest” anymore. But \mathbf{x}^* can still be thought of as a target, or a direction in which travelers would feel pressure to move.

So, what should λ be? As discussed above, there are problems if you shift too few travelers, and potentially even bigger problems if you shift too many. MSA adopts a reasonable middle ground: initially, we shift a lot of travelers, but as the algorithm progresses, we shift fewer and fewer until we settle down on the average. The hope is that this avoids both the problems of shifting too few (at first, we’re taking big steps, so hopefully we get somewhere close to equilibrium quickly) and of shifting too many (eventually, we’ll only be moving small amounts of flow so there is no worry of infinite cycling).

Specifically, on the i -th iteration, MSA uses $\lambda = 1/(i + 1)$. So, the first time through, half of the travelers are shifted to the current shortest paths. The second time through, a third of the people shift to the current shortest paths (and two thirds stay on their current path). On the third iteration, a fourth of the people shift to new paths, and so on. (MSA can also be applied with different step size rules; see Exercise 11.)

At this point, it’s worth using the Beckmann formulation to show that the choice of an all-or-nothing assignment for \mathbf{x}^* has mathematical justification, in addition to the intuitive interpretation of shifting towards shortest path. Let \mathbf{x} be the current set of link flows, and \mathbf{x}^* an all-or-nothing assignment. As a result of a shift of size λ , the Beckmann function will be changed from $f(\mathbf{x})$ to $f(\mathbf{x}')$, and we want to show that it’s possible to choose λ in some way to guarantee $f(\mathbf{x}') \leq f(\mathbf{x})$. That is, we want to show that we can reduce the Beckmann function (and thus move closer to the equilibrium solution) by taking a (correctly-sized) step in the direction $\mathbf{x}^* - \mathbf{x}$. Define $f(\mathbf{x}(\lambda)) = f((1 - \lambda)\mathbf{x} + \lambda\mathbf{x}^*)$ to be the Beckmann function after taking a step of size λ . Using the multivariate chain rule, the derivative of $f(\lambda)$ is

$$\frac{df}{d\lambda} = \sum_{(i,j) \in A} \frac{\partial f}{\partial x_{ij}} \frac{dx_{ij}}{d\lambda} = \sum_{(i,j) \in A} t_{ij}((1 - \lambda)x_{ij} + \lambda x_{ij}^*)(x_{ij}^* - x_{ij})$$

Evaluating this derivative at $\lambda = 0$ gives

$$\frac{df}{d\lambda}(0) = \sum_{(i,j) \in A} t_{ij}(x_{ij})(x_{ij}^* - x_{ij})$$

Now, \mathbf{x}^* was specifically chosen to put all vehicles on the shortest paths at travel times $\mathbf{t}(\mathbf{x})$, and so $\sum_{(i,j) \in A} x_{ij}^* t_{ij}(x_{ij}) \leq \sum_{(i,j) \in A} x_{ij} t_{ij}(x_{ij})$, and therefore $\frac{df}{d\lambda}(0) \leq 0$. Furthermore, if we are not at the equilibrium solution already, $\sum_{(i,j) \in A} x_{ij}^* t_{ij}(x_{ij})$ is strictly less than $\sum_{(i,j) \in A} x_{ij} t_{ij}(x_{ij})$. This implies $\frac{df}{d\lambda}(0) < 0$ or, equivalently, we can decrease the Beckmann function if we take a small enough step in the direction $\mathbf{x}^* - \mathbf{x}$, by shifting people from longer paths onto shorter ones.

Two examples of MSA are shown below. A proof of convergence is sketched in Exercise 12.

Small network example Here we solve the small example of Figure 7.1 by MSA, using the relative gap γ_1 to measure how close we are to equilibrium.

Initialization. Find the shortest paths: with no travelers on the network, the top link has a travel time of 10, and the bottom link has a travel time of 20. Therefore the top link is the shortest path, so $\mathbf{x}^* = [50 \ 0]$. We take this to be the initial solution $\mathbf{x} \leftarrow \mathbf{x}^* = [50 \ 0]$. Recalculating the travel times, we have $t_1 = 10 + x_1 = 60$ and $t_2 = 20 + x_2 = 20$ (or, in vector form, $\mathbf{t} = [60 \ 20]$).

Iteration 1. With the new travel times, the shortest path is now the bottom link, so $\kappa = 20$ and the relative gap is

$$\gamma_1 = \frac{\mathbf{t} \cdot \mathbf{x}}{\kappa \cdot \mathbf{d}} - 1 = \frac{50 \times 60 + 0 \times 20}{20 \times 50} - 1 = 2$$

This is far too big, so we continue with the second iteration. If everyone were to take the new shortest path, the flows would be $\mathbf{x}^* = [0 \ 50]$. Because this is the first iteration, we shift 1/2 of the travelers onto this path, so $\mathbf{x} \leftarrow (1/2)\mathbf{x}^* + (1/2)\mathbf{x} = [0 \ 25] + [25 \ 0] = [25 \ 25]$. The new travel times are thus $\mathbf{t} = [35 \ 45]$.

Iteration 2. With the new travel times, the shortest path is now the top link, so $\kappa = 35$ and the relative gap is

$$\gamma_1 = \frac{\mathbf{t} \cdot \mathbf{x}}{\kappa \cdot \mathbf{d}} - 1 = \frac{25 \times 35 + 25 \times 45}{35 \times 50} - 1 = 0.143$$

If everyone were to take the new shortest path, the flows would be $\mathbf{x}^* = [50 \ 0]$. Because this is the second iteration, we shift 1/3 of the travelers onto this path, so $\mathbf{x} \leftarrow (1/3)\mathbf{x}^* + (2/3)\mathbf{x} = [50/3 \ 0] + [50/3 \ 50/3] = [100/3 \ 50/3]$. The new travel times are thus $\mathbf{t} = [43.33 \ 36.67]$.

Iteration 3. With the new travel times, the shortest path is now the bottom link, so $\kappa = 36.67$ and the relative gap is $\gamma_1 = 0.121$. A bit better, but still too big, so we carry on. Here $\mathbf{x}^* = [0 \ 50]$, $\mathbf{x} \leftarrow (1/4)\mathbf{x}^* + (3/4)\mathbf{x} = [0 \ 50/4] + [25 \ 50/4] = [25 \ 25]$. The new travel times are $\mathbf{t} = [35 \ 45]$. Note that we have returned to the same solution found in Iteration 1. *Don't despair; this just means the last shift was too big. Next time we'll shift fewer vehicles (because λ is smaller with each iteration).*

Iteration 4. With the new travel times, the shortest path is now the top link, so $\kappa = 35$ and the relative gap is $\gamma_1 = 0.143$. The new target is $\mathbf{x}^* = [50 \ 0]$, $\mathbf{x} \leftarrow (1/5)\mathbf{x}^* + (4/5)\mathbf{x} = [30 \ 20]$. The new travel times are $\mathbf{t} = [40 \ 40]$. With the new travel times, the shortest path is the top link, so $\kappa = 40$ and the relative gap is $\gamma_1 = 0$, so we stop. *In fact, either path could have been chosen for the shortest path. Whenever there is a tie between shortest paths, you are free to choose among them.*

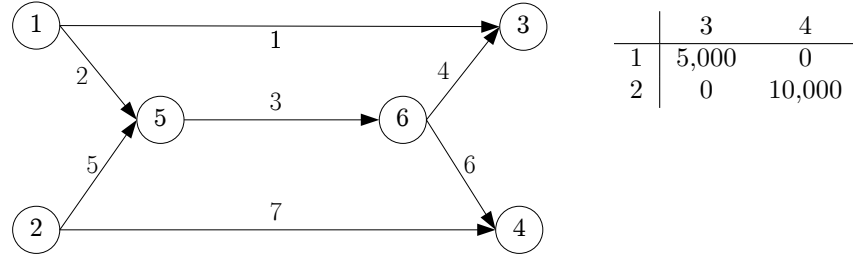


Figure 7.2: Larger example with two OD pairs. (Link numbers shown.)

Larger network example Here we apply MSA to a slightly larger network with two OD pairs, shown in Figure 7.2, where each link has the link performance function $t(x) = 10 + x/100$.

There are four paths in this network; for OD pair (1,3) these are denoted $[1, 3]$ and $[1, 5, 6, 3]$ according to their link numbers, and for OD pair (2,4) these are $[2, 5, 6, 4]$ and $[2, 4]$. In this example, we'll calculate the average excess cost, rather than the relative gap.

Initialization. Find the shortest paths: with no travelers on the network, paths $[1, 3]$, $[1, 5, 6, 3]$, $[2, 5, 6, 4]$, and $[2, 4]$ respectively have travel times of 10, 30, 30, and 10. Therefore $[1, 3]$ is shortest for OD pair (1,3), and $[2, 4]$ is shortest for OD pair (2,4), so $\mathbf{x}^* = [5000 \ 0 \ 0 \ 0 \ 0 \ 0 \ 10000]^2$. Since is the first iteration, we simply set

$$\mathbf{x} \leftarrow \mathbf{x}^* = [5000 \ 0 \ 0 \ 0 \ 0 \ 0 \ 10000]$$

Recalculating the travel times, we have

$$\mathbf{t} = [60 \ 10 \ 10 \ 10 \ 10 \ 10 \ 110]$$

Iteration 1. With the new travel times, the shortest path for (1,3) is now $[1, 5, 6, 3]$, with a travel time of 30, so $\kappa^{13} = 30$. Likewise, the new shortest path for (2,4) is $[2, 5, 6, 4]$, so $\kappa^{24} = 30$ and the average excess cost is

$$\begin{aligned} AEC &= \frac{\mathbf{t} \cdot \mathbf{x} - \boldsymbol{\kappa} \cdot \mathbf{d}}{\mathbf{d} \cdot \mathbf{1}} \\ &= \frac{5000 \times 60 + 10000 \times 110 - 30 \times 5000 - 30 \times 10000}{5000 + 10000} = 63.33 \end{aligned}$$

This is far too big and suggests that the average trip is 63 minutes slower than the shortest paths available! If everyone were to take the new shortest paths, the flows would be

$$\mathbf{x}^* = [0 \ 5000 \ 15000 \ 5000 \ 10000 \ 10000 \ 0]$$

²For each OD pair, we add the total demand from the OD matrix onto each link in the shortest path.

(Be sure you understand how we calculated this.) Because this is iteration 1, we shift 1/2 of the travelers onto this path, so

$$\mathbf{x} \leftarrow (1/2)\mathbf{x}^* + (1/2)\mathbf{x} = [2500 \quad 2500 \quad 7500 \quad 2500 \quad 5000 \quad 5000 \quad 5000]$$

The new travel times are thus

$$\mathbf{t} = [35 \quad 35 \quad 85 \quad 35 \quad 60 \quad 60 \quad 60]$$

Iteration 2. With the new travel times, the shortest path for (1,3) is now [1, 3], with $\kappa^{13} = 35$. The new shortest path for (2,4) is [2, 4], so $\kappa^{24} = 60$ and the average excess cost is

$$\begin{aligned} AEC = & \frac{\left(2500 \times 35 + 2500 \times 35 + 7500 \times 85 + 2500 \times 35 + 5000 \times 60 \right. \\ & \left. + 5000 \times 60 + 5000 \times 60 - 35 \times 5000 - 60 \times 10000 \right)}{15000} \\ & = 68.33 \end{aligned}$$

This is still big (and in fact worse), but we persistently continue with the second iteration. If everyone were to take the new shortest paths, the flows would be

$$\mathbf{x}^* = [5000 \quad 0 \quad 0 \quad 0 \quad 0 \quad 0 \quad 10000]$$

so

$$\mathbf{x} \leftarrow (1/3)\mathbf{x}^* + (2/3)\mathbf{x} = [3333 \quad 1667 \quad 5000 \quad 1667 \quad 3333 \quad 3333 \quad 6667]$$

The new travel times are

$$\mathbf{t} = [43.3 \quad 26.7 \quad 60 \quad 26.7 \quad 43.3 \quad 43.3 \quad 76.7]$$

Iteration 3. With the new travel times, the shortest path for (1,3) is still [1, 3], with $\kappa^{13} = 43.3$, and the shortest path for (2,4) is still [2, 4] with $\kappa^{24} = 76.7$ and the average excess cost is $AEC = 23.6$. Continuing the fourth iteration, as before

$$\mathbf{x}^* = [5000 \quad 0 \quad 0 \quad 0 \quad 0 \quad 0 \quad 10000]$$

so

$$\mathbf{x} \leftarrow (1/4)\mathbf{x}^* + (3/4)\mathbf{x} = [3750 \quad 1250 \quad 3750 \quad 1250 \quad 2500 \quad 2500 \quad 7500]$$

The new travel times are

$$\mathbf{t} = [47.5 \quad 22.5 \quad 47.5 \quad 22.5 \quad 35 \quad 35 \quad 85]$$

Iteration 4. With the new travel times, the shortest path for (1,3) is still [1, 3], with $\kappa^{13} = 47.5$, and the shortest path for (2,4) is still [2, 4] with $\kappa^{24} = 85$ and the average excess cost is $AEC = 9.42$. Continuing the fifth iteration, as before

$$\mathbf{x}^* = [5000 \quad 0 \quad 0 \quad 0 \quad 0 \quad 0 \quad 10000]$$

so

$$\mathbf{x} \leftarrow (1/5)\mathbf{x}^* + (4/5)\mathbf{x} = [4000 \quad 1000 \quad 3000 \quad 1000 \quad 2000 \quad 2000 \quad 8000]$$

The new travel times are

$$\mathbf{t} = [50 \quad 20 \quad 40 \quad 20 \quad 30 \quad 30 \quad 90]$$

Iteration 5. With the new travel times, the shortest path for (1,3) is still [1, 3], with $\kappa^{13} = 50$, and the shortest path for (2,4) is still [2, 4] with $\kappa^{24} = 90$ and the average excess cost is $AEC = 3.07$. *Note that the shortest paths have stayed the same over the last three iterations. This means that we really could have shifted more flow than we actually did. The Frank-Wolfe algorithm, described in the next section, fixes this problem.* We have

$$\mathbf{x}^* = [5000 \quad 0 \quad 0 \quad 0 \quad 0 \quad 0 \quad 10000]$$

so

$$\mathbf{x} \leftarrow (1/6)\mathbf{x}^* + (5/6)\mathbf{x} = [4167 \quad 833 \quad 2500 \quad 833 \quad 1667 \quad 1667 \quad 8333]$$

The new travel times are

$$\mathbf{t} = [51.7 \quad 18.3 \quad 35 \quad 18.3 \quad 26.7 \quad 26.7 \quad 93.3]$$

Iteration 6. With the new travel times, the shortest path for (1,3) is still [1, 3], with $\kappa^{13} = 51.7$, but the shortest path for (2,4) is now [2, 5, 6, 4] with $\kappa^{24} = 88.3$. The average excess cost is $AEC = 3.80$. *Note that the OD pairs are no longer behaving “symmetrically,” the shortest path for (1,3) stayed the same, but the shortest path for (2,4) has changed.* We have

$$\mathbf{x}^* = [5000 \quad 0 \quad 10000 \quad 0 \quad 10000 \quad 10000 \quad 0]$$

so

$$\mathbf{x} \leftarrow (1/7)\mathbf{x}^* + (6/7)\mathbf{x} = [4286 \quad 714 \quad 3571 \quad 714 \quad 2857 \quad 2857 \quad 7142]$$

The new travel times are

$$\mathbf{t} = [52.9 \quad 17.1 \quad 45.7 \quad 17.1 \quad 38.6 \quad 38.6 \quad 81.4]$$

This process continues over and over until the average excess cost is sufficiently small. Even with such a small network, MSA requires a very long time to converge. An average excess cost of 1 is obtained after eleven iterations, 0.1 after sixty-three iterations, 0.01 after three hundred thirty-two, and the rate of convergence only slows down from there.

7.2.2 Frank-Wolfe

One of the biggest drawbacks with MSA is that it has a fixed step size. Iteration i moves exactly $1/(i+1)$ of the travelers onto the new shortest paths, no matter how close or far away we are from the equilibrium. Essentially, MSA decides its course of action before it even gets started, then sticks stubbornly to the plan of moving $1/(i+1)$ travelers each iteration. The Frank-Wolfe (FW) algorithm fixes this problem by using an *adaptive* step size. At each iteration, FW calculates exactly the right amount of flow to shift to get as close to equilibrium as possible.

We might try to do this by picking λ to minimize the relative gap or average excess cost, but this turns out to be harder to compute. Instead, we pick λ to solve a “restricted” VI where the feasible set is the line segment connecting \mathbf{x} and \mathbf{x}^* . It turns out that this is the same as choosing λ to minimize the Beckmann function (7.3) along this line segment. Both approaches for deriving the step size λ are discussed below.

Define X' to be the link flows lying on the line segment between \mathbf{x} and \mathbf{x}^* . That is, $X' = \{\mathbf{x}' : \mathbf{x} = \lambda\mathbf{x}^* + (1-\lambda)\mathbf{x} \text{ for some } \lambda \in [0, 1]\}$. The restricted VI is: find $\hat{\mathbf{x}}' \in X'$ such that $\mathbf{t}(\hat{\mathbf{x}}') \cdot (\hat{\mathbf{x}}' - \mathbf{x}') \leq 0$ for all $\mathbf{x}' \in X'$.

This VI is simple enough to be solved as a single equation. The set X has two endpoints (\mathbf{x} and \mathbf{x}^* , corresponding to $\lambda = 0$ and $\lambda = 1$, respectively). For now, assume that the solution $\hat{\mathbf{x}}'$ to the VI is not at one of these endpoints.³ In this case, the force vector $-\mathbf{t}(\hat{\mathbf{x}}')$ is perpendicular to the direction $\mathbf{x}^* - \mathbf{x}$. (Figure 7.3), so $-\mathbf{t}(\hat{\mathbf{x}}') \cdot (\mathbf{x}^* - \mathbf{x}) = 0$. Writing this equation out in terms of individual components, we need to solve

$$\sum_{ij} t_{ij}(\hat{x}'_{ij}) (x^*_{ij} - x_{ij}) = 0 \quad (7.12)$$

or equivalently

$$\sum_{ij} t_{ij}(\hat{x}'_{ij}) x^*_{ij} = \sum_{ij} t_{ij}(\hat{x}'_{ij}) x_{ij} \quad (7.13)$$

The same equation can be derived based on the Beckmann function. Recall the discussion above, where we wrote the function $f(\mathbf{x}(\lambda)) = f((1-\lambda)\mathbf{x} + \lambda\mathbf{x}^*)$ to be the value of the Beckmann function after taking a step of size λ , and furthermore found the derivative of $f(\mathbf{x}(\lambda))$ to be

$$\frac{df}{d\lambda} = \sum_{(i,j) \in A} t_{ij}((1-\lambda)x_{ij} + \lambda x^*_{ij}) (x^*_{ij} - x_{ij}) \quad (7.14)$$

It is not difficult to show that $f(\mathbf{x}(\lambda))$ is a convex function of λ , so we can find its minimum by setting the derivative equal to zero, which occurs if the condition

$$\sum_{(i,j) \in A} x^*_{ij} t_{ij}((1-\lambda)x_{ij} + \lambda x^*_{ij}) = \sum_{(i,j) \in A} x_{ij} t_{ij}((1-\lambda)x_{ij} + \lambda x^*_{ij}) \quad (7.15)$$

³Exercise 13 asks you to show that the solution methods provided below will still give the right answer even in these cases.

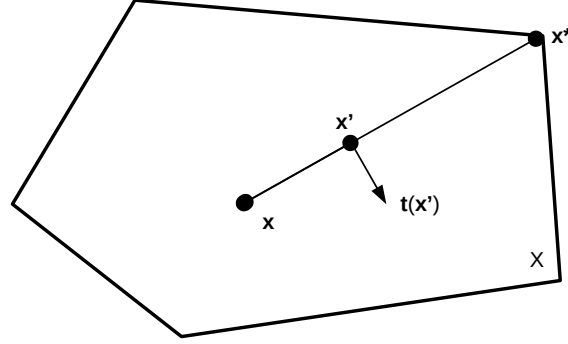


Figure 7.3: A solution to the restricted VI in the Frank-Wolfe method.

is satisfied, which is the same as (7.13).

Study these equations carefully: the coefficients x_{ij} and x_{ij}^* are constants and do not change with λ ; the only part of this condition which is affected by λ are the travel times. You can interpret this equation as trying to find a balance between \mathbf{x} and \mathbf{x}^* in the following sense: different values of λ correspond to shifting a different number of travelers from their current paths to shortest paths, which will result in different travel times on all the links. You want to pick λ so that, *after* you make the switch, both the old paths \mathbf{x} and the old shortest paths \mathbf{x}^* are equally attractive in terms of their travel times.

It is convenient to write equation (7.12) or (7.15) as a function in terms of λ as follows:

$$\zeta'(\lambda) = \sum_{ij} t_{ij}(\lambda x_{ij}^* + (1 - \lambda)x_{ij}) (x_{ij}^* - x_{ij}) = 0 \quad (7.16)$$

which we need to solve for $\lambda \in [0, 1]$ and $\zeta(\lambda)$ is used as a shorthand for $f(\mathbf{x}(\lambda))$. Since the link performance functions are typically nonlinear, we cannot expect to be able to solve this equation analytically to get an explicit formula for λ . General techniques such as Newton's method or an equation solver can be used; but it's not too difficult to use an enlightened trial-and-error method such as a binary search or bisection. Details for all of these line search techniques were presented in Section 4.4.1, and will not be repeated here.

To summarize: in FW, \mathbf{x}^* is an all-or-nothing assignment (just as with MSA). The difference is that λ is chosen to solve (7.16). So, there is a little bit more work at each iteration (we have to solve an equation for λ instead of using a precomputed formula as in MSA), but the reward is much faster convergence to the equilibrium solution. Two examples of FW now follow; you are asked to provide a proof of correctness in Exercise 14 and 15.

7.2.3 Small network example

Here we solve the small example of Figure 7.1 by FW. Some steps are similar to MSA, and therefore omitted. Here, when we do the bisection method, we do five interval reductions, so we are within $1/2^5 = 1/32$ of the correct λ^* value. When solving by computer, you would usually perform more steps than this, because the bisection calculations are very fast.

Initialization. As before, we load everybody on the initial shortest path, so $\mathbf{x} = \mathbf{x}^* = [50 \ 0]$ and $\mathbf{t} = [60 \ 20]$

Iteration 1. As before, the relative gap is $\gamma_1 = 2$. With the new shortest paths, $\mathbf{x}^* = [0 \ 50]$. Begin the bisection method.

Bisection Iteration 1. Initially $\lambda^* \in [0, 1]$. Calculate $\zeta'(1/2) = (0 - 50) \times (10 + 25) + (50 - 0) \times (20 + 25) = 500 > 0$ so we discard the upper half.

Bisection Iteration 2. Now we know $\lambda^* \in [0, 1/2]$. Calculate

$$\zeta'(1/4) = (0 - 50) \times (10 + 37.5) + (50 - 0) \times (20 + 12.5) = -750 < 0$$

so we discard the lower half.

Bisection Iteration 3. Now we know $\lambda^* \in [1/4, 1/2]$. Calculate

$$\zeta'(3/8) = (0 - 50) \times (10 + 31.25) + (50 - 0) \times (20 + 18.75) = -125 < 0$$

so we discard the lower half.

Bisection Iteration 4. Now we know $\lambda^* \in [3/8, 1/2]$. Calculate

$$\zeta'(7/16) = (0 - 50) \times (10 + 28.125) + (50 - 0) \times (20 + 21.875) = 187.5$$

so we discard the upper half.

Bisection Iteration 5. Now we know $\lambda^* \in [3/8, 7/16]$. Calculate

$$\zeta'(13/32) = (0 - 50) \times (10 + 29.6875) + (50 - 0) \times (20 + 20.3125) = 31.25$$

so we discard the upper half.

From here we take the midpoint of the last interval $[3/8, 13/32]$ to estimate $\lambda^* \approx 25/64 = 0.390625$, so $\mathbf{x} = (25/64)\mathbf{x}^* + (39/64)\mathbf{x} = [30.47 \ 19.53]$ and $\mathbf{t} = [40.47 \ 39.53]$.

Iteration 2. The relative gap is calculated as $\gamma_1 = 0.014$. (This is an order of magnitude smaller than the relative gap MSA found by this point.) The shortest paths are still $\mathbf{x}^* = [0 \ 50]$, and we begin bisection.

Bisection Iteration 1. Initially $\lambda^* \in [0, 1]$. Calculate

$$\zeta'(1/2) = 900 > 0$$

so we discard the upper half.

Bisection Iteration 2. Now we know $\lambda^* \in [0, 1/2]$. Calculate

$$\zeta'(1/4) = 435 > 0$$

so we discard the upper half.

Bisection Iteration 3. Now we know $\lambda^* \in [0, 1/4]$. Calculate

$$\zeta'(1/8) = 203 > 0$$

so we discard the upper half.

Bisection Iteration 4. Now we know $\lambda^* \in [0, 1/8]$. Calculate

$$\zeta'(1/16) = 87 > 0$$

so we discard the upper half.

Bisection Iteration 5. Now we know $\lambda^* \in [0, 1/16]$. Calculate

$$\zeta'(1/32) = 29 > 0$$

so we discard the upper half.

The midpoint of the final interval is $\lambda^* \approx 1/64$, so $\mathbf{x} = (1/64)\mathbf{x}^* + (63/64)\mathbf{x} = [29.99 \ 20.01]$ and $\mathbf{t} = [39.99 \ 40.01]$.

Iteration 3. The relative gap is now $\gamma_1 = 0.00014$, so we quit and claim we have found flows that are “good enough” (the difference in travel times between the routes is less than a second).

Alternately, using calculus, we could have identified λ^* during the second iteration as *exactly* 0.40, which would have found the exact equilibrium after only one step.

7.2.4 Large network example

Here we apply FW to the network shown in Figure 7.2, using the same notation as in the MSA example.

Initialization. Path $[1, 3]$ is shortest for OD pair $(1, 3)$, and path $[2, 4]$ is shortest for OD pair $(2, 4)$, so

$$\mathbf{x}^* = [5000 \ 0 \ 0 \ 0 \ 0 \ 0 \ 10000]$$

and

$$\mathbf{x} = \mathbf{x}^* = [5000 \ 0 \ 0 \ 0 \ 0 \ 0 \ 10000]$$

Recalculating the travel times, we have

$$\mathbf{t} = [60 \ 10 \ 10 \ 10 \ 10 \ 10 \ 110]$$

Iteration 1. With the new travel times, the shortest path for (1,3) is now [1, 5, 6, 3], and the new shortest path for (2,4) is [2, 5, 6, 4], so $AEC = 63.33$. If everyone were to take the new shortest paths, the flows would be

$$\mathbf{x}^* = [0 \quad 5000 \quad 15000 \quad 5000 \quad 10000 \quad 10000 \quad 0]$$

Begin the bisection method to find the right combination of \mathbf{x}^* and \mathbf{x} .

Bisection Iteration 1. Initially $\lambda^* \in [0, 1]$. Calculate

$$\begin{aligned} \zeta'(1/2) &= (0 - 5000) \times (10 + 2500/100) + \dots \\ &\quad + (0 - 10000) \times (10 + 5000/100) = 2050000 > 0 \end{aligned}$$

so we discard the upper half.

Bisection Iteration 2. Now we know $\lambda^* \in [0, 1/2]$. Calculate

$$\begin{aligned} \zeta'(1/4) &= (0 - 5000) \times (10 + 3750/100) + \dots \\ &\quad + (0 - 10000) \times (10 + 7500/100) = 550000 > 0 \end{aligned}$$

so we discard the upper half.

Bisection Iteration 3. Now we know $\lambda^* \in [0, 1/4]$. Calculate

$$\begin{aligned} \zeta'(1/8) &= (0 - 5000) \times (10 + 4375/100) + \dots \\ &\quad + (0 - 10000) \times (10 + 8750/100) = -200000 < 0 \end{aligned}$$

so we discard the lower half.

Bisection Iteration 4. Now we know $\lambda^* \in [1/8, 1/4]$. Calculate

$$\begin{aligned} \zeta'(3/16) &= (0 - 5000) \times (10 + 4062/100) + \dots \\ &\quad + (0 - 10000) \times (10 + 8125/100) = 175000 > 0 \end{aligned}$$

so we discard the upper half.

Bisection Iteration 5. Now we know $\lambda^* \in [1/8, 3/16]$. Calculate

$$\begin{aligned} \zeta'(5/32) &= (0 - 5000) \times (10 + 4219/100) + \dots \\ &\quad + (0 - 10000) \times (10 + 8437/100) = -12500 < 0 \end{aligned}$$

so we discard the lower half.

The final interval is $\lambda^* \in [5/32, 3/16]$, so the estimate is $\lambda^* = 11/64$ and

$$\mathbf{x} = \frac{11}{64}\mathbf{x}^* + \frac{53}{64}\mathbf{x} = [4141 \quad 859 \quad 2578 \quad 859 \quad 1719 \quad 1719 \quad 8281]$$

The new travel times are thus

$$\mathbf{t} = [51.4 \quad 18.6 \quad 35.8 \quad 18.6 \quad 27.2 \quad 27.2 \quad 92.8]$$

Iteration 2. With the new travel times, the shortest path for (1,3) is now [1, 3], but the shortest path for (2,4) is still [2, 5, 6, 4]. The relative gap is $AEC = 2.67$ (roughly 30 times smaller than the corresponding point in the MSA algorithm!) We have

$$\mathbf{x}^* = [5000 \quad 0 \quad 10000 \quad 0 \quad 10000 \quad 10000 \quad 0]$$

We begin the bisection method to find the right combination of \mathbf{x}^* and \mathbf{x} .

Bisection Iteration 1. Initially $\lambda^* \in [0, 1]$. Calculate

$$\zeta'(1/2) = 637329 > 0$$

so we discard the upper half.

Bisection Iteration 2. Now we know $\lambda^* \in [0, 1/2]$. Calculate

$$\zeta'(1/4) = 154266 > 0$$

so we discard the upper half.

Bisection Iteration 3. Now we know $\lambda^* \in [0, 1/4]$. Calculate

$$\zeta'(1/8) = 36063 > 0$$

so we discard the upper half.

Bisection Iteration 4. Now we know $\lambda^* \in [0, 1/8]$. Calculate

$$\zeta'(1/16) = 7741 > 0$$

so we discard the upper half.

Bisection Iteration 5. Now we know $\lambda^* \in [0, 1/16]$. Calculate

$$\zeta'(1/32) = 1302 > 0$$

so we discard the upper half.

The final interval is $\lambda^* \in [0, 1/32]$, so the estimate is $\lambda^* = 1/64$ and

$$\mathbf{x} = \frac{1}{64}\mathbf{x}^* + \frac{63}{64}\mathbf{x} = [4154 \quad 845 \quad 2694 \quad 845 \quad 1848 \quad 1848 \quad 8152]$$

The new travel times are thus

$$\mathbf{t} = [51.5 \quad 18.5 \quad 36.9 \quad 18.5 \quad 28.5 \quad 28.5 \quad 91.5]$$

At this point, the average excess cost is around 1.56 min; note that FW is able to decrease the relative gap much faster than MSA. However, we're still quite far from equilibrium if you compute the actual path travel times. In this case, even though we're allowing the step size to vary for each iteration, we are forcing travelers from all OD pairs to shift in the same proportion. In reality, OD pairs farther from equilibrium should see bigger flow shifts, and OD pairs closer to equilibrium should see smaller ones. This can be remedied by more advanced algorithms.

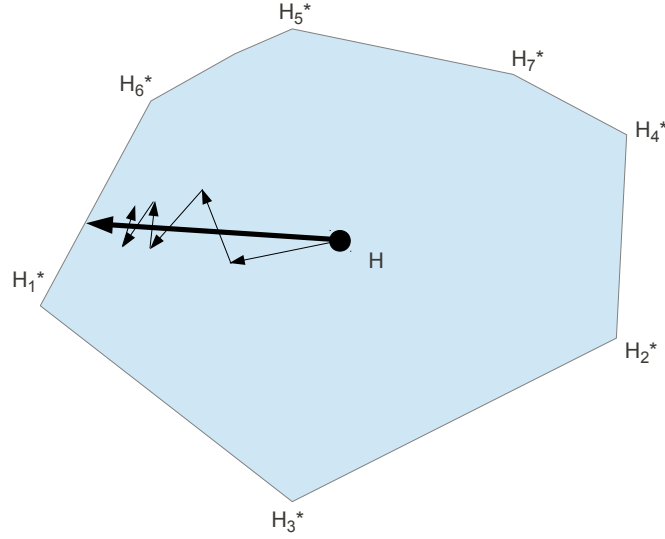


Figure 7.4: Frank-Wolfe can only move towards extreme points of the feasible region.

7.2.5 Conjugate Frank-Wolfe

If the distinction between FW and MSA is that FW chooses the step size λ in a more clever way, the distinction between conjugate Frank-Wolfe (CFW) and plain FW is that CFW chooses the target \mathbf{x}^* in a more clever way. To understand why CFW is more clever, we first need to understand why using the all-or-nothing assignment as target can be problematic.

Viewed in terms of the set of feasible link assignments X , the all-or-nothing assignments correspond to corner points of X . That is, Frank-Wolfe must limit itself to the corner points of the feasible region when determining where to move. In Figure 7.4, FW is constrained to follow the trajectory shown by the thin lines, and is unable to take a direct step like that indicated by the thick arrow. While these directions are effective in the early iterations, as the algorithm approaches the equilibrium point its converge slows down dramatically, and “zigzagging” behavior is observed.

So, how can we choose the target solution in a smarter way, so that steps in the direction of the target still move toward equilibrium, while granting more flexibility than the use of an all-or-nothing assignment? Conjugate Frank-Wolfe provides one approach towards doing so.

Understanding CFW requires introducing the concept of conjugacy, which is done here. Temporarily ignoring the context of the traffic assignment problem, assume that we are trying to find the minimum point of a convex quadratic function $f(x_1, x_2)$ of two variables, when there are no constraints. Figure 7.5 shows a few examples of these functions. Any quadratic function of two variables

can be written in the form

$$f(\mathbf{x}) = \mathbf{x}^T \mathbf{Q} \mathbf{x} + \mathbf{b}^T \mathbf{x}, \quad (7.17)$$

where $\mathbf{x} = [x_1 \ x_2]$ and \mathbf{b} are two-dimensional vectors, and \mathbf{Q} is a 2×2 matrix. Figure 7.5 shows the \mathbf{Q} matrix and \mathbf{b} vector corresponding to each example.

How would you go about finding the minimum of such a function? Given some initial solution (x_1, x_2) , one idea is to fix x_1 as a constant, and find the value of x_2 which minimizes f . Then, we can fix x_2 , and find the value of x_1 which minimizes f , and so on. This process will converge to the minimum, as shown in Figure 7.6, but in general this convergence is only asymptotic, and the process will never actually reach the minimum. The exception is when \mathbf{Q} is the identity matrix, as in Figure 7.6(a). In this case, the exact optimum is reached in only two steps.

In fact, it is possible to reach the exact optimum in only two steps even when \mathbf{Q} is not the identity matrix, by changing the search directions. The process described above (alternately fixing x_1 , and then x_2) can be thought of as alternating between searching in the direction $[0 \ 1]$, then searching in the direction $[1 \ 0]$. As shown in Figure 7.7, by making a different choice for the two search directions, the minimum can always be obtained in exactly two steps.

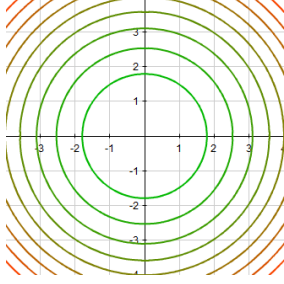
This happens if the two directions \mathbf{d}_1 and \mathbf{d}_2 are *conjugate*, that is, if

$$\mathbf{d}_1^T \mathbf{Q} \mathbf{d}_2 = 0 \quad (7.18)$$

Conjugacy generalizes the concept of orthogonality (or perpendicularity). If \mathbf{Q} is the identity matrix, equation (7.18) reduces to $\mathbf{d}_1 \cdot \mathbf{d}_2 = 0$, the definition of perpendicular vectors.

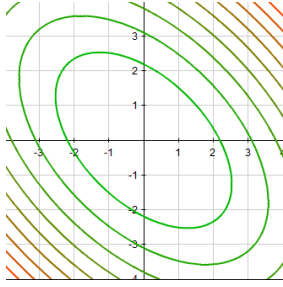
Now, returning to the traffic assignment problem, we want to use the concept of conjugacy to choose a more intelligent search direction. In particular, we want the target \mathbf{x}^* to be chosen so that the search direction $\mathbf{x}^* - \mathbf{x}$ is conjugate to the previous search direction. Before the derivation, there are a few differences between traffic assignment and the unconstrained quadratic program used to introduce conjugacy which should be addressed.

- The Beckmann function is not a function of two variables. This is not a huge problem. Finding the unconstrained minimum of a quadratic function of n variables requires n conjugate steps (so, in the above examples with two variables, two steps sufficed), due to a result known as the expanding subspace theorem. Of course...
- The Beckmann function is in general not a quadratic function. (Can you think of a case when it is?) Instead of the matrix \mathbf{Q} , we will instead use the Hessian of the Beckmann function Hf . Therefore, we cannot hope for exact convergence in a finite number of iterations. However, when the solution gets closer and closer to equilibrium, the Beckmann function can be better and better approximated by a quadratic by taking the first two terms of its Taylor series. This is good news, because zigzagging in



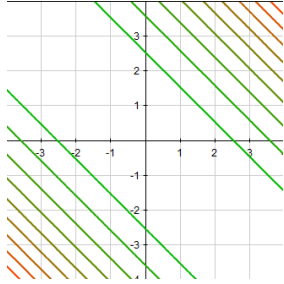
$$\mathbf{Q} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \quad \mathbf{b} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}$$

(a) $f(x_1, x_2) = x_1^2 + x_2^2$



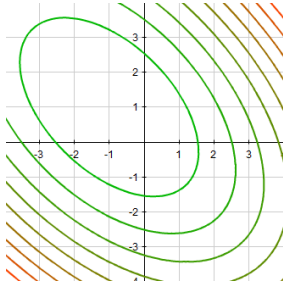
$$\mathbf{Q} = \begin{bmatrix} 1 & 1/2 \\ 1/2 & 1 \end{bmatrix} \quad \mathbf{b} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}$$

(b) $f(x_1, x_2) = x_1^2 + x_2^2 + x_1x_2$



$$\mathbf{Q} = \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} \quad \mathbf{b} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}$$

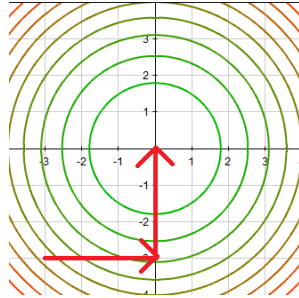
(c) $f(x_1, x_2) = x_1^2 + x_2^2 + 2x_1x_2$



$$\mathbf{Q} = \begin{bmatrix} 1 & 1/2 \\ 1/2 & 1 \end{bmatrix} \quad \mathbf{b} = \begin{bmatrix} 1 \\ -1 \end{bmatrix}$$

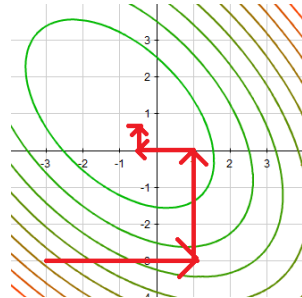
(d) $f(x_1, x_2) = x_1^2 + x_2^2 + x_1x_2 + x_1 - x_2$

Figure 7.5: Four examples of convex quadratic functions of the form $f(\mathbf{x}) = \mathbf{x}^T \mathbf{Q} \mathbf{x} + \mathbf{b}^T \mathbf{x}$



$$\mathbf{Q} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \quad \mathbf{b} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}$$

(a) $f(x_1, x_2) = x_1^2 + x_2^2$



$$\mathbf{Q} = \begin{bmatrix} 1 & 1/2 \\ 1/2 & 1 \end{bmatrix} \quad \mathbf{b} = \begin{bmatrix} 1 \\ -1 \end{bmatrix}$$

(b) $f(x_1, x_2) = x_1^2 + x_2^2 + x_1x_2 + x_1 - x_2$

Figure 7.6: Searching in orthogonal directions finds the optimum in two steps only if $\mathbf{Q} = \mathbf{I}$.

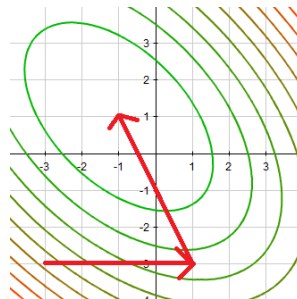


Figure 7.7: Searching in conjugate directions always leads to the optimum in two steps.

plain Frank-Wolfe becomes worse and worse as the equilibrium solution is approached.

- The optimization problem (7.3)–(7.6) has constraints. Again, this is not a major problem, since the feasible region is convex. As long as we ensure that the target point \mathbf{x}^* is feasible, any choice of $\lambda \in [0, 1]$ will retain feasibility.

So, how can we make sure that the new target vector \mathbf{x}^* is chosen so that the search direction is conjugate to the previous direction, and that \mathbf{x}^* is feasible? Since X is a convex set, feasibility can be assured by choosing \mathbf{x}^* to be a convex combination of the old target vector ($\mathbf{x}_{\text{old}}^*$) and the all-or-nothing assignment \mathbf{x}^{AON} :

$$\mathbf{x}^* = \alpha \mathbf{x}_{\text{old}}^* + (1 - \alpha) \mathbf{x}^{\text{AON}} \quad (7.19)$$

for some $\alpha \in [0, 1]$. Choosing $\alpha = 0$ would make the all-or-nothing assignment the target (as in plain FW), while choosing $\alpha = 1$ would make the target in this iteration the same as in the last. In fact, α should be chosen so that the new direction is conjugate to the last, that is

$$(\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\mathbf{x}^* - \mathbf{x}) = 0 \quad (7.20)$$

where \mathbf{H} is the Hessian of the Beckmann function evaluated at the current solution \mathbf{x} . Substituting (7.19) into (7.20), we can solve for α as follows:

$$(\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\mathbf{x}^* - \mathbf{x}) = 0 \quad (7.21)$$

$$\iff (\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\alpha \mathbf{x}_{\text{old}}^* + (1 - \alpha) \mathbf{x}^{\text{AON}} - \alpha \mathbf{x} - (1 - \alpha) \mathbf{x}) = 0 \quad (7.22)$$

$$\iff (\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\alpha(\mathbf{x}_{\text{old}}^* - \mathbf{x}) + (1 - \alpha)(\mathbf{x}^{\text{AON}} - \mathbf{x})) = 0 \quad (7.23)$$

or, after rearrangement,

$$\alpha[(\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\mathbf{x}_{\text{old}}^* - \mathbf{x}^{\text{AON}})] = -(\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\mathbf{x}^{\text{AON}} - \mathbf{x}) \quad (7.24)$$

$$\iff \alpha = \frac{(\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\mathbf{x}^{\text{AON}} - \mathbf{x})}{(\mathbf{x}_{\text{old}}^* - \mathbf{x})^T \mathbf{H}(\mathbf{x}_{\text{AON}}^* - \mathbf{x}_{\text{old}}^*)} \quad (7.25)$$

Now, for TAP, the Hessian takes a specific form. Since the Beckmann function is

$$f(\mathbf{x}) = \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (7.26)$$

its gradient is simply the vector of travel times at the current flows

$$\nabla f(\mathbf{x}) = \text{vect}\{t_{ij}(x_{ij})\} \quad (7.27)$$

and its Hessian is the diagonal matrix of travel time derivatives at the current flows

$$Hf(\mathbf{x}) = \text{diag}\{t'_{ij}(x_{ij})\} \quad (7.28)$$

So, the matrix products in equation (7.25) can be written out explicitly, giving

$$\alpha = \frac{\sum_{(i,j) \in A} ((x_{old}^*)_{ij} - x_{ij})(x_{ij}^{AON} - x_{ij})t'_{ij}}{\sum_{(i,j) \in A} ((x_{old}^*)_{ij} - x_{ij})(x_{ij}^{AON} - (x_{old}^*)_{ij})t'_{ij}} \quad (7.29)$$

where the derivatives t'_{ij} are evaluated at the current link flows x_{ij} .

Almost there! A careful reader may have some doubts about the formula in (7.29). First, it is possible that the denominator can be zero, and division by zero is undefined. Second, to ensure feasibility of \mathbf{x}^* , we need $\alpha \in [0, 1]$, even though it is not obvious that this formula always lies in this range (and in fact, it need not do so). Furthermore, $\alpha = 1$ is undesirable, because then the current target point is the same as the target point in the last iteration. If the previous line search was exact, there will be no further improvement and the algorithm will be stuck in an infinite loop. Finally, what should you do for the first iteration, when there is no “old” target \mathbf{x}^* ?

To address the first issue, the easiest approach is to simply set $\alpha = 0$ if the denominator of (7.29) is zero (i.e., if the formula is undefined, simply take a plain FW step by using the all-or-nothing solution as the target). As for the second and third issues, if the denominator is nonzero we can project the right-hand side of (7.29) onto the interval $[0, 1 - \epsilon]$ where $\epsilon > 0$ is some small tolerance value. That is, if equation (7.29) would give a value greater than $1 - \epsilon$, set $\alpha = 1 - \epsilon$; if it would give a negative value, use zero. Finally, for the first iteration, simply use the all-or-nothing solution as the target: $\mathbf{x}^* = \mathbf{x}^{AON}$.

So, to summarize the discussion, choose α in the following way. If it is the first iteration or the denominator of (7.29) is zero, set $\alpha = 0$. Otherwise set

$$\alpha = \text{proj}_{[0, 1 - \epsilon]} \left(\frac{\sum_{(i,j) \in A} t'_{ij}((x_{old}^*)_{ij} - x_{ij})(x_{ij}^{AON} - x_{ij})}{\sum_{(i,j) \in A} t'_{ij}((x_{old}^*)_{ij} - x_{ij})(x_{ij}^{AON} - (x_{old}^*)_{ij})} \right) \quad (7.30)$$

Then the target solution \mathbf{x}^* is calculated using (7.19). The value of the step size λ is chosen in the same way as in Frank-Wolfe, by performing a line search (e.g., using bisection or Newton’s method) to solve (7.16).

Large network example Here we apply CFW to the network shown in Figure 7.2, using the same notation as in the FW and MSA examples. The tolerance ϵ is chosen to be a small positive constant, 0.01 in the following example.

Initialization. Generate the initial solution by solving an all-or-nothing assignment. Path [1, 3] is shortest for OD pair (1,3), and path [2, 4] is shortest for OD pair (2,4), so

$$\mathbf{x}^{AON} = [5000 \quad 0 \quad 0 \quad 0 \quad 0 \quad 0 \quad 10000]$$

and

$$\mathbf{x} = \mathbf{x}^{AON} = [5000 \quad 0 \quad 0 \quad 0 \quad 0 \quad 0 \quad 10000]$$

Recalculating the travel times, we have

$$\mathbf{t} = [60 \quad 10 \quad 10 \quad 10 \quad 10 \quad 10 \quad 110]$$

Iteration 1. Proceeding in the same way as in the large network example for Frank-Wolfe, the all-or-nothing assignment in this case is

$$\mathbf{x}^{\text{AON}} = [0 \quad 5000 \quad 0 \quad 10000 \quad 15000 \quad 5000 \quad 10000]$$

which is used as the target \mathbf{x}^* since this is the first iteration of CFW. Repeating the same line search process, the optimal value of λ is 19/120, producing the new solution

$$\mathbf{x} = [4208 \quad 792 \quad 8417 \quad 1583 \quad 2375 \quad 792 \quad 1583]$$

Iteration 2. Again as with regular Frank-Wolfe, based on the updated travel times the all-or-nothing assignment is

$$\mathbf{x}^{\text{AON}} = [5000 \quad 0 \quad 0 \quad 10000 \quad 10000 \quad 0 \quad 10000]$$

However, here FW and CFW take different paths. Rather than using \mathbf{x}^{AON} as the target vector, CFW generates a conjugate search direction. First calculate the right-hand side of (7.29). Since for this problem $t_{ij} = 1/100$ for all links (regardless of the flow), the formula is especially easy to compute using \mathbf{x} and \mathbf{x}^* from the previous iteration and \mathbf{x}^{AON} as just now computed. The denominator of (7.29) is nonzero, and the formula gives -2.37 ; projecting onto the set $[0, 1 - \epsilon]$ thus gives $\alpha = 0$. So, calculating the target \mathbf{x}^* from equation (7.19) with $\alpha = 0$ we have

$$\mathbf{x}^* = [5000 \quad 0 \quad 0 \quad 10000 \quad 10000 \quad 0 \quad 10000]$$

and, using a line search between \mathbf{x} and \mathbf{x}^* , find that $\lambda = 0.0321$ is best, resulting in

$$\mathbf{x} = [4233 \quad 766 \quad 8146 \quad 1853 \quad 2620 \quad 766 \quad 1854]$$

Iteration 3. With the new flows \mathbf{x} , the travel times are now

$$\mathbf{t} = [52.3 \quad 17.7 \quad 91.5 \quad 28.5 \quad 36.2 \quad 17.7 \quad 28.5]$$

and the all-or-nothing assignment is

$$\mathbf{x}^{\text{AON}} = [5000 \quad 0 \quad 10000 \quad 0 \quad 0 \quad 0 \quad 0]$$

Calculating the right-hand side of (7.29), we see that the denominator is nonzero, and the formula gives 0.198, which can be used as is since it lies in $[0, 1 - \epsilon]$. So, calculating the target \mathbf{x}^* from equation (7.19) with $\alpha = 0.198$ we have

$$\mathbf{x}^* = [5000 \quad 0 \quad 8024 \quad 1976 \quad 1976 \quad 0 \quad 1976] .$$

Note that unlike any of the other link-based algorithms in this section, the target flows are not an all-or-nothing assignment (i.e., not an extreme point of X). Performing a line search between \mathbf{x} and \mathbf{x}^* , find that $\lambda = 0.652$ is best, resulting in

$$\mathbf{x} = [4733 \quad 267 \quad 8067 \quad 1933 \quad 2200 \quad 267 \quad 1933]$$

which solves the equilibrium problem exactly, so we terminate.

In this example, CFW found the exact equilibrium solution in three iterations. This type of performance is not typical (even though it is generally faster than regular FW or MSA). In this example, the link performance functions are linear, so the Beckmann function is quadratic. Iterative line searches with conjugate directions lead to the exact solution of quadratic programs in a finite number of iterations, as suggested by the above discussion. This performance cannot be assured with other types of link performance functions.

An even faster algorithm known as *biconjugate* Frank-Wolfe chooses its target so that the search direction is conjugate to both of the previous two search directions. This method converges better than CFW, but is not explained here because the details are a little more complicated even though the idea is the same.

7.3 Path-Based Algorithms

This section introduces equilibrium algorithms which work in the space of path flows H , rather than the space of link flows X . These tend to be faster, especially when high precision solutions are needed, but they require more computer memory. Furthermore, achieving the full potential of these algorithms for rapid convergence requires considerably more programming skill than for link-based algorithms.

Before explaining path-based algorithms, it is worth explaining why link-based algorithms are slow to converge. The following criticisms are specifically aimed at the Frank-Wolfe algorithm, but apply to other link-based algorithms as well.

It treats all OD pairs equally. If an OD pair is close to equilibrium, only a small flow shift among its paths is needed, while if an OD pair is far from equilibrium, a larger flow shift is needed. The Frank-Wolfe method uses the same λ shift for all OD pairs regardless of how close or far away each one is from equilibrium.

It uses a restricted set of search directions. If one imagines the space of feasible traffic assignments, the “all-or-nothing” \mathbf{x}^* solutions generated by the Frank-Wolfe algorithm represent extreme points or corners of this region. In other words, the Frank-Wolfe algorithm is only capable of moving towards a corner. Initially, this is fine, but as one nears the optimal solution, this results in extensive zig-zagging when a much more direct

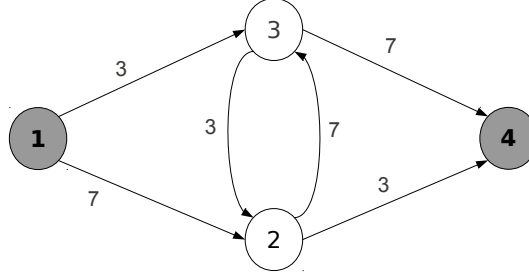


Figure 7.8: A persistent cycle in Frank-Wolfe. (Link flows shown.)

path exists (Figure 7.4). Conjugate Frank-Wolfe is a bit better in this regard, but still uses a fairly restrictive set of target points.

It is unable to erase cyclic flows. Consider the network in Figure 7.8, with the flows as shown. Such a flow might easily arise if $[1, 2, 3, 4]$ is the shortest path during the first iteration of Frank-Wolfe, $[1, 3, 2, 4]$ is the shortest path during the second, and $\lambda = 0.3$. With only one OD pair, it is impossible for both links $(2, 3)$ and $(3, 2)$ to be used at equilibrium, as discussed in Section 6.2.3. However, the Frank-Wolfe method will always leave some flow on both links unless $\lambda = 1$ at any iteration (which is exceedingly rare, especially in later iterations when λ is typically very close to zero).

These difficulties can all be avoided by tracking the path flows \mathbf{h} , rather than the link flows \mathbf{x} . The path flows contain much more information, tracking flow by origin and destination, as opposed to link flows which are aggregated together. On balance, the number of elements in the path flow vector is many orders of magnitude larger than that of the link flows, easily numbering in the millions for realistic networks. Algorithms which require us to first list off all paths in the network are not tractable. Instead, path-based algorithms only track the paths which an OD pair actually uses, that is the set $\hat{\Pi}^{rs} = \{\pi \in \pi^{rs} : h^\pi > 0\}$.⁴ This is often referred to as the set of *working paths* for each OD pair.⁵ A rough description of path-based algorithms can then be described as

1. Initialize $\hat{\Pi}^{rs} \leftarrow \emptyset$ for all OD pairs.
2. Repeat the following steps for each OD pair (r, s) :
 - (a) Find the shortest path $\hat{\pi}_{rs}$. Add it to $\hat{\Pi}^{rs}$ if it's not already used.

⁴Those familiar with other types of optimization problems might recognize this as a column generation scheme.

⁵You may recognize similarities with the “trial-and-error” method from Chapter 5. Path-based algorithms are essentially a more clever form of this method.

- (b) Shift travelers among paths to get closer to equilibrium.
 - (c) Update travel times.
3. Drop paths from $\hat{\Pi}^{rs}$ if they are no longer used; return to step 2 unless a convergence criterion is satisfied.

On the surface, this scheme looks quite similar to the link-based methods presented earlier. Why might it converge faster? Recall the three factors described above. First, each OD pair is now being treated independently. With MSA and Frank-Wolfe, the same step-size λ was applied across all links (and therefore, across all OD pairs). If one OD pair is very close to equilibrium, while another is far away, we should probably make a finer adjustment to the first OD pair, and a larger adjustment to the second one. Link-based methods allow no such finesse, and instead bluntly apply the same λ to all origins. In practice, this means that λ becomes very small after a few iterations: we can't move very far without disturbing an OD pair which is already close to equilibrium. As a result of this, it takes a really long time to solve OD pairs which are far from equilibrium. Equivalently, the set of search directions is broader in that we can vary the step size by OD pairs; and lastly, it is quite possible to erase cyclic flows in a path-based context, because the extra precision allows us to take larger steps.

Step 3 is where different path-based algorithms differ. This section describes two path-based algorithms: projected gradient and gradient projection (yes, these are different). The names are similar because both algorithms use the same two ingredients: exploiting the fact that the *gradient* is the direction of steepest ascent (and therefore, in a minimization problem, we should move in the opposite direction to descend as quickly as possible); and having to consider the constraints in the problem by using a *projection* operation to stay in the feasible set. (Recall from Chapter 3 that projection involves finding the point within a set which lies closest to another point.) Where they differ is in the order these steps are applied.

In the projected gradient method, we first calculate the gradient, then project the gradient onto the feasible set so that our search direction respects the demand constraint. We thus move along a “projected” version of the gradient.

In the gradient projection method, we first take a step in the opposite direction of the gradient, which will typically result in an infeasible point. We then apply a projection operation to return to the feasible set. The projection is done *after* the flow shift, not before: we do the projection after we make use of the gradient, so to speak.

7.3.1 Projected gradient

The projected gradient algorithm is based on the Beckmann formulation, applying an algorithm of Rosen from nonlinear optimization. The basic idea is as follows: at any point, the gradient of a function is a vector pointing in the direction of steepest increase. Since equilibria correspond to minimum points,

moving in the direction of the negative gradient will travel along the direction of steepest descent. However, we need to be careful not to leave the region of feasible path flows H , making sure that our search direction still satisfies the demand constraint and retains nonnegative path flows.

Writing the Beckmann function in terms of path flows \mathbf{h} , rather than link flows as is customary, we have

$$f(\mathbf{h}) = \sum_{(i,j) \in A} \int_0^{\sum_{\pi \in \Pi} \delta_{ij}^{\pi} h^{\pi}} t_{ij}(x) dx \quad (7.31)$$

and its partial derivative with respect to any path flow variable is

$$\frac{\partial f}{\partial h^{\pi}} = \sum_{(i,j) \in A} \delta_{ij}^{\pi} t_{ij} \left(\sum_{\pi' \in \Pi} \delta_{ij}^{\pi'} h^{\pi'} \right) = \sum_{(i,j) \in A} \delta_{ij}^{\pi} t_{ij}(x_{ij}) = c^{\pi} \quad (7.32)$$

so the direction of steepest descent \mathbf{s} is the negative gradient:

$$\mathbf{s} = -\nabla_{\mathbf{h}} f = -\text{vect}(c^{\pi}) \quad (7.33)$$

Assuming that the current path flow solution \mathbf{h} is feasible, we must move in a direction that does not violate any of the constraints, only using the working paths $\hat{\Pi}^{rs}$. For instance, if the demand constraint $\sum_{\pi \in \hat{\Pi}^{rs}} h^{\pi} = d^{rs}$ is satisfied for all OD pairs (r, s) , it must remain so after taking a step in the direction $\Delta \mathbf{h}$:

$$\sum_{\pi \in \hat{\Pi}^{rs}} (h^{\pi} + \Delta h^{\pi}) = d^{rs} \quad (7.34)$$

This in turn implies

$$\sum_{\pi \in \hat{\Pi}^{rs}} \Delta h^{\pi} = d^{rs} - \sum_{\pi \in \hat{\Pi}^{rs}} h^{\pi} \quad (7.35)$$

$$\sum_{\pi \in \hat{\Pi}^{rs}} \Delta h^{\pi} = 0 \quad (7.36)$$

So, we need to find the projection of the steepest descent direction \mathbf{s} onto the space $\Delta H \equiv \{\Delta \mathbf{h} : \sum_{\pi \in \hat{\Pi}^{rs}} \Delta h^{\pi} = 0\}$. It turns out that this projection has a remarkably simple closed form expression.

Proposition 7.1. *Using $\bar{c}^{rs} = (1/|\hat{\Pi}^{rs}|) \sum_{\pi \in \hat{\Pi}^{rs}} c^{\pi}$ to denote the average travel time of the working paths for OD pair (r, s) , the direction \mathbf{s}' with components $s'_{\pi} = c^{\pi} - \bar{c}^{rs}$ is the projection of the steepest descent direction \mathbf{s} onto the set ΔH .*

Proof. To show that \mathbf{s}' is the projection of \mathbf{s} onto ΔH , we must show that $\mathbf{s}' \in \Delta H$, and that $\mathbf{s} - \mathbf{s}'$ is orthogonal to ΔH . Regarding the first part, we have

$$\sum_{\pi \in \hat{\Pi}^{rs}} s'_{\pi} = \sum_{\pi \in \hat{\Pi}^{rs}} (c^{\pi} - \bar{c}^{rs}) = \sum_{\pi \in \hat{\Pi}^{rs}} c^{\pi} - |\hat{\Pi}^{rs}| \bar{c}^{rs} = 0 \quad (7.37)$$

so $\mathbf{s}' \in \Delta H$.

Regarding the second part, let $\Delta \mathbf{h}$ be any vector in ΔH . We now show that $(\mathbf{s} - \mathbf{s}') \cdot \Delta \mathbf{h} = 0$. We have

$$(\mathbf{s} - \mathbf{s}') \cdot \Delta \mathbf{h} = \sum_{\pi \in \hat{\Pi}^{rs}} [c^\pi - (c^\pi - \bar{c}^{rs})] \Delta h^\pi \quad (7.38)$$

$$= -\bar{c}^{rs} \sum_{\pi \in \hat{\Pi}^{rs}} \Delta h^\pi \quad (7.39)$$

$$= 0 \quad (7.40)$$

since $\Delta \mathbf{h} \in \Delta H$. \square

So, given current path flows \mathbf{h}^{rs} for OD pair (r, s) , we use $\Delta \mathbf{h}^{rs} = \text{vect}(c^\pi - \bar{c}^{rs})$ as the search direction. To update the path flows $\mathbf{h}^{rs} \leftarrow \mathbf{h}^{rs} + \mu \Delta \mathbf{h}^{rs}$, we need an expression for the step size. For any path $\pi \in \hat{\Pi}^{rs}$ for which $\Delta h^\pi < 0$, the new path flows would be infeasible if $\mu > h^\pi / \Delta h^\pi$. Therefore, the largest possible step size is

$$\bar{\mu} = \min_{\pi \in \hat{\Pi}^{rs}; \Delta h^\pi < 0} \frac{h^\pi}{\Delta h^\pi} \quad (7.41)$$

The actual step size $\mu \in [0, \bar{\mu}]$ should be chosen to minimize the Beckmann function. This can be done either through bisection or one or more iterations of Newton's method.

7.3.2 Gradient projection

The gradient projection algorithm provides an alternative way to generate $\Delta \mathbf{h}$. The names projected gradient and gradient projection are very similar, but do not be confused. In the projected gradient algorithm, we first calculated a gradient, then projected the gradient onto the space of feasible directions, then took a step in the opposite direction of the projected gradient. In gradient projection, the last two steps are reversed: we calculate the gradient, take a step in the opposite direction, *then* project the solution onto the set of feasible path flows. Projected gradient methods search the interior of the feasible region; gradient projection can search along the boundary as well as the interior.

For an OD pair (r, s) and a given set of working paths Π^{rs} , the gradient of the Beckmann function was found to be

$$\nabla_{\mathbf{h}} f = \text{vect}(c^\pi) \quad (7.42)$$

So, a first attempt at gradient projection would be to update $\mathbf{h} \leftarrow \text{proj}_H(\mathbf{h} - \nabla_{\mathbf{h}} f)$. Unfortunately, projecting onto the set H is not particularly easy. If we apply a suitable change of variables, though, the projection can be made much easier.

Define the *basic path* for OD pair (r, s) to be a path $\hat{\pi}_{rs}$ with minimum travel time. All of the other paths are called *nonbasic paths*. We can eliminate the

basic path flow variable by expressing it in terms of the nonbasic path flows:

$$h^{\hat{\pi}_{rs}} = d^{rs} - \sum_{\pi \in \hat{\Pi}^{rs} : \pi \neq \hat{\pi}_{rs}} h^{\pi}. \quad (7.43)$$

Substituting this into the path-based Beckmann function (7.31), the partial derivative with respect to one of the nonbasic path flows is now

$$\frac{\partial \hat{f}}{\partial h^{\pi}} = c^{\pi} - c^{\hat{\pi}_{rs}} \quad \forall (r, s) \in Z^2, \pi \in \hat{\Pi}^{rs} - \{\hat{\pi}_{rs}\}, \quad (7.44)$$

denoting the Beckmann function as \hat{f} instead of f because the function has been modified by using (7.43) to eliminate some of the path flow variables.

So, the change in path flows will be the negative of the gradient. Since the gradient is given by (7.44) and $c^{\pi} \geq c^{\hat{\pi}_{rs}}$ (because the basic path is by definition the shortest one), moving in this direction means that every nonbasic path flow will *decrease*. Since the transformation (7.43) eliminated the demand satisfaction constraint, the only remaining constraint is that the nonbasic path flow variables be nonnegative. Projecting onto this set is trivial: if any of the nonbasic path flow variables is negative after taking a step, simply set it to zero. At this point, the basic path flow can be calculated through equation (7.43).

Furthermore, the larger the difference in travel times, the larger the corresponding element of the gradient will be. This suggests that more flow be shifted away from paths with higher travel times. We can go a step further, and estimate directly how much flow should be shifted from a nonbasic path to a basic path to equalize the travel times, using Newton's method.

Let Δh denote the amount of flow we shift *away* from non-basic path π and *onto* the basic path $\hat{\pi}$, and let $c_{\pi}(\Delta h)$ and $c_{\hat{\pi}}(\Delta h)$ denote the travel times on path π and $\hat{\pi}$ after we make such a shift. We want to choose Δh so these costs are equal, that is, so

$$g(\Delta h) = c_{\pi}(\Delta h) - c_{\hat{\pi}}(\Delta h) = 0; \quad (7.45)$$

that is, g is simply the difference in travel times between the two paths. To apply Newton's method, we need to find the derivative of g with respect to Δh .

Using the relationships between link travel times and path travel times, we have

$$g(\Delta h) = \sum_{(i,j) \in \mathcal{A}} (\delta_{ij}^{\pi} - \delta_{ij}^{\hat{\pi}}) t_{ij}(x_{ij}(\Delta h))$$

so

$$g'(\Delta h) = \sum_{(i,j) \in \mathcal{A}} (\delta_{ij}^{\pi} - \delta_{ij}^{\hat{\pi}}) \frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{d\Delta h}$$

by the chain rule. For each arc, there are four possible cases:

Case I : $\delta_{ij}^{\pi} = \delta_{ij}^{\hat{\pi}} = 0$, that is, neither path π nor path $\hat{\pi}$ uses link (i, j) . Then $(\delta_{ij}^{\pi} - \delta_{ij}^{\hat{\pi}}) \frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{d\Delta h} = 0$ and this link does not contribute to the derivative.

Case II : $\delta_{ij}^\pi = \delta_{ij}^{\hat{\pi}} = 1$, that is, both paths π and path $\hat{\pi}$ use link (i, j) .

Then $(\delta_{ij}^\pi - \delta_{ij}^{\hat{\pi}}) \frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{d\Delta h} = 0$ and this link again does not contribute to the derivative. (Another way to think of it: since both paths use this arc, its total flow will not change if we shift travelers from one path to another.)

Case III : $\delta_{ij}^\pi = 1$ and $\delta_{ij}^{\hat{\pi}} = 0$, that is, path π uses link (i, j) , but path $\hat{\pi}$ does not. Then $(\delta_{ij}^\pi - \delta_{ij}^{\hat{\pi}}) \frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{\Delta h} = \frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{d\Delta h} = -\frac{dt_{ij}}{dx_{ij}}$ since $\frac{dx_{ij}}{d\Delta h} = -1$.

Case IV : $\delta_{ij}^\pi = 0$ and $\delta_{ij}^{\hat{\pi}} = 1$, that is, path $\hat{\pi}$ uses link (i, j) , but path π does not. Then $(\delta_{ij}^\pi - \delta_{ij}^{\hat{\pi}}) \frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{\Delta h} = -\frac{dt_{ij}}{dx_{ij}} \frac{dx_{ij}}{d\Delta h} = -\frac{dt_{ij}}{dx_{ij}}$ since $\frac{dx_{ij}}{d\Delta h} = 1$.

Putting it all together, the only terms which contribute to the derivative are the links which are used by either π or $\hat{\pi}$, but not both. Let A_1 , A_2 , A_3 , and A_4 denote the sets of links falling into the four cases listed above. Then

$$g'(\Delta h) = - \sum_{(i,j) \in A_3 \cup A_4} \frac{dt_{ij}}{dx_{ij}}$$

which is simply the negative sum of the derivatives of these links, evaluated at the current link flows.

Then, starting with an initial guess of $\Delta h = 0$, one step of Newton's method gives an improved guess of

$$\Delta h = 0 - g(0)/g'(0) = \frac{c_\pi - c_{\hat{\pi}}}{\sum_{a \in A_3 \cup A_4} \frac{dt_{ij}}{dx_{ij}}}$$

That is, the recommended Newton shift is given by the difference in path costs, divided by the sum of the derivatives of the link performance functions for links used by one path or the other, but not both. Therefore, the updated nonbasic and basic path flows are given by

$$h_{\hat{\pi}} \leftarrow h_{\hat{\pi}} + \frac{c_\pi - c_{\hat{\pi}}}{\sum_{a \in A_3 \cup A_4} \frac{dt_{ij}}{dx_{ij}}}$$

and

$$h_\pi \leftarrow h_\pi - \frac{c_\pi - c_{\hat{\pi}}}{\sum_{a \in A_3 \cup A_4} \frac{dt_{ij}}{dx_{ij}}}$$

This process is repeated for every nonbasic path.

This is demonstrated on the example in Figure 7.2 as follows.

Iteration 1, Step 1. Initially $\hat{\Pi}^{13} = \hat{\Pi}^{24} = \emptyset$.

Iteration 1, Step 2a, OD pair (1,3). Find the shortest path for (1,3): with no travelers on the network, the top link has a travel time of 10. This is not in the set of used paths, so include it: $\hat{\Pi}^{13} = \{[1, 3]\}$.

Iteration 1, Step 2b, OD pair (1,3). Since there is only one used path, we simply have $h_{[1,3]}^{13} = 5000$.

Iteration 1, Step 2c, OD pair (1,3). Update travel times:

$$\mathbf{t} = [60 \quad 10 \quad 10 \quad 10 \quad 10 \quad 10 \quad 10]$$

Iteration 1, Step 2a, OD pair (2,4). Find the shortest path for (2,4): with no travelers on the network, link 7 has a travel time of 10. This is not in the set of used paths, so include it: $\hat{\Pi}^{24} = \{[2, 4]\}$.

Iteration 1, Step 2b. Since there is only one used path, $h_{[2,4]}^{24} = 10000$.

Iteration 1, Step 2c. Update travel times:

$$\mathbf{t} = [60 \quad 10 \quad 10 \quad 10 \quad 10 \quad 10 \quad 110]$$

Iteration 1, Step 3. All paths are used, so return to step 2. The relative gap is $\gamma_1 = 2.11$.

Iteration 2, Step 2a, OD pair (1,3). The shortest path is now $[1, 5, 6, 3]$. This is not part of the set of used paths, so we add it: $\hat{\Pi}^{13} = \{[1, 3], [1, 5, 6, 3]\}$.

Iteration 2, Step 2b, OD pair (1,3). The difference in travel times between the paths is 30 minutes; and the sum of the derivatives of links 1, 2, 3, and 4 is 0.04. So we shift $30/0.04 = 750$ vehicles from $[1, 3]$ to $[1, 5, 6, 3]$, $h_{[1,3]}^{13} = 4250$ and $h_{[1,5,6,3]}^{13} = 750$.

Iteration 2, Step 2c, OD pair (1,3). Update travel times:

$$\mathbf{t} = [52.5 \quad 17.5 \quad 17.5 \quad 17.5 \quad 10 \quad 10 \quad 110]$$

Note that the two paths have exactly the same cost after only one step! This is because Newton's method is exact for linear functions.

Iteration 2, Step 2a, OD pair (2,4). The shortest path is now $[2, 5, 6, 4]$. This is not part of the set of used paths, so we add it: $\hat{\Pi}^{24} = \{[2, 4], [2, 5, 6, 4]\}$.

Iteration 2, Step 2b, OD pair (2,4). The difference in travel times between the paths is 72.5 minutes; and the sum of the derivatives of links 5, 3, 6, and 7 is 0.04. So we shift $72.5/0.04 = 1812.5$ vehicles from $[2, 4]$ to $[2, 5, 6, 4]$, $h_{[2,4]}^{24} = 8187.5$ and $h_{[2,5,6,4]}^{24} = 1812.5$.

Iteration 2, Step 2c, OD pair (2,4). Update travel times:

$$\mathbf{t} = [52.5 \quad 17.5 \quad 35.625 \quad 17.5 \quad 28.125 \quad 28.125 \quad 91.875]$$

Note that the two paths again have exactly the same cost. However, the equilibrium for the first OD pair has been disturbed.

Iteration 2, Step 3. All paths are used, so return to step 2. The relative gap is $\gamma_1 = 0.0115$.

Iteration 3, Step 2a, OD pair (1,3). The shortest path is now [1,3], which is already in the set of used paths, so nothing to do here.

Iteration 3, Step 2b, OD pair (1,3). The difference in travel times between the paths is 18.125 minutes; and the sum of the derivatives of links 1, 2, 3, and 4 is 0.04. So we shift $18.125/0.04 = 453$ vehicles from [1,5,6,3] to [1,3], $h_{[1,3]}^{13} = 4703$ and $h_{[1,5,6,3]}^{13} = 297$.

Iteration 3, Step 2c, OD pair (1,3). Update travel times:

$$\mathbf{t} = [57.0 \quad 13.0 \quad 31.1 \quad 13.0 \quad 28.1 \quad 28.1 \quad 91.9]$$

Again the first OD pair is at equilibrium, up to rounding error.

Iteration 3, Step 2a, OD pair (2,4). The shortest path is again [2,5,6,4]. This is already in the set of used paths, so nothing to do here.

Iteration 3, Step 2b, OD pair (2,4). The difference in travel times between the paths is 4.6 minutes; and the sum of the derivatives of links 5, 3, 6, and 7 is 0.04. So we shift $4.6/0.04 = 114$ vehicles from [2,4] to [2,5,6,4], $h_{[2,4]}^{24} = 8073$ and $h_{[2,5,6,4]}^{24} = 1927$.

Iteration 3, Step 2c, OD pair (2,4). Update travel times:

$$\mathbf{t} = [57.0 \quad 13.0 \quad 32.2 \quad 13.0 \quad 29.3 \quad 29.3 \quad 90.7]$$

Iteration 3, Step 3. All paths are used, so return to step 2. The relative gap is $\gamma_1 = 0.00028$.

Note that after three iterations of gradient projection, the relative gap is two orders of magnitude smaller than that from the Frank-Wolfe algorithm. Although not demonstrated here for reasons of space, the performance of gradient projection relative to Frank-Wolfe actually *improves* from here on out. Frank-Wolfe usually does most of its work in the first few iterations, and then converges *very* slowly after that.⁶ On the other hand, gradient projection maintains a steady rate of progress throughout, with a nearly constant proportionate decrease in gap from one iteration to the next.

7.4 Bush-Based Algorithms

As seen in the previous section, path-based algorithms converge faster than link-based algorithms, and allow us find much more accurate equilibrium solutions in a fraction of the time. The prime disadvantage is a huge memory requirement, with potentially millions of paths available for use in large networks. Furthermore, many of these paths are “redundant” in some way, as they overlap: the same sequence of links might be used by many different paths.

⁶It’s been said that Frank-Wolfe converges, but just barely.

Bush-based algorithms (also known as origin-based algorithms) try to address these limitations. Rather than treating each OD pair separately (as link-based algorithms do), bush-based algorithms simultaneously consider every destination associated with a single origin. Instead of considering every possible used path (as path-based algorithms do, and there are potentially very many of these), they maintain a set of links called a *bush*, which are the only links which travelers from that origin are permitted to use. One can think of a bush as the set of links obtained by superimposing all of the paths used by travelers starting from an origin. In particular, a bush must be:

- *Connected*; that is, using only links in the bush, it is possible to reach every node which was reachable in the original network.
- *Acyclic*; that is, no path using only bush links can pass the same node more than once. This is not restrictive, because travelers trying to minimize their own travel time would never cycle back to the same node, and greatly speeds up the algorithm, because acyclic networks are much simpler and admit much faster methods for finding shortest paths and other quantities of interest.

There is no particular reason why bushes have to be “origin-based” rather than “destination-based,” and all of the results in this section can be derived in a parallel way for bushes terminating at a common destination, rather than starting at a common origin.

An example of a bush is shown in panels (a) and (b) of Figure 7.9, where the thickly-shaded links are part of the bush, and the lightly-shaded links are not. The bush in panel (a) is a special type of bush known as a *tree*, which has exactly one path from the origin to every node. The thickly-shaded links in panel (c) do not form a bush, because it is not connected; there is no way to reach the nodes at the bottom of the network only using bush links. Likewise, the thick links in panel (d) do not form a bush either, because a cycle exists and it would be possible to revisit some nodes multiple times using the bush links (find them!).

Notice that because a bush is acyclic, it can never include both directions of a two-way link. This implies that at equilibrium, *on every street travelers from the same origin must all be traveling in the same direction*. (This follows from Proposition 6.3 in Section 6.2.3). Interestingly, link-based and path-based algorithms cannot enforce this requirement easily; and this is yet another reason that bush-based algorithms are a good option for solving equilibrium. There is one bush for every origin; this means that if there are z origins and m links, we need to keep track of at most zm values. By contrast, a link-based approach (such as Frank-Wolfe) requires storage of only m values to represent a solution, while a path-based approach could conceivably require $z^2 2^m$ values.⁷

All bush-based algorithms operate according to the same general scheme:

1. Start with initial bushes for each origin (the shortest path tree with free-flow times is often used as a starting point).

⁷These values are very approximate, but give you an idea of the scale.

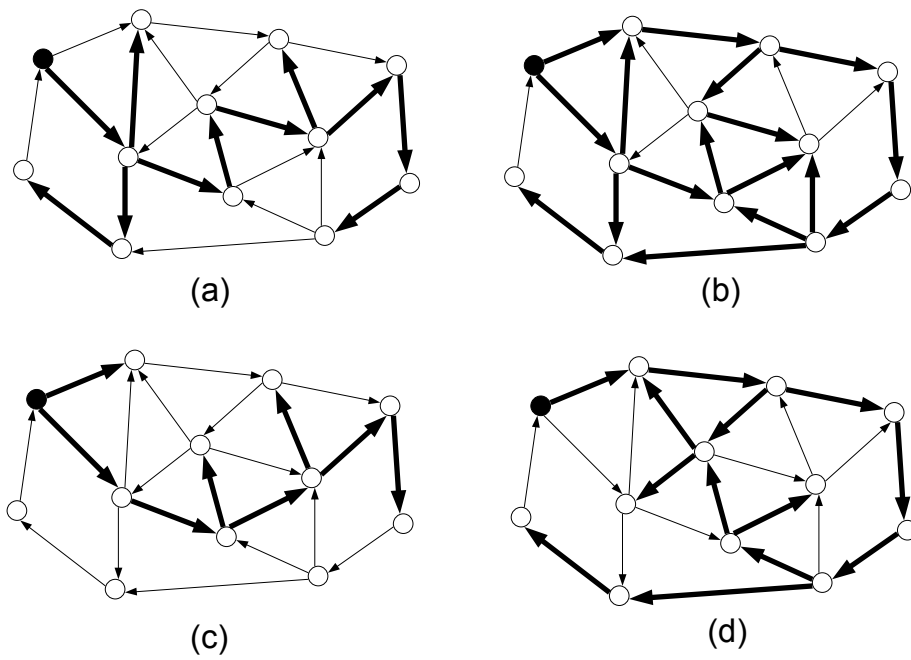


Figure 7.9: Examples of bushes (panels (a) and (b)) and non-bushes (panels (c) and (d)).

2. Shift flows within each bush to bring each origin closer to equilibrium.
3. Improve the bushes by adding links which can reduce travel times, and by removing unused links. Return to step 2.

Step 1 is fairly self-explanatory: collecting the shortest paths from an origin to every destination into one bush results in a connected and acyclic set of links, and we can simply load the travel demand to each destination on the unique path in the bush.

Step 2 is where most bush-based algorithms differ.⁸ This section explains three different methods for shifting flows within a bush: Algorithm B, origin-based assignment (OBA), and linear user cost equilibrium (LUCE). Each of these algorithms is based on labels calculated for bush nodes and links, and Section 7.4.1 defines each of these. The three algorithms themselves are presented in Section 7.4.2.

Step 3 requires a little bit of thought to determine how to adjust the bush links themselves to allow movement towards equilibrium when Step 2 is performed on the new bushes. Section 7.4.3 shows how this can be done.

7.4.1 Bush labels

Algorithm B, OBA, and LUCE all make use of “bush labels” associated with each node and link in a bush. All of these labels can be calculated in a straightforward, efficient manner using the topological order within each bush. As you read this section, pay attention to which of these labels are calculated in ascending topological order (i.e., the formulas only involve nodes with lower topological order) and which are calculated in descending order (based only on nodes with higher topological order.) In the discussion below, assume that a bush is given and fixed, with origin node r . The set of bush links is denoted \mathcal{B} and, for the purposes of this section, assume that the travel times t_{ij} and travel time derivatives t'_{ij} of all bush links are given and constant. All of these labels are only defined for bush links, and the formulas only involve bush links. From the standpoint of Step 2 of bush-based algorithms, non-bush links are completely irrelevant. Table 7.1 shows all of the labels defined in this section, and which labels are used in which algorithms, and in the bush-updating steps described in Section 7.4.3.

We start with two different ways to represent the travel patterns on each bush, starting with x labels. The label $x_{ij}^{\mathcal{B}}$ associated with each link indicates the number of travelers starting at node r (the root of bush \mathcal{B}) and traveling on bush link (i, j) . The superscript \mathcal{B} indicates that we are only referring to the flow on this link associated with the bush \mathcal{B} , and the *total* link flow x_{ij} is the sum of $x_{ij}^{\mathcal{B}}$ across all bushes \mathcal{B} . However, using these superscripts tends to clutter formulas, and often times it is clear that we are only referring to flows within the context of a specific bush. In this case, we can simply write x_{ij} with

⁸If this is starting to sound familiar, good! How to shift flows to move closer to equilibrium is also where link-based and path-based algorithms differ most.

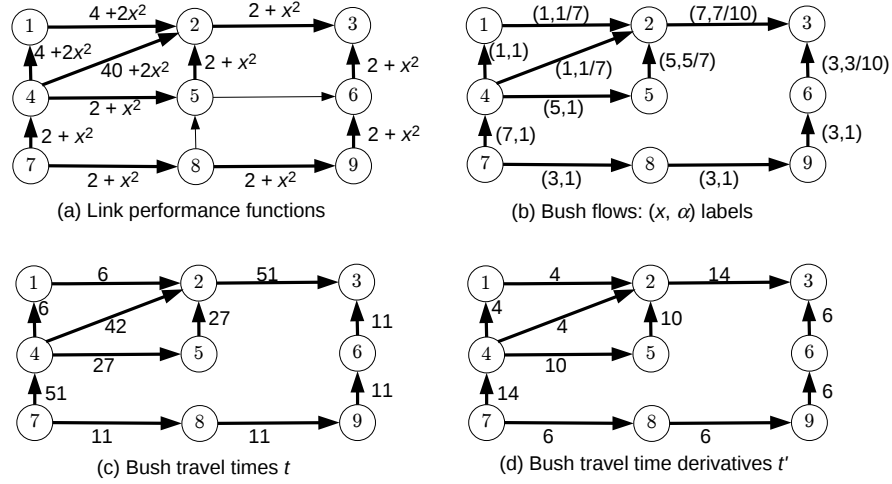


Figure 7.10: Example bush used to demonstrate label calculation.

it being understood that this label refers to the flow on a bush link. Within this section, we are only concerned with a single bush and the superscript will be omitted for brevity.

The network in Figure 7.10(a) will be used as to demonstrate the labels introduced in this section. The thick links comprise the bush, and the link performance functions for all links in the network are shown. The origin in this case is node 7, and the demand is 10 vehicles from node 7 to node 3. You can verify that a topological ordering of the nodes on this bush is 7, 8, 9, 6, 4, 1, 5, 2, 3.

The corresponding label x_i associated with each node indicates the total number of vehicles using node i on the bush \mathcal{B} , including flow which is terminating at node i . The flow conservation equations relating x_{ij} and x_i labels are as follows:

$$x_i = \sum_{(h,i) \in \mathcal{B}} x_{hi} = d^{ri} + \sum_{(i,j) \in \mathcal{B}} x_{ij} \quad \forall i \in N \quad (7.46)$$

where the first expression defines the node flow in terms of *incoming* link flows, and the second in terms of *outgoing* link flows. The two definitions are equivalent.

It is sometimes convenient to refer to the *fraction* of the flow at a node coming from a particular link. For any node i with positive node flow ($x_i > 0$), define α_{hi} to be the proportion of the node flow contributed by the incoming link (h, i) , that is

$$\alpha_{hi} = x_{hi} / x_i \quad (7.47)$$

Clearly each α_{hi} is nonnegative, and, by flow conservation, the sum of the α_{hi} values entering each node i is one. The definition of α_{hi} is slightly trickier

when $x_i = 0$, because the formula (7.47) then involves a division by zero. To accommodate this case, we adopt this rule: *when $x_i = 0$, the proportions α_{hi} may take any values whatsoever, as long as they are nonnegative and sum to one.* It is important to be able to define α_{hi} values even in this case, because the flow-shifting algorithms may cause x_i to become positive, and in this case we need to know how to distribute this new flow among the incoming links.

If we are given α labels for each bush link, it is possible to calculate the resulting node and link flows x_i and x_{ij} , using this recursion:

$$x_i = d^r + \sum_{(i,j) \in \mathcal{B}} x_{ij} \quad \forall i \in N \quad (7.48)$$

$$x_{ij} = \alpha_{ij} x_j \quad \forall (i,j) \in \mathcal{B} \quad (7.49)$$

The sum in (7.48) is empty for the node with the highest topological order. So we can start there, and then proceed with the calculations in backward topological order.

Figure 7.10(b) shows the x and α labels for the example bush. You should verify that the formulas (7.48) and (7.49) are consistent with these labels. The link travel times and link travel time derivatives are shown in panels (c) and (d) of this figure.

As has been used earlier in the text for shortest path algorithms, L is used to denote travel times on shortest paths. The superscript \mathcal{B} can be appended to these labels when it is necessary to indicate that these labels are for shortest paths *specifically on the bush \mathcal{B}* , although it will usually be clear from context which bush is meant. This section will omit such a superscript to avoid cluttering the formulas, and it should be understood that L means the shortest path *only on the bush under consideration*. The same convention will apply to the other labels in this section. There are L labels associated with each node i , and with each link $(i,j) \in \mathcal{B}$: L_i denotes the distance on the shortest path from r to i using only bush links, and L_{ij} is the travel time which would result if you follow the shortest path to node i , then take link (i,j) . These labels are calculated using these equations:

$$L_r = 0 \quad (7.50)$$

$$L_{ij} = L_i + t_{ij} \quad \forall (i,j) \in \mathcal{B} \quad (7.51)$$

$$L_i = \min_{(h,i) \in \mathcal{B}} \{L_{hi}\} \quad \forall i \neq r \quad (7.52)$$

The U labels are used to denote travel times on longest paths within the bush, and are calculated in a similar way. Like the L labels, U labels are calculated for both nodes and bush links, using the formulas:

$$U_r = 0 \quad (7.53)$$

$$U_{ij} = U_i + t_{ij} \quad \forall (i,j) \in \mathcal{B} \quad (7.54)$$

$$U_i = \max_{(h,i) \in \mathcal{B}} \{U_{hi}\} \quad \forall i \neq r \quad (7.55)$$

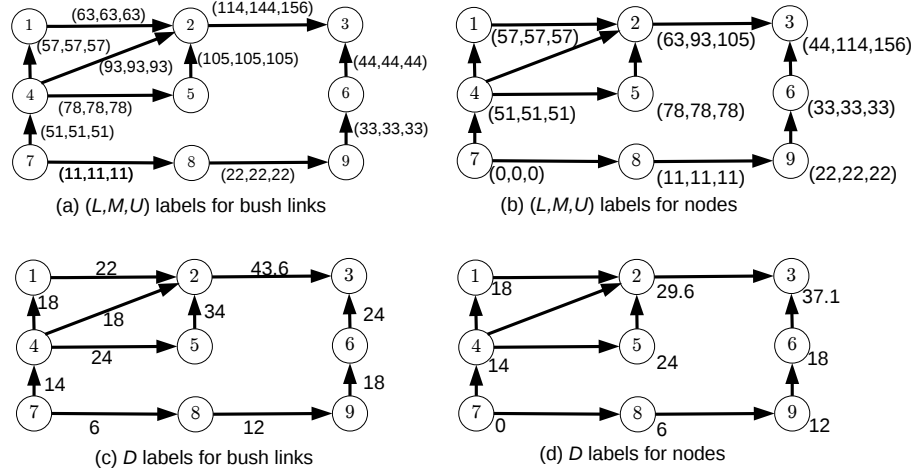


Figure 7.11: Continuation of label calculation demonstration.

The M labels are used to denote the *average* travel times within the bush, recognizing that some travelers will be on longer paths and other travelers will be on shorter ones. The node label M_i represents the average travel time between origin r and node i across all bush paths connecting these nodes, weighted by the number of travelers using each of these paths. The label M_{ij} indicates the average travel time of vehicles after they finish traveling on link (i, j) , again averaging across all of the bush paths starting at the origin r and ending with link (i, j) . These can be calculated as follows:

$$M_r = 0 \quad (7.56)$$

$$M_{ij} = M_i + t_{ij} \quad \forall (i, j) \in \mathcal{B} \quad (7.57)$$

$$M_i = \sum_{(h,i) \in \mathcal{B}} \alpha_{hi} M_{hi} \quad \forall i \neq r \quad (7.58)$$

Figure 7.11 shows the L , U , and M labels corresponding to the example bush in panels (a) and (b). Panel (a) shows the labels associated with links (L_{ij} , M_{ij} , and U_{ij}), while panel (b) shows the labels associated with nodes (L_i , M_i , and U_i).

It is also useful to know how the average travel times (the M labels) will change with a marginal increase in flow on a particular link or through a particular node. The set of D labels are used to represent this. As we will see shortly, these play the role of the travel time derivatives⁹, which are needed in flow shifting rules based on Newton's method (like those described for path-based algorithms in the previous section.) Ideally, the D_{ij} labels would represent the

⁹For a mnemonic, you can associate the labels L , U , and M with the *lower*, *upper*, and *mean* travel times to nodes, while D refers to the *derivative*.

Table 7.1: Bush labels used in different algorithms.

Label	B	OBA	LUCE	Updating
x	X			
α		X	X	
L	X			X
U	X			X
M		X	X	
D		X	X	

derivative of M_j with respect to α_{ij} . However, there is no known method for calculating these derivatives in an efficient way (“efficient” meaning they can be calculated in a single topological pass). The following formulas for D_{ij} can be used as an approximation that can be calculated in a single topological pass.

$$D_r = 0 \quad (7.59)$$

$$D_{ij} = D_i + t'_{ij} \quad \forall (i, j) \in \mathcal{B} \quad (7.60)$$

$$D_i = \sum_{(h,i) \in \mathcal{B}} \alpha_{hi}^2 D_{hi} + \sum_{(h,i) \in \mathcal{B}} \sum_{\substack{(g,i) \in \mathcal{B} \\ (g,i) \neq (h,i)}} \alpha_{hi} \alpha_{gi} \sqrt{D_{hi} D_{gi}} \quad \forall i \neq r \quad (7.61)$$

Figure 7.11 shows the D labels associated with the example bush in panels (c) and (d). Panel (c) shows the D_{ij} labels associated with links, and panel (d) shows the D_i labels associated with nodes.

7.4.2 Shifting flows on a bush

With the labels defined in the previous subsection, we are now ready to state the Algorithm B, OBA, and LUCE flow shifting procedures. Any of these can be used for Step 2 of the generic bush-based algorithm presented at the start of the section. You will notice that all of these procedures follow the same general form: calculate bush labels (different labels for different algorithms) in *forward* topological order, then scan each node in turn in *reverse* topological order. When scanning a node, use the labels to identify vehicles entering the node from higher-cost approaches, and shift them to paths using lower-cost approaches. Update the x and/or α labels accordingly, then proceed to the previous node topologically until the origin has been reached.

All three of these algorithms also make use of *divergence nodes* (also called *last common nodes* or *pseudo-origins* in the literature) as a way to limit the scope of these updates. While the definition of divergence nodes is slightly different in these algorithms, the key idea is to find the “closest” node which is common to all of the paths travelers are being shifted among. The rest of this subsection details these definitions and the role they play.

Algorithm B

Algorithm B identifies the longest and shortest paths to reach a node, and shifts flows between them to equalize their travel times. That is, when scanning a node i , only two paths (the longest and shortest) are considered. It is easy to determine these paths using the L and U labels, tracing back the shortest and longest paths by identifying the links used for the minimum or maximum in equations (7.52) and (7.55). Once these paths are identified, the divergence node a is the last node common to both of these paths. The shortest and longest path segments between nodes a and i form a *pair of alternate segments*; let σ_L and σ_U denote these path segments. Within Algorithm B, flow is shifted from the longest path to the shortest path, using Newton's method to determine the amount of flow to shift:

$$\Delta h = \frac{(U_i - U_d) - (L_i - L_d)}{\sum_{(g,h) \in \sigma_L \cup \sigma_U} t'_{gh}} \quad (7.62)$$

There is also the constraint $\Delta h \leq \min_{(i,j) \in \sigma_U} x_{ij}$ which must be imposed to ensure that all links retain nonnegative flow after the shift. This flow is subtracted from each of the links in the longest path segment σ_U , and added to each of the links in the shortest path segment σ_L . The derivation of this formula parallels that of gradient projection in Section 7.3.2. By shifting flow from the longer path to the shorter path, we will either (i) equalize the travel times on the two paths or (ii) shift all the flow onto the shorter path, and have it still be faster than the longer path. In either case, we move closer to satisfying the equilibrium condition.

The steps of Algorithm B are as follows:

1. Calculate the L and U labels in forward topological order.
2. Let i be the topologically last node in the bush.
3. Scan i by performing the following steps:
 - (a) Use the L and U labels to determine the divergence node a and the pair of alternate segments σ_L and σ_U .
 - (b) Calculate Δh using equation (7.62) (capping Δh at $\min_{(i,j) \in \sigma_U} x_{ij}$ if needed).
 - (c) Subtract Δh from the x label on each link in σ_U , and add Δh to the x label on each link in σ_L .
4. If $i = r$, go to the next step. Otherwise, let i be the previous node topologically and return to step 3.
5. Update all travel times t_{ij} and derivatives t'_{ij} using the new flows x (remembering to add flows from other bushes.)

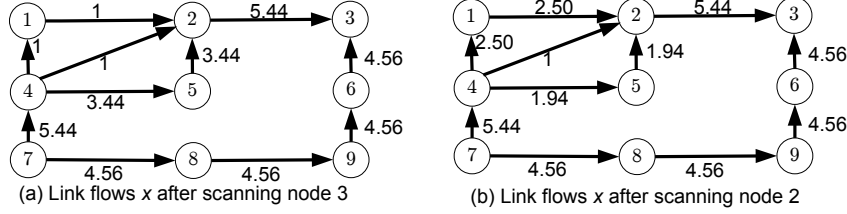


Figure 7.12: Algorithm B flow shifts.

Demonstrating on the example in Figures 7.10 and 7.11, we start with the L and U labels as shown in Figure 7.11(a) and (b), and start by letting $i = 3$, the last node topologically. The longest path in the bush from the origin ($r = 7$) to node $i = 3$ is $[7,4,5,2,3]$, and the shortest path is $[7,8,9,6,3]$, as can be easily found from the L and U labels. The divergence node is the last node common to both of these paths, which in this case is the origin, so $a = 7$. Equation (7.62) gives $\Delta h = 1.56$, so we shift this many vehicles away from the longest path and onto the shortest path, giving the flows in Figure 7.12(a).

The second-to-last node topologically is node 2, and we repeat this process. The longest and shortest paths from the origin to node 2 in the bush are $[7,4,5,2]$ and $[7,4,1,2]$, respectively.¹⁰ The last node common to both of these paths is 4, so the divergence node is $a = 4$ and we shift flow between the pair of alternate segments $[4,5,2]$ and $[4,1,2]$. Using equation (7.62) gives $\Delta h = 1.50$. Shifting this many vehicles from the longer segment to the shorter one gives the flows in Figure 7.12(b).

The previous node topologically is node 5. Since there is only one incoming bush link to node 5, there is nothing for Algorithm B to do. To see why, notice that the longest and shortest bush paths are $[7,4,5]$ and $[7,4,5]$. The divergence node would be $a = 5$, which is the same as i , and the “pair of alternate segments” is the empty paths $[5]$ and $[5]$. Intuitively, since there is only one way to approach node 5, there are no “alternate routes” to divert incoming flow. In fact, the same is true for all of the previous nodes topologically (1, 4, 6, 9, 8, 7 in the reverse of the order given above.), so there are no more flow shifts on the bush.

Origin-based assignment (OBA)

Rather than considering just two paths at a time, as Algorithm B did, OBA shifts flow among many paths simultaneously. This can be both a blessing (in that its moves affect more paths simultaneously) and a curse (in that its moves are not as sharp as Algorithm B, which can focus on the two paths with the

¹⁰Notice that we have not yet updated the travel time labels based on the shift at node 7. You may do so if you wish, but it is not required for Algorithm B to work, and in all of the examples in this section travel times are not updated until all flow shifts are complete for the bush.

greatest travel time difference). Because OBA shifts flow among many paths, it makes use of the average cost labels M and the derivative labels D , rather than the shortest and longest path costs (and the direct link travel times derivatives) used by Algorithm B. OBA also makes use of a Newton-type shift, dividing a difference in travel times by an approximation of this difference's derivative.

When scanning a node i in OBA, first identify a least-travel time approach, that is, a link (\hat{h}, i) such that $M_{\hat{h}i} \leq M_{hi}$ for all other approaches (h, i) . This link is called the *basic approach* in analogy to the basic path concept used in path-based algorithms. Then, for each nonbasic approach (h, i) , the following amount of flow is shifted from x_{hi} to $x_{\hat{h}i}$:

$$\Delta x_{hi} = \frac{M_{hi} - M_{\hat{h}i}}{D_{hi} + D_{\hat{h}i} - 2D_a}, \quad (7.63)$$

with the constraint that $\Delta x_{hi} \leq x_{hi}$ to prevent negative flows. It can be shown that this formula would equalize the mean travel times on the two approaches if they were linear functions. In reality, they are not, but the formula is still used as an approximation.

Here a is the divergence node, defined for OBA as the node with the highest topological order which is common to *all* paths in the bush from r to i , excluding i itself. This is a reasonable place to truncate the search, since shifting flow among path segments between a and i will not affect the flows on any earlier links in the bush. This is a generalization of the definition used for Algorithm B, needed since there are more than two paths subject to the flow shift.

After applying the shift Δx_{hi} to each of the links entering node i , we have to update flows on the other links between a and i to maintain flow conservation. This is done by assuming that the α values stay the same elsewhere, meaning that any increase or decrease in the flow passing through a node propagates backward to its incoming links in proportion to the contribution each incoming link provides to that total flow. Thus, the x labels can be recalculated using equations (7.48) and (7.49) to links and nodes topologically between a and i .

The steps of OBA are as follows:

1. Calculate the M and D labels in forward topological order.
2. Let i be the topologically last node in the bush.
3. Scan i by performing the following steps:
 - (a) Determine the divergence node a corresponding to node i , and the basic approach (\hat{h}, i) .
 - (b) For each nonbasic approach (h, i) calculate Δx_{hi} using equation (7.63) (with $\Delta x_{hi} \leq x_{hi}$), subtract Δx_{hi} from x_{hi} and add it to $x_{\hat{h}i}$.
 - (c) Update α_{hi} for every bush link terminating at i .
 - (d) Apply equations (7.48) and (7.49) to all nodes and links topologically between a and i , updating their x_{ij} and x_i values.

4. If $i = r$, go to the next step. Otherwise, let i be the previous node topologically and return to step 3.
5. Update all travel times t_{ij} and derivatives t'_{ij} using the new flows x (remembering to add flows from other bushes.)

Again demonstrating on the example in Figure 7.10, we use the M and D labels shown in Figure 7.11. As with Algorithm B, we scan nodes in reverse topological order. Starting with node 3, we compare the labels M_{23} and M_{63} to determine the basic approach. Since M_{63} is lower, this is the least-cost approach and named basic. For OBA, the divergence node corresponding to node 3 is node 7, since it is the only node common to all paths between nodes 7 and 3, except for node 3 itself. Formula (7.63) is then applied to determine the amount of flow Δx_{23} that should be shifted from approach (2,3) to the basic approach (6,3), producing $\Delta x_{23} = 1.48$. Shifting this flow updates the α values as shown in Figure 7.13(b). To maintain flow conservation, we also need to adjust the flow on other links in the bush. Holding the α proportions fixed at all other links, these changes at node 3 are propagated back proportional to the original flows, producing the link flows in Figure 7.13(a).

Next scanning node 2, we examine M_{12} , M_{42} , and M_{52} to determine the basic approach: since M_{12} is smallest, it is deemed the basic approach. The divergence node a corresponding to node 2 is node 4, since all paths from 7 to 2 pass through node 4 (and no other node with higher topological order). We now apply two shifts, calculating Δx_{42} and Δx_{52} with equation (7.63). This equation gives $\Delta x_{42} = 2.5$, but since $x_{42} = 0.79$ it is only possible to move 0.79 units of flow from (4,2) to (1,2). For approach (5,2), $\Delta x_{52} = 1.50$, so we move 1.50 units of flow from (5,2) to (1,2). Assuming that the α values at all other nodes remains constant, propagating this change back produces the x and α labels in Figure 7.13(c) and (d). Notice that we only had to update flows between the divergence node ($a = 4$) and the node being scanned ($i = 2$), since no other links would be affected by shifting flow among approaches to node 2.

As with Algorithm B, no changes are made for the remaining nodes scanned, since they only have one incoming link. This incoming link is trivially the “basic” approach, and there are no nonbasic approaches to shift flow from.

Linear user cost equilibrium (LUCE)

The LUCE algorithm is similar in structure to OBA, using the same definition of a divergence node and the same sets of labels (M and D). Where it differs is in how nodes are scanned. When LUCE scans a node, it attempts to solve a “local” user equilibrium problem, based on linear approximations to the approach travel times. To be concrete, when node i is scanned, the labels M_{hi} and D_{hi} respectively indicate the average travel time for travelers arriving at i via (h,i) , and the derivative of this average travel time as additional vehicles arrive via this link. Both of these are calculated based on the current bush flows, so if the flow on approach (h,i) changes by Δx_{hi} , then the new average travel time

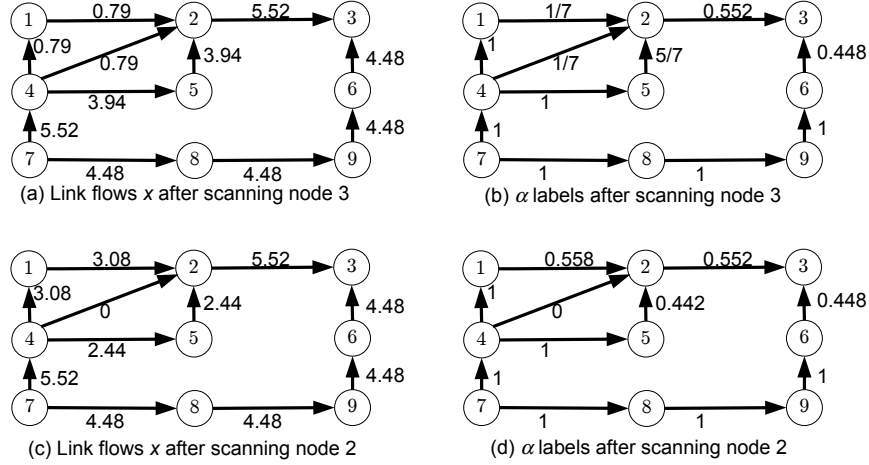


Figure 7.13: OBA flow shifts.

M'_{hi} can be approximated by

$$M'_{hi} \approx M_{hi} + \Delta x_{hi} D_{hi}. \quad (7.64)$$

The LUCE algorithm chooses Δx_{hi} values for each approach according to the following principles:

- (a) Flow conservation must be obeyed, that is, $\sum_{(h,i) \in \mathcal{B}} \Delta x_{hi} = 0$.
- (b) No link flow can be made negative, that is, $x_{hi} + \Delta x_{hi} \geq 0$ for all $(h, i) \in \mathcal{B}$.
- (c) The M'_{hi} values should be equal and minimal for any approach with positive flow after the shift, that is, if $x_{hi} + \Delta x_{hi} > 0$, then M'_{hi} must be less than or equal to the M' label for any other approach to i .

The Δx_{hi} values satisfying this principles can be found using a “trial and error” algorithm, like that introduced in Section 5.2.1. Choose a set of approaches (call it \mathcal{A}), and set $\Delta x_{hi} = -x_{hi}$ for all (h, i) not in this set \mathcal{A} . For the remaining approaches, solve the linear system of equations which set M'_{hi} equal for all $(h, i) \in \mathcal{A}$ and which have $\sum_{(h,i) \in \mathcal{A}} \Delta x_{hi} = 0$. The number of equations will equal the number of approaches in \mathcal{A} . After obtaining such a solution, you can verify whether the three principles are satisfied. Principle (a) will always be satisfied. If principle (b) is violated, approaches with negative $x_{hi} + \Delta x_{hi}$ should be removed from \mathcal{A} . If principle (c) is violated, then some approach not in \mathcal{A} has a lower M' value, and that approach should be added to \mathcal{A} . In either of the latter two cases, the entire process should be repeated with the new \mathcal{A} set.

The steps of LUCE are as follows:

1. Calculate the M and D labels in forward topological order.
2. Let i be the topologically last node in the bush.
3. Scan i by performing the following steps:
 - (a) Determine the divergence node a corresponding to node i , and the basic approach (\hat{h}, i) .
 - (b) Use the process described above to find Δx_{hi} values satisfying the local equilibrium principles (a)–(c) above, and add Δx_{hi} to x_{hi} for each approach to i .
 - (c) Update α_{hi} for every bush link terminating at i .
 - (d) Use equations (7.48) and (7.49) to all nodes and links topologically between a and i , updating their x_{ij} and x_i values.
4. If $i = r$, go to the next step. Otherwise, let i be the previous node topologically and return to step 3.
5. Update all travel times t_{ij} and derivatives t'_{ij} using the new flows x (remembering to add flows from other bushes.)

Applying LUCE to the same example, we again start by scanning node 2. Assuming that both approaches (2,3) and (6,3) will continue to be used, we solve the following equations simultaneously for Δx_{23} and Δx_{63} :

$$\begin{aligned} M_{23} + D_{23}\Delta x_{23} &= M_{63} + D_{63}\Delta x_{63} \\ \Delta x_{23} + \Delta x_{63} &= 0 \end{aligned}$$

Substituting the M and D labels from Figure 7.11 and solving, we obtain $\Delta x_{23} = -1.48$ and $\Delta x_{63} = +1.48$. Updating flows as in OBA (using the divergence node $a = 7$ and assuming all α values at nodes other than 3 are fixed) gives the x and α labels shown in Figure 7.14(a) and (b). Notice that this step is exactly the same as the first step taken by OBA. The interpretation of LUCE and solving a local, linearized equilibrium problem provides insight into how the OBA flow shift formula (7.63) was derived.

The shift is slightly different when there are three approaches to a node, as happens when we proceed to scan node 2. First assuming that all three approaches will be used, we solve these three equations simultaneously, enforcing flow conservation and that the travel times on the three approaches should be the same:

$$\begin{aligned} M_{12} + D_{12}\Delta x_{12} &= M_{42} + D_{42}\Delta x_{42} \\ M_{52} + D_{52}\Delta x_{52} &= M_{42} + D_{42}\Delta x_{42} \\ \Delta x_{12} + \Delta x_{42} + \Delta x_{52} &= 0 \end{aligned}$$

Substituting values from Figure 7.11 and solving these equations simultaneously gives $\Delta x_{12} = 2.82$, $\Delta x_{42} = -1.85$, and $\Delta x_{52} = -0.97$. This is problematic, since

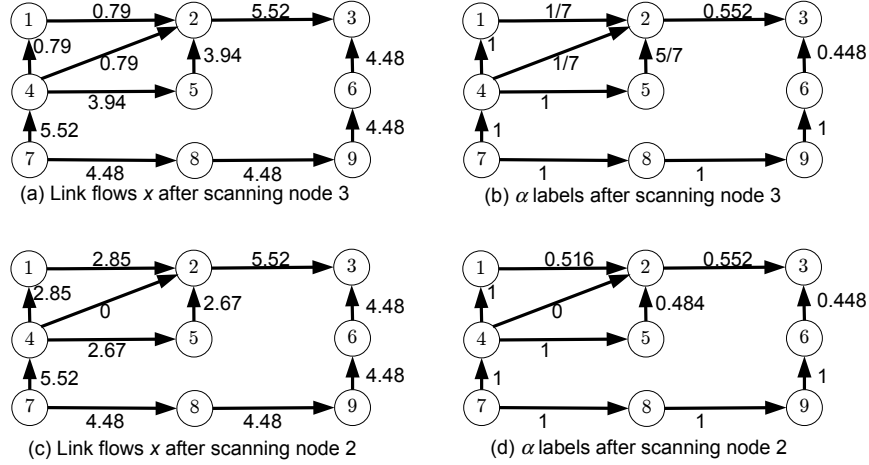


Figure 7.14: LUCE flow shifts.

$x_{42} = 0.79$ and it is impossible to reduce its flow further by 1.85. This means that approach (4,2) should not be used, so we fix $\Delta x_{42} = -0.79$ as a constant and re-solve the system of equations for Δx_{12} and Δx_{52} :

$$\begin{aligned} M_{12} + D_{12}\Delta x_{12} &= M_{42} - 0.79D_{42} \\ M_{52} + D_{52}\Delta x_{52} &= M_{42} - 0.79D_{42} \\ \Delta x_{12} - 0.79 + \Delta x_{52} &= 0 \end{aligned}$$

This produces $\Delta x_{12} = 2.06$, and $\Delta x_{52} = -1.27$, alongside the fixed value $\Delta x_{42} = -0.79$. Updating flows on other links as in OBA produces the x and α labels shown in Figure 7.14(b) and (c). No other shifts occur at lower topologically-ordered nodes, because there is only one incoming link, and flow conservation demands that no flow increase or decrease take place.

7.4.3 Improving bushes

After shifting flow for the current bushes as described in the previous section, the next step is to determine whether the bushes themselves need to change (adding or dropping links) to allow us to move closer to equilibrium on the entire network.

To begin, any link which is unused can be dropped from a bush without disturbing the solution (unless it is needed for connectivity); and it can always be added back later if we need to. As far as which links to add, one approach is to add any “shortcut” links, defined based on the bush travel times. If we solve each bush precisely to equilibrium in Step 2, then the distance from the origin r to each node i is the same on any used path; call this L_i . A “shortcut” link (i, j)

is any link for which $L_i + t_{ij} < L_j$, where the L labels have been re-calculated after eliminating the zero-flow bush links. Such a link can justifiably be called a shortcut, because it provides a faster path to node j than currently exists in the bush.

If this is our rule for adding links to the bush, it is easy to see that no cycles are created. Reasoning by contradiction, assume that a cycle *is* created, say, $[i_1, i_2, \dots, i_k, i_1]$. If link travel times are strictly positive, it is easy to see that $L_{i_{j+1}} > L_{i_j}$ for any successive pair of nodes in the cycle (either (i_j, i_{j+1}) was in the previous bush, in which case $L_{i_j} + t_{ij i_{j+1}} = L_{i_{j+1}}$ by equilibrium; or (i_j, i_{j+1}) was just added, in which case $L_{i_j} + t_{ij i_{j+1}} < L_{i_{j+1}}$). Applying this identity cyclically, we have $L_{i_1} < L_{i_2} < \dots < L_{i_k} < L_{i_1}$, a contradiction since we cannot have $L_{i_1} < L_{i_1}$.

The requirement that we solve each bush exactly to equilibrium in step 2 is vital to ensuring no cycles are created. If this is not the case (and in practice, we can only solve equilibrium approximately), another criterion is needed. To see why, the argument used to show that no cycles would be created relied critically on the assumption the difference in minimum cost labels for adjacent nodes in the cycle was exactly the travel time of the connecting link (except possibly for a link just added), which is only possible if all used paths have equal and minimal travel time. Another one, which is almost as easy to implement, is to add links to the bush based on the *maximum* travel time to a node using bush links, as opposed to the *minimum* travel time.

Assume that we have re-calculated the longest path labels U_i to each bush node, after eliminating unused bush links as described above. We now define a “modified shortcut” link as any link for which $U_i + t_{ij} < U_j$ (this is the same definition as a shortcut link but with maximum costs labels used instead of minimum cost labels). Using a similar argument, we can show that adding modified shortcut links cannot create cycles, even if the bush is not at equilibrium. Arguing again by contradiction, assume that the cycle $[i_1, i_2, \dots, i_k, i_1]$ is created by adding modified shortcut links, and consider in turn each adjacent pair of labels L_{i_j} and $L_{i_{j+1}}$. By the definition of maximum cost, for any link which was in the bush before we have $U_{i_j} + t_{ij i_{j+1}} = U_{i_{j+1}}$. For any link just added, we have $U_{i_j} + t_{ij i_{j+1}} < U_{i_{j+1}}$. Applying this identity cyclically, we have $U_{i_1} \leq U_{i_2} \leq \dots \leq U_{i_k} \leq U_{i_1}$. Furthermore, since there were no cycles in the bush during the previous iteration, at least one of these links must be new, and for this link the inequality must be strict.

As an example, consider the bush updates which occur after performing the LUCE example in the previous section. Figure 7.15 shows the updated link travel times in panel (a), and the remaining bush links after zero-flow links are removed in panel (b). Re-calculating L and U labels with the new travel times and bush topology gives the values in Figure 7.15(c). At this point the three unused bush links (4,2), (5,6), and (8,5), are examined to determine whether $U_i + t_{ij} < U_j$ for any of them. This is true for (5,6) and (8,5), since $41.6 + 2 < 66.3$ and $22.1 + 2 < 41.6$, but false for (4,2) since $32.5 + 42 \geq 72.9$. So, (5,6) and (8,5) are added to the bush, as shown in Figure 7.15(d). From here, one can return to the flow shifting algorithm to update flows further.

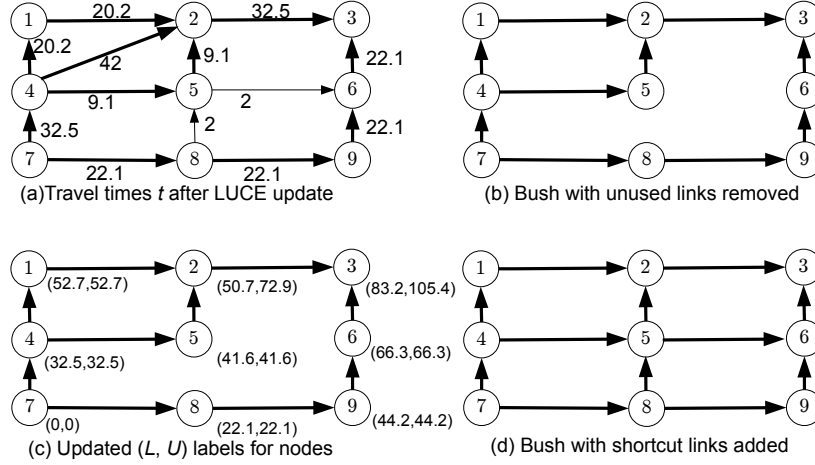


Figure 7.15: Updating a bush by removing unused links and adding shortcuts.

7.5 Likely Path Flow Algorithms

Recall from Section 6.2.2 that there are generally an infinite number of path flow vectors \mathbf{h} which satisfy the principle of user equilibrium. This contrasts with the fact that the equilibrium link flow vector $\hat{\mathbf{x}}$ is unique as long as link performance functions are increasing. That section introduced the principle of entropy maximization as a way to select a equilibrium path flow believed to be the most likely to occur in practice. The entropy-maximizing link flows solve the optimization problem

$$\max_{\mathbf{h}} \quad - \sum_{(r,s) \in Z^2} \sum_{\pi \in \Pi^{rs}} h_{\pi} \log(h_{\pi}/d^{rs}) \quad (7.65)$$

$$\text{s.t.} \quad \sum_{\pi \in \Pi} \delta_{ij}^{\pi} h_{\pi} = \hat{x}_{ij} \quad \forall (i,j) \in A \quad (7.66)$$

$$\sum_{\pi \in \Pi^{rs}} h_{\pi} = d^{rs} \quad \forall (r,s) \in Z^2 \quad (7.67)$$

$$h_{\pi} \geq 0 \quad \forall \pi \in \Pi \quad (7.68)$$

That section also introduced the proportionality condition, showing in Theorem 6.4 that the ratio of the flows on any two paths connecting the same OD pair depends *only* on the pairs of alternate segments where the paths differ, not what that OD pair happens to be, or on any links where the paths coincide. The proportionality condition is easier to verify, and entropy-maximizing solutions must satisfy proportionality. Proportionality does not imply entropy maximization, but in practical terms they seem nearly equivalent. Therefore,

this section focuses on finding proportional solutions, hence the term “likely path flow algorithms” rather than “*most* likely.”

The proof that entropy maximization implies proportionality was based on the Lagrangian of the entropy-maximization problem:

$$\begin{aligned} \mathcal{L}(\mathbf{h}, \boldsymbol{\beta}, \boldsymbol{\gamma}) = & - \sum_{(r,s) \in Z^2} \sum_{\pi \in \hat{\Pi}^{rs}} h_\pi \log \left(\frac{h_\pi}{d^{rs}} \right) + \sum_{(i,j) \in A} \beta_{ij} \left(\hat{x}_{ij} - \sum_{\pi \in \Pi} \delta_{ij}^\pi h_\pi \right) \\ & + \sum_{(r,s) \in Z^2} \gamma_{rs} \left(d^{rs} - \sum_{\pi \in \hat{\Pi}^{rs}} h_\pi \right). \end{aligned} \quad (7.69)$$

After some algebraic manipulation, we obtained the formula

$$h^\pi = K_{rs} \exp \left(- \sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij} \right), \quad (7.70)$$

where K_{rs} is a proportionality constant associated with the OD pair (r, s) corresponding to path π . The value of K_{rs} can be found from the constraint that total path flows must equal demand, $\sum_{\pi \in \hat{\Pi}^{rs}} h^\pi = d^{rs}$. The rest of the equation shows that the entropy-maximizing path flows are determined solely by the Lagrange multipliers β_{ij} associated with each link, and therefore that the ratio between two path flows for the same OD pair only depends on the links where they differ.

Here we describe three ways to calculate likely path flows. The first method is “primal,” and operates directly on the path flow vector \mathbf{h} itself. The second method is “dual,” operating on the Lagrange multipliers $\boldsymbol{\beta}$ in the entropy maximization problem, which can then be used to determine the path flows. Both of these methods presume that an equilibrium link flow solution $\hat{\mathbf{x}}$ has already been found. The third method, “traffic assignment by paired alternative segments” (TAPAS), is an algorithm which simultaneously solves for equilibrium and likely paths.

Primal and dual methods for likely path flows have two major issues in common:

1. The sets of equilibrium paths $\hat{\Pi}_{rs}$ must be determined in some way. If $\hat{\mathbf{x}}$ is exactly a user equilibrium, we can identify the shortest path for each OD pair, and then set $\hat{\Pi}_{rs}$ to be all the paths with that same cost between r and s . In practice, though, we cannot solve for the equilibrium link flows exactly, but only to a finite precision. So we cannot guarantee that all of the used paths for an OD pair have the same travel time, and we cannot guarantee that the unused paths are actually unused at the true equilibrium solution.
2. The system of constraints (7.66) and (7.67) has many more variables than equations; the number of equilibrium paths is typically much larger than

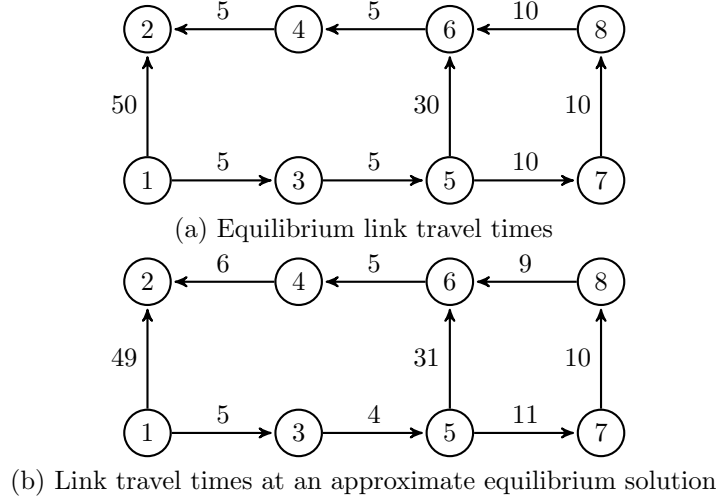


Figure 7.16: Link travel times at the true equilibrium solution (top) and at an approximate equilibrium solution (bottom).

the number of OD pairs and network links. For primal methods, we need to find and specify these “degrees of freedom” so we know how to adjust path flows while keeping the solution feasible (matching equilibrium link flows and the OD matrix). For dual methods, it means that many of the link flow constraints (7.66) are redundant, and their corresponding β_{ij} Lagrange multipliers are not needed. If these redundant constraints and Lagrange multipliers are included, then there are infinitely many β values that correspond to entropy maximization.

To address the first issue, we can include all paths in $\hat{\Pi}$ that are within some threshold ϵ of the shortest-path travel time for each OD pair. Some care must be taken to ensure that the resulting path sets are “consistent,” in that travelers from different origins and destinations consider the same sets of alternatives, when their path sets overlap. Figure 7.16 illustrates the difficulties in obtaining a consistent solution. The top panel of the figure shows the travel times at the true equilibrium solution, and the bottom panel shows the travel times at an approximate equilibrium solution, such as the one obtained after stopping one of the algorithms in this chapter after a finite number of iterations.

In this figure, assume there are two OD pairs, one from node 1 to node 2, and another from node 3 to node 4. At the true equilibrium solution, there are three equal-cost paths between nodes 1 and 2 — $[1, 2]$, $[1, 3, 5, 6, 4, 2]$, and $[1, 3, 5, 7, 8, 6, 4, 2]$ — and two equal-cost paths between nodes 3 and 4 — $[3, 5, 6, 4]$ and $[3, 5, 7, 8, 6, 4]$. Travelers whose paths cross nodes 5 and 6 consider two possible alternative routes between these nodes ($[3, 5, 6, 4]$ and $[3, 7, 8, 6, 4]$), regardless of which OD pair they came from. So defining $\hat{\Pi}_{12}$ and $\hat{\Pi}_{34}$ to include the equal-cost paths listed above is consistent.

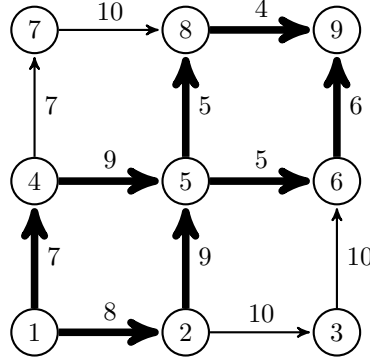


Figure 7.17: Example demonstrating rejection and acceptance gaps. Thick links are parts of used paths at the equilibrium solution. Link labels are travel times.

Now consider the bottom panel of the figure. There is now a unique shortest path for each OD pair — $[1, 2]$ and $[3, 5, 6, 4]$ — but this does not reflect the true equilibrium, simply the imprecision in an approximate solution. The threshold ϵ can be set to reflect this. Assume that $\epsilon = 1$, so that any path within 1 minute of the shortest path travel time is included in $\hat{\Pi}$. This choice gives the path sets $\hat{\Pi}_{12} = \{[1, 2], [1, 3, 5, 6, 4, 2]\}$ and $\hat{\Pi}_{34} = \{[3, 5, 6, 4], [3, 5, 7, 8, 6, 4]\}$. This choice is *not* consistent, because it assumes travelers passing between nodes 5 and 6 consider different choices depending on their OD pair: travelers starting at node 1 only consider the segment $[5, 6]$, while travelers starting at node 3 consider both $[5, 6]$ and $[5, 7, 8, 6]$ as options. If travelers are choosing routes to minimize cost, it should not matter what their origin or destination is. Increasing ϵ to a larger value would address this problem, but in a larger network would run the risk of including paths which are not used at the true equilibrium solution.

So some care must be taken in how ϵ is chosen. In the approximate equilibrium solution, we can define the *acceptance gap* g_a to be the greatest difference between a used path's travel time and the shortest path travel time for its OD pair, and the *rejection gap* g_r to be the smallest difference between the travel time of an unused path, and the shortest path for an OD pair. In Figure 7.17, the OD pairs are from 1 to 6 and 4 to 9, and the links are labeled with their travel times. The thick lines show the links used by these OD pairs, and the thin lines show unused links. For travelers between nodes 1 and 6, the used paths have travel times of 21 (shortest) and 22 minutes, and the unused path has a travel time of 26 minutes. For travelers between nodes 4 and 9, the used paths have travel times 18 (which is shortest) and 20 minutes, while the unused path has a travel time of 21 minutes. The acceptance gap g_a for this solution is 2 minutes (difference between 20 and 18), and the rejection gap is 5 (difference between 26 and 21).

It is possible to show that if ϵ is at least equal to the acceptance gap, but less than half of the rejection gap, then the resulting sets of paths $\hat{\Pi}_{rs}$ are consistent

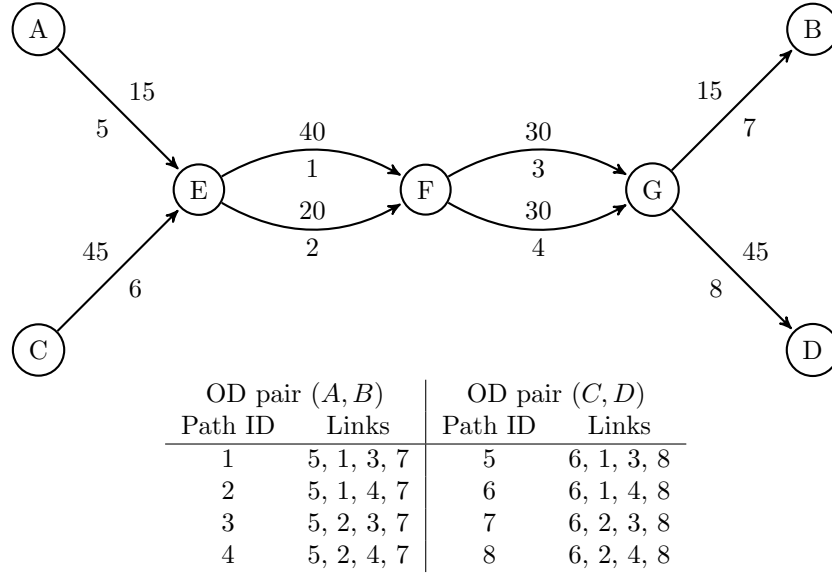


Figure 7.18: Example demonstrating redundancies and “degrees of freedom” when choosing a path flow solution. Links are labeled with equilibrium flows (above) and link IDs (below).

for the proportionality condition. That is, we need

$$g_a \leq \epsilon \leq \frac{g_r}{2}. \quad (7.71)$$

In the network of Figure 7.17, any ϵ choice between 2 and 2.5 will thus lead to a consistent set of paths. Exercise 21 gives a more formal definition of consistency and asks you to prove this statement. It is not possible to choose such an ϵ unless the equilibrium problem is solved with enough precision for the acceptance gap to be less than half of the rejection gap. To fully maximize entropy, rather than just satisfying proportionality, demands a more stringent level of precision in the equilibrium solution.

The second issue involves redundancies in the set of equations enforcing the OD matrix and equilibrium link flow constraints. Properly resolving this issue requires using linear algebra to analyze the structure of the set of equations (and this in fact is the key to bridging the gap between proportionality and entropy maximization), but for proportionality a simpler approach is possible.

A redundancy in a system of equations can be interpreted as a “degree of freedom,” a dimension along which a solution can be adjusted. Consider the network in Figure 7.18, which has two OD pairs (A to B, and C to D). The equilibrium link flows are shown in the figure, along with the link IDs and an indexing of the eight paths. The OD matrix and link flow constraints are

reflected in the six equations

$$h_1 + h_2 + h_3 + h_4 = 15 \quad (7.72)$$

$$h_5 + h_6 + h_7 + h_8 = 45 \quad (7.73)$$

$$h_1 + h_2 + h_5 + h_6 = 40 \quad (7.74)$$

$$h_3 + h_4 + h_7 + h_8 = 20 \quad (7.75)$$

$$h_1 + h_3 + h_5 + h_7 = 30 \quad (7.76)$$

$$h_2 + h_4 + h_6 + h_8 = 30 \quad (7.77)$$

Equations (7.72) and (7.73) reflect the constraints that the total demand among all paths from A to B, and from C to D, must equal the respective values in the OD matrix. Equations (7.74)–(7.77) reflect the constraints that the flow on links 1–4 must match their equilibrium values. Similar equations for links 5–8 are omitted, since they are identical to (7.72) and (7.73), as can easily be verified.

This system has eight variables but only six equations, so there must be at least two independent variables — in fact, there are four, since some of the six equations are redundant. For example, adding equations (7.72) and (7.73) gives you the same result as adding equations (7.74) and (7.75), so one of them — say, (7.75) can be eliminated. Likewise, equation (7.77) can be eliminated, since adding (7.72) and (7.73) is the same as adding (7.76) and (7.77). These choices are not unique, and there are other equivalent ways of expressing the same redundancies.

Each of these redundancies corresponds to an independent way to adjust the path flows without affecting either total path flows between OD pairs, or total link flows. A primal method uses these redundancies to adjust the path flows directly, increasing the entropy of the solution without sacrificing feasibility of the original path flow solution. A dual method uses these redundancies to eliminate unnecessary link flow constraints — for instance, equations (7.75) and (7.77) in the example above — so that the entropy-maximizing β_{ij} values are unique.

7.5.1 Primal method

To find a proportional solution with a primal method, we need (1) the equilibrium path set $\hat{\Pi}$; (2) the equilibrium link flows $\hat{\mathbf{x}}$; (3) an initial path flow solution \mathbf{h}^0 , and (4) a list of “redundancies” in the link flow constraints, each of which corresponds to a way to change path flows while maintaining feasibility. The algorithm then adjusts the \mathbf{h} values, increasing the entropy at each iteration, until termination.¹¹

The equilibrium path set and link flows were already described above. Depending on the algorithm used to solve for the equilibrium link flows, a path

¹¹There are alternative solution representations that can make this algorithm much faster, but would make the explanations more complicated. You are encouraged to think about how to implement this algorithm efficiently, without having to list all paths explicitly.

flow solution may already be available. If you solved for equilibrium with a path-based algorithm, its solution already contains the flows on each path. If you used a bush-based algorithm, a corresponding path flow can be found using the “within origins” adjustment technique described below. If you used a link-based algorithm, it is not as easy to directly identify a path flow solution from the final output — but it may be possible to track a path flow solution as the algorithm progresses (for instance, in Frank-Wolfe, each “target” all-or-nothing solution can be clearly identified with a path flow solution, and this can be averaged with previous path flow solutions using the same λ values).

To achieve proportionality, it is not necessary to identify *all* of the redundancies in the link flow constraints (it would be needed to fully maximize entropy). It is enough to identify *pairs of alternate segments* between two nodes which are used by multiple OD pairs. In Figure 7.18, links 1 and 2 are alternate segments between nodes E and F, and links 3 and 4 are alternate segments between nodes F and G.

The primal algorithm alternates between two steps: (1) adjusting \mathbf{h} to achieve proportionality for each origin r ; and (2) for each pair of alternative segments, adjusting \mathbf{h} to achieve proportionality between those nodes. These steps are described below; you can verify that each one of these steps preserves feasibility of the solution, and increases entropy. If the solution remains unchanged after performing both of these steps, then proportionality has been achieved and we terminate. Other stopping criteria can be introduced, by measuring the deviation from proportionality and terminating once this deviation is sufficiently small.

Proportionality within origins

Section 6.2.3 described how a user equilibrium solution can be described by the total flow from each origin r on each link (i, j) , denoted x_{ij}^r , and how at equilibrium the links with positive x_{ij}^r values must form an acyclic subnetwork, a bush. It is easy to obtain such a solution if we have a path flow solution \mathbf{h} :

$$x_{ij}^r = \sum_{s \in Z} \sum_{\pi \in \hat{\Pi}_{rs}} \delta_{ij}^{\pi} h^{\pi}. \quad (7.78)$$

With these values, we can calculate the fraction of the flow from origin r approaching any node j from one specific link entering that node (i, j) :

$$\alpha_{ij}^r = x_{ij}^r / \sum_{(h,j) \in \Gamma^{-1}(j)} x_{hj}^r \quad (7.79)$$

with α_{ij}^r defined arbitrarily if the denominator is zero. (This use of α is the same as in the bush-based algorithms described in Section 7.4).

We can ensure that the proportionality condition holds between *all* paths associated with an origin r by updating the path flows according to the formula

$$h^{\pi} \leftarrow d^{rs} \prod_{(i,j) \in \pi} \alpha_{ij}^r \quad \forall s \in Z, \pi \in \hat{\Pi}_{rs}, \quad (7.80)$$

that is, by applying the *aggregate* approach proportions across all paths from this origin to each individual path. One can show updating the path flows with this formula will not change either the total OD flows or link flows, maintaining feasibility, and will also increase entropy if there is any change in \mathbf{h} .

This process is repeated for each origin r .

Proportionality between origins

To achieve proportionality between different origins, we consider pairs of alternate segments between two nodes, used by multiple OD pairs. In Figure 7.18, an example of a pair of alternate segments is links 1 and 2. In a larger network, these alternate segments can contain multiple links. These links connect the same two nodes (E and F), and are parts of equal-cost paths for both OD pairs (A, B) and (C, D) . If we move some flow from link 1 to link 2 from OD pair (A, B) , and move the equivalent amount of flow from link 2 to link 1 from OD pair (C, D) , we will not disturb the total link flows or OD flows, but will change \mathbf{h} and allow us to increase entropy.

Let $\sigma = \{\sigma_1, \sigma_2\}$ be a pair of alternate segments, and let $Z^2(\sigma)$ be the set of OD pairs which have paths using both of the alternate segments in σ . For each segment σ_i , the segment flow $g(\sigma_i)$ is defined as the sum of flows on all paths using that segment:

$$g(\sigma_i) = \sum_{(r,s) \in Z^2(\sigma)} \sum_{\pi \in \hat{\Pi}_{rs} : \sigma_i \subseteq \pi} h^\pi, \quad (7.81)$$

where the notation $\sigma_i \subseteq \pi$ means that all of the links in the segment σ_i are in the path π .

If the equilibrium path set $\hat{\Pi}$ is consistent in the sense of (7.71), then any path π which uses one segment of the pair has a “companion” path which is identical, except it uses the other segment of the pair, denoted $\pi^c(\sigma)$. For example, in Figure 7.18, if σ is the pair of alternate segments between nodes E and F, the companion of path 1 is path 3, and the companion of path 6 is path 8.

To achieve proportionality between origins for the alternate segments in σ , we calculate the ratios between $g(\sigma_i)$ values and apply the same ratios to the path flows for each OD pair using this set of alternate segments:

$$h^\pi \leftarrow (h^\pi + h^{\pi^c(\sigma)}) \frac{g(\sigma_i)}{g(\sigma_1) + g(\sigma_2)} \quad \forall (r, s) \in Z^2(\sigma), i \in \{1, 2\}, \pi \in \hat{\Pi}_{rs} \wedge \pi \supseteq \sigma_i. \quad (7.82)$$

It is again possible to show that updating path flows with this formula leaves total OD flows and link flows fixed, and can only increase entropy.

Example

This section shows how the primal algorithm can solve the example in Figure 7.18. Assume that the initial path flow solution is $h_1 = 15$, $h_5 = 15$,

$h_6 = 10$, $h_8 = 20$, and all other path flows zero. As shown in the first column of Table 7.2, this solution satisfies the OD matrix and the resulting link flows match the equilibrium link flows, so it is feasible. The entropy of this solution, calculated using (7.65), is 47.7.¹² Two pairs of alternate segments are identified: links 1 and 2 between nodes E and F, and links 3 and 4 between nodes F and G.

This table summarizes the progress of the algorithm in successive columns; you may find it helpful to refer to this table when reading this section. The bottom section of the table shows the origin-based link flows corresponding to the path flow solution, calculated using (7.78).

The first iteration applies the within-origin formula (7.80) to origin A, and to origin C. To apply the formula to origin A, the origin-based proportions are first calculated with equation (7.79): $\alpha_1^A = 1$, $\alpha_2^A = 0$, $\alpha_3^A = 1$, and $\alpha_4^A = 0$, and thus $h_1 \leftarrow 15$, $h_2 \leftarrow 0$, $h_3 \leftarrow 0$, and $h_4 \leftarrow 0$. (There is no change.) For origin C, we have $\alpha_1^C = 5/9$, $\alpha_2^C = 4/9$, $\alpha_3^C = 1/3$, and $\alpha_4^C = 2/3$, and thus $h_5 \leftarrow 8\frac{1}{3}$, $h_6 \leftarrow 16\frac{2}{3}$, $h_7 \leftarrow 6\frac{2}{3}$, and $h_8 \leftarrow 13\frac{1}{3}$. The entropy of this new solution has increased to 59.6.

Next, we apply the between-origin formula to the pair of alternate segments between nodes F and G. Paths 1, 3, 5, and 7 use link 3; and their companion paths (2, 4, 6, and 8, respectively) use link 4. The segment flow for link 3 is the sum of the path flows that use it (30), and similarly the segment flow for link 4 is also 30. Thus formula (7.82) requires paths 1 and 2 to have equal flow, paths 3 and 4 to have equal flow, and so on for each path and its companion. Redistributing the flow between paths and companions in this way gives the result in the third column of Table 7.2, and the entropy has increased to 72.5.

The between-origin formula is applied a second time to the other pair of alternate segments, between nodes E and F. Paths 1, 2, 5, and 6 use link 1, and the companion paths using link 2 are 3, 4, 7, and 8, respectively. The segment flows for links 1 and 2 are 40 and 20, so for each path and its companion, the path using link 1 should have twice the flow of the path using link 2. The fourth column of Table 7.2 shows the results, and the entropy has increased again to 79.8.

This solution achieves proportionality (and in fact maximizes entropy). The proportionality conditions can either be checked directly, or noticed when running the algorithm a second time does not change any path flows. For larger networks with a more complicated structure, the algorithm generally requires multiple iterations, and only converges to proportionality in the limit.

7.5.2 Dual method

An alternative approach involves the optimality conditions directly. Recall from equation (7.70) that entropy-maximizing (and thus proportional) flow on each path is $K_{rs} \exp\left(-\sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij}\right)$, where K_{rs} is an OD-specific constant chosen so that the path flows sum to the total demand d_{rs} . We can adopt this

¹²When computing this formula, $0 \log 0$ is taken to be zero, since $\lim_{x \rightarrow 0^+} x \log x = 0$.

Table 7.2: Demonstration of primal proportionality algorithm.

		Initial solution	Within-origin	F/G PAS	E/F PAS
Entropy		47.7	59.6	72.5	79.8
Path flows	h_1	15	15	7.5	5
	h_2	0	0	7.5	5
	h_3	0	0	0	2.5
	h_4	0	0	0	2.5
	h_5	15	8.33	12.5	15
	h_6	10	16.67	12.5	15
	h_7	0	6.67	10	7.5
	h_8	20	13.33	10	7.5
OD flows	(A, B)	15	15	15	15
	(C, D)	45	45	45	45
Total link flows	x_1	40	40	40	40
	x_2	20	20	20	20
	x_3	30	30	30	30
	x_4	30	30	30	30
From origin A	x_1^A	15	15	15	10
	x_2^A	0	0	0	5
	x_3^A	15	15	7.5	7.5
	x_4^A	0	0	7.5	7.5
From origin B	x_1^B	25	25	25	30
	x_2^B	20	20	20	15
	x_3^B	15	15	22.5	22.5
	x_4^B	30	30	22.5	22.5

condition as a *formula* for \mathbf{h} , and no matter what values are chosen for β_{ij} , the path flows we calculate satisfy proportionality. The difficulty is that they will generally not be feasible, unless the link flows \mathbf{x} corresponding to \mathbf{h} happen to equal their equilibrium values $\hat{\mathbf{x}}$. A *dual* algorithm tries to adjust these β_{ij} values until $\mathbf{x} = \hat{\mathbf{x}}$, at which point we terminate with proportional (and in fact entropy-maximizing) path flows.

This contrasts with the *primal* approach in the previous section, which always maintained feasibility (the link flows in Table 7.2 never changed from the equilibrium values) and worked toward optimality, expressed in equation (7.70). The algorithm in this section always maintains optimality, and works toward feasibility. Unlike the primal algorithm, the path flows in the dual algorithm are not feasible until termination.

The idea behind the algorithm is simple enough: start with initial values for β_{ij} on each link; calculate \mathbf{h} from equation (7.70); calculate \mathbf{x} from \mathbf{h} ; and see which links have too much flow or too little flow. Adjust the β_{ij} values accordingly, and iterate until the flow on every link is approximately equal to its equilibrium value. There are a few details to take care of: how to find an initial solution, when to terminate, how to adjust the β_{ij} values, and dealing with redundancies in the system of constraints. The first two details are fairly simple. Any initial solution will do; $\beta = \mathbf{0}$ is simplest. Terminate when \mathbf{x} is “close enough” to $\hat{\mathbf{x}}$ according to some measure.)

For adjusting the β_{ij} values, notice from equation (7.70) that increasing β_{ij} will decrease x_{ij} , and vice versa. So a natural update rule is

$$\beta_{ij} \leftarrow \beta_{ij} + \alpha(x_{ij} - \hat{x}_{ij}) \quad (7.83)$$

where α is a step size, and $x_{ij} - \hat{x}_{ij}$ is the difference between the link flows currently implied by β , and the equilibrium values. In addition to this intuitive interpretation, this search direction is also proportional to the gradient of the least-squares function

$$\phi(\mathbf{x}) = \sum_{(i,j) \in A} (x_{ij} - \hat{x}_{ij})^2. \quad (7.84)$$

This function is zero only at a feasible solution, and (7.83) is a steepest descent direction in terms of \mathbf{x} .¹³

Equation (7.84) can also be used to set the step size α . One can select a trial sequence of α values (say, 1, 1/2, 1/4, 1/8, ...), evaluating each α value in turn and stopping once the new \mathbf{x} values reduce (7.84). A more sophisticated step size rule chooses α using Newton’s method, to approximately maximize entropy. Newton’s method also has the advantage of scaling the step size based on the effect changes in β have on link flows. Exercise 23 develops this approach in more detail.

A last technical detail concerns redundancies in the system of link flow equations, as discussed at the end of Section 7.5. Redundancies in the link flow

¹³To be precise, it is not a steepest descent direction in terms of β . An alternative derivation of equation (7.83) is explored in Exercise 22.

equations mean that the β_{ij} values maximizing entropy may not be unique. To resolve this issue, redundant link flow constraints can be removed, and their β_{ij} values left fixed at zero. Practical experience shows that this can significantly speed convergence.

Example

The dual method is now demonstrated on the same example as the primal algorithm; see Figure 7.18. You may find it helpful to refer to Table 7.3 when reading this section to track the progress of the algorithm. The format is similar to Table 7.2, except for additional rows showing the β_{ij} values.

This table summarizes the progress of the algorithm in successive columns; you may find it helpful to refer to this table when reading this section. The bottom section of the table shows the origin-based link flows corresponding to the path flow solution, calculated using (7.78).

To begin, as discussed at the end of Section 7.5, two of the link flow constraints are redundant, and their β_{ij} values are fixed at zero. Assume that links 2 and 4 are chosen for this purpose, so $\beta_2 = \beta_4 = 0$ throughout the algorithm. (The algorithm would perform similarly for other choices of the two redundant links; note that we are also continuing to ignore the link flow constraints associated with links 5–8, since these are identical to the OD matrix constraints (7.72) and (7.73).)

Initially, $\beta_1 = \beta_3 = 0$. This means that $\exp\left(-\sum_{(i,j) \in A} \delta_{ij}^\pi \beta_{ij}\right) = 1$ for all paths, so $h_1 = h_2 = h_3 = h_4 = K_{AB}$ and $h_5 = h_6 = h_7 = h_8 = K_{CD}$. To ensure that the sum of each OD pairs' path flows equals the total demand, we need $K_{AB} = 3.75$ and $K_{CD} = 11.25$, and equation (7.70) gives the flows on each path, as shown in the Iteration 0 column of Table 7.3.

The table also shows the link flows corresponding to this solution: links 1–4 all have 30 vehicles, whereas the equilibrium solution has $x_1 = 40$ and $x_2 = 20$. The least-squares function (7.84) has the value $(40 - 30)^2 + (20 - 30)^2 = 200$. Trying an initial step size of $\alpha = 1$ would give $\beta_1 = 0 + 1 \times (30 - 40) = -10$. The other β_{ij} values are unchanged: β_3 remains at zero because it has the correct link flow, while β_2 and β_4 are permanently fixed at zero because their link flow constraints were redundant. Re-applying equation (7.70) with this new value of β_1 (and recalculating K_{AB} and K_{CD} to satisfy the OD matrix) would give $x_1 = 60$, $x_2 = 0$, and $x_3 = x_4 = 30$. This has a larger least-squares function than before (800 vs. 200), so we try again with $\alpha = 1/2$. This is slightly better (the least-squares function is 768), but still worse than the current solution. After two more trials, we reach $\alpha = 1/8$, which produces a lower mismatch (88).

This step size is accepted, and we proceed to the next iteration. The path and link flows are shown in the Iteration 1 column of Table 7.3. The flows on links 1 and 2 are closer to their equilibrium values than before. Continuing as before, we find that $\alpha = 1/8$ is again the acceptable step size with the new link flow values, so $\beta_1 \leftarrow -1.25 + \frac{1}{8}(36.2 - 40) = -0.42$, producing the values in the

Table 7.3: Demonstration of dual proportionality algorithm.

		Iteration 0	1	2	\dots	∞
Entropy		83.2	73.4	81.9		79.8
	β_1	0	-1.25	-0.42		-0.692
	β_3	0	0	0		0
Path flows	h_1	3.75	5.83	4.53		5
	h_2	3.75	5.83	4.53		5
	h_3	3.75	1.67	2.97		2.5
	h_4	3.75	1.67	2.97		2.5
	h_5	11.25	17.49	13.58		15
	h_6	11.25	17.49	13.58		15
	h_7	5.01	8.92	6.53		7.5
	h_8	5.01	8.92	6.53		7.5
OD flows	(A, B)	15	15	15		15
	(C, D)	45	45	45		45
Total link flows	x_1	30	46.6	36.2		40
	x_2	30	13.4	23.8		20
	x_3	30	30	30		30
	x_4	30	30	30		30
From origin A	x_1^A	7.5	11.7	9.1		10
	x_2^A	7.5	3.3	5.9		5
	x_3^A	7.5	7.5	7.5		7.5
	x_4^A	7.5	7.5	7.5		7.5
From origin B	x_1^B	22.5	35.0	27.1		30
	x_2^B	22.5	10.0	17.8		15
	x_3^B	22.5	22.5	22.5		22.5
	x_4^B	22.5	22.5	22.5		22.5

Iteration 2 column. Over additional iterations, the algorithm converges to the final values shown in the rightmost column.

It is instructive to compare the dual algorithm in Table 7.3 with the primal algorithm in Table 7.2. Notice how the dual algorithm always maintains proportionality, and the link flows gradually converge to their equilibrium algorithms. By contrast, the primal algorithm maintains the link flows at their equilibrium values, and gradually converges to proportionality. The entropy also does not change monotonically, and at times it is higher than the maximum entropy value — this can only happen for an infeasible solution.

7.5.3 Traffic assignment by paired alternative segments

The primal and dual methods described in the preceding sections assumed that user equilibrium link flows were already available, and then found likely path flows as a post-processing step. Traffic assignment by paired alternative segments (TAPAS) is an algorithm which finds the equilibrium solution and pro-

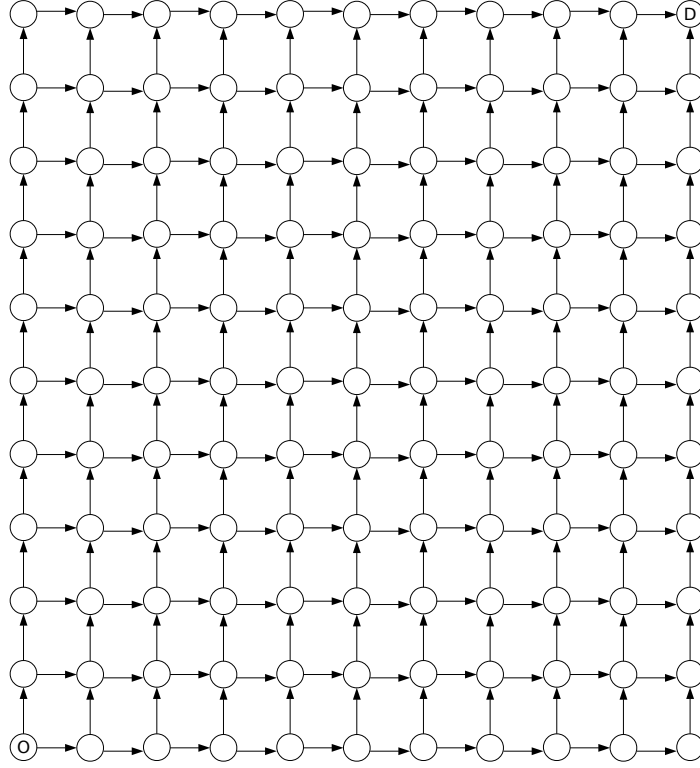


Figure 7.19: A grid network consisting of 100 city blocks.

portional path flows simultaneously. Interestingly, accomplishing both tasks at once does not slow down the algorithm. TAPAS is in fact among the fastest of the traffic assignment algorithms currently known.

Recall from Section 7.4 that path- and bush-based algorithms find the equilibrium solution by shifting flow from longer paths to shorter ones, and that these paths often differ on a relatively small set of links (in gradient projection, we denoted these by the set $A_3 \cup A_4$; in bush-based algorithms, by the concept of a divergence node). The main insights of TAPAS are that these algorithms tend to shift flow repeatedly between the same sets of links, and that these links are common to paths used by between different origins and destinations. As a result, it makes sense to store these path segments from one iteration to the next, rather than having to expend effort finding them again and again. Furthermore, since these links are common to multiple origins, we can apply proportionality concepts at the same time to find a high-entropy path flow solution.

Pairs of alternative segments can also form a concise representation of the equilibrium conditions. In the grid network of Figure 7.19, the number of paths between the origin in the lower-left and the destination in the upper-right is

rather large (in fact there are 184,756) even though the network is a relatively modestly-sized grid of ten rows and columns. If all paths are used at equilibrium, expressing the equilibrium condition by requiring the travel times on all paths be equal requires 184,755 equations.

The network can also be seen as a hundred “city blocks,” each of which can either be traversed in the clockwise direction (north, then east) or in the counterclockwise direction (east, then north). If travel times on all paths in the network are identical, then the travel time around each block must be the same in the clockwise and counterclockwise directions. In fact, the converse is true as well: if the travel time is the same around every block in both orientations, then the travel times on all paths in the network are the same as well. We can thus express equality of all network paths with only 100 equations!

Furthermore, other nodes in the network may serve as origins and destinations, not just nodes at two corners. Expressing equilibrium in terms of path travel time equality requires additional equations for each new OD pair, but the same 100 equations expressing equality of travel times around each block are sufficient no matter how many nodes serve as origins or destinations.

This discussion implies that most of the 184,755 path travel-time equations are redundant. It is not trivial to identify these linear dependencies, but methods based on paired alternative segments are a way to do so. The two ways to travel around each block can be seen as a pair of alternative segments between their southwestern and northeastern nodes, and the set of all such pairs of alternative segments is “spanning” in the sense that one can shift flow between any two paths with the same origin and destination by shifting flow between a sequence of pairs of alternative segments. There are some subtleties involving the equivalence of equilibrium conditions on pairs of alternative segments and on paths (see Exercise 24), but this example shows how they can often simplify the search for a user equilibrium solution.

The TAPAS algorithm represents network flows aggregated by origin, with x_{ij}^r denoting the flow on link (i, j) which started at origin r , following equation (6.30). The algorithm also involves a set of pairs of alternative segments (PASs). Each PAS ζ is defined by two path segments σ_1^ζ and σ_2^ζ starting and ending at the same nodes, and by a list of *relevant origins* Z_ζ indicating a subset of zones which have flow on both path segments at the current solution. The steps of the algorithm involve maintaining a set of PASs (creating new ones, updating relevant origins, and optionally discarding inactive ones), and adjusting the link flows x_{ij}^r to move towards user equilibrium and proportionality. Notice that TAPAS does not store the path flows \mathbf{h} explicitly, for computational reasons. Instead, the path flows are represented *implicitly*, obtained from the x_{ij}^r values using the proportional split formulas previously introduced in Section 7.5.1:

$$\alpha_{ij}^r = x_{ij}^r / \sum_{(h,j) \in \Gamma^{-1}(j)} x_{ij}^h \quad (7.85)$$

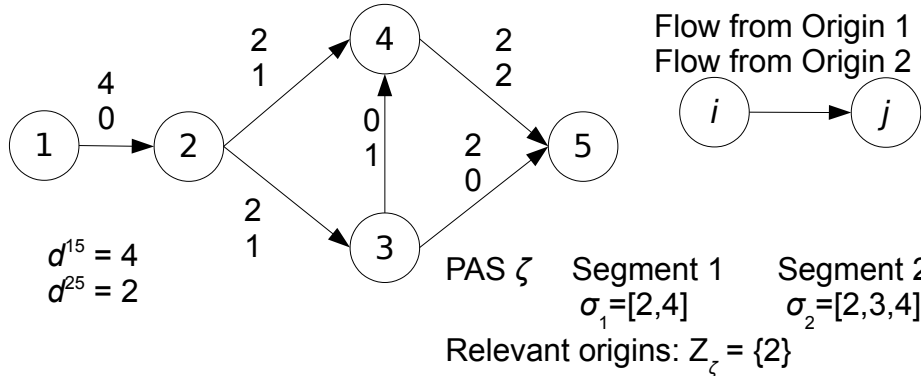


Figure 7.20: Example for demonstrating pairs of alternative segments.

where

$$x_{ij}^r = \sum_{s \in Z} \sum_{\pi \in \hat{\Pi}_{r,s}} \delta_{ij}^\pi h^\pi. \quad (7.86)$$

An example of a PAS is shown in Figure 7.20. For each link, the upper and lower labels give the flows on that link from Origin 1 and Origin 2, respectively. There are two path segments: $[2, 4]$ and $[2, 3, 4]$, and there is one relevant origin (Origin 2). You might expect that this PAS is also relevant to Origin 1; and indeed at the ultimate equilibrium solution this will be true. However, in large networks it is not immediately obvious which PASs are relevant to which origins, and the TAPAS algorithm must discover this during its steps.

There are three main components to the algorithm: PAS management, flow shifts, and proportionality adjustments. PAS management involves identifying new PAS, updating the lists of relevant origins, and removing inactive ones. Flow shifts move the solution closer to user equilibrium, by shifting vehicles from longer paths to shorter ones. Proportionality adjustments maintain the total link flows at their current values, but adjust the origin-specific link flows to increase the entropy of the path flow solution implied by (7.85). One possible way to perform these steps is as follows; the rest of this subsection fills out the details of each step.

1. Find an initial origin-disaggregated solution, and initialize the set of PASs to be empty.
2. Update the set of PASs by determining whether new ones should be created, or whether existing ones are relevant to more origins.
3. Perform flow shifts within existing PASs.
4. Perform proportionality adjustments within existing PASs.
5. Check for convergence, and return to step 2 unless done.

This algorithmic description may appear vague. Like many of the fastest algorithms currently available, the performance of the algorithm depends on successfully balancing these three components of the algorithm. The right amount of time to spend on each component is network- and problem-specific, and implementations that make such decisions adaptively, based on the progress of the algorithm, can work well.

Updating the set of pairs of alternative segments

In its second step, the TAPAS algorithm must update the set of PASs. Since flow shifts mainly occur within PASs, finding the user equilibrium solution relies on being able to identify new alternative routes which are shorter than the ones currently being used. Given the origin-disaggregated solution \mathbf{x}^r for each origin r , we can search for routes in the following way.

Solving a shortest path algorithm over the entire network produces a tree rooted at an origin r , containing paths to every node. (Note that already having an origin-disaggregated solution can greatly accelerate the process of finding shortest paths; see Exercise 48 from Chapter 2.) At equilibrium, essentially all of these links should be used.¹⁴ By comparing this tree to the links which are used in the disaggregate solution \mathbf{x}^r , we can identify any links in the shortest path tree not currently being used by an origin r even though there is flow to their head nodes.

More specifically, let \hat{A}^r denote the set of these links. A link (i, j) is in \hat{A}^r if (i, j) is in the shortest path tree rooted at r , if $x_{ij}^r = 0$ (it is currently unused by origin r), and if there is some other link (h, j) for which x_{hj}^r (flow from origin r is reaching the head node j in another way). We then look for a PAS ζ whose two segments σ_1^ζ and σ_2^ζ end with the links (h, j) and (i, j) , which will allow us to shift flow onto the shortest path segment.

Two possibilities exist: either such a PAS already exists (in which case we add origin r to the relevant set Z_ζ if it is not already listed), or we create a new one. To create a new PAS, we must choose two path segments σ_1^ζ and σ_2^ζ which start and end at the same node. The first segment σ_1^ζ should consist of links with positive flow from origin r , and the second one σ_2^ζ should consist of links from the shortest path tree. We also know they must both end at node j , but must choose an appropriate node for them to start (a divergence node). As discussed in Section 7.4.2, there are several ways to choose divergence nodes. For TAPAS, an ideal divergence node results in short path segments σ_1^ζ and σ_2^ζ . Short segments both result in faster computation, and intuitively are more likely to be relevant to more origins.

Putting these concepts together, we can search for a divergence node a for which there is a segment σ_1^ζ starting at a , ending with link (h, j) , and only using links with positive flow from origin r ; and a segment σ_2^ζ starting at a , ending with link (i, j) , and only using links in the current shortest path tree. Among

¹⁴The only exceptions would be to links leading to nodes not used by this origin, or if multiple paths are tied for being shortest with one of them having zero flow.

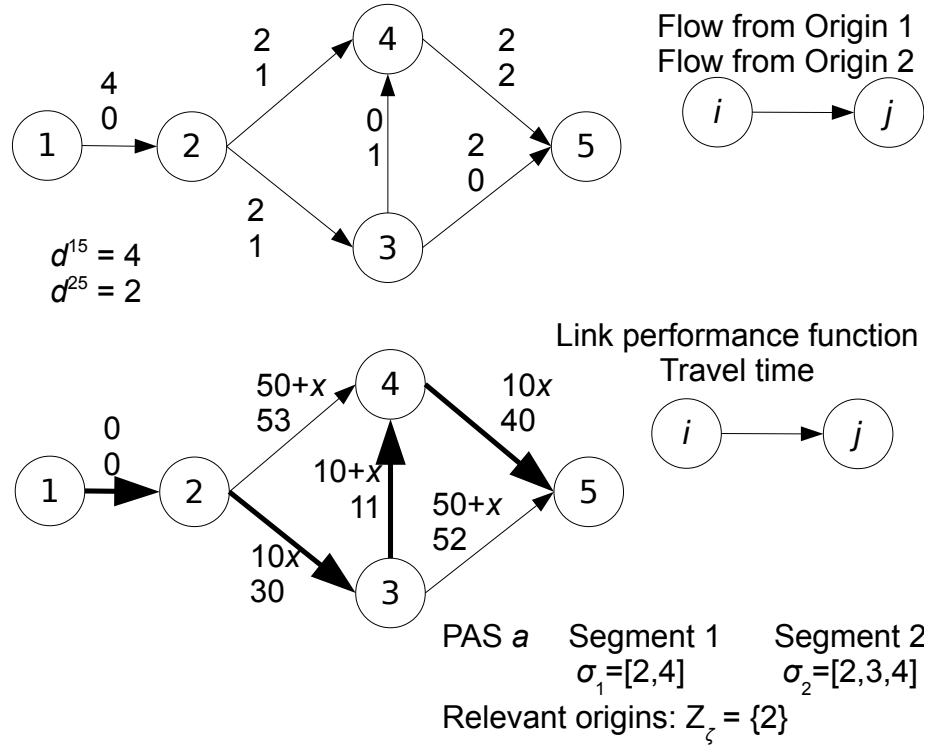


Figure 7.21: Creating a new PAS. The top panel shows the origin-specific flows, the bottom panel the link performance functions and current times. Bold links are the shortest path tree for Origin 1.

all such divergence nodes and segments, we want one for which the two path segments are short.

As an example, consider again the network from Figure 7.20. The link performance functions and current travel times are shown in Figure 7.21. The bold links show the shortest path tree rooted at Origin 1, and see that there are two links used by this origin which are not part of this tree: links (2, 4) and (3, 5). For link (2, 4), we see that the two segments of the PAS a from the previous example include a segment of links used by this origin [2, 4], and a segment of links from the shortest path tree [2, 3, 4] that have a common divergence node 2. At this point, we declare Origin 1 relevant to this PAS by adding it to the set Z_a .

For link (3, 5), we need to create a new PAS. There are two possibilities for choosing segments: one choice is [3, 5] and [3, 4, 5]; and the other choice is [2, 3, 5] and [2, 3, 4, 5] (both involve a segment of used links and a segment from the shortest path tree, starting at a common divergence node). The first one is preferred, because it has fewer links — and in fact the common link (2, 3)

in the second PAS is irrelevant, since shifting flow between segments will not change flow on such a link at all). By being shorter, there are potentially more relevant origins. If node 3 were also an origin, it could be relevant to the first choice of segments, but not the second. Therefore we create a new PAS b , and set $\sigma_b^1 = [3, 5]$, $\sigma_b^2 = [3, 4, 5]$, and $Z_b = \{1\}$. (The choice of which segment is the first and second one is arbitrary.)

Repeating the same process with the shortest path tree from Origin 2, we verify that it is relevant to PAS a (which it already is), and add it as relevant to PAS b , so $Z_a = Z_b = \{1, 2\}$. (Both origins are now relevant to both PASs).

Flow shifts

TAPAS uses flow shifts to find an equilibrium solution. There are two types of flow shifts: the most common involves shifting flow between the two segments on an existing PAS. The second involves identifying and eliminating cycles of used links for particular origins.

For the first type, assume we are given a PAS ζ , and without loss of generality assume that the current travel time on the first segment σ_1^ζ is greater than that on the second σ_2^ζ . We wish to shift flow from the first segment to the second one to either equalize their travel times, or to shift all the flow to the second path if it is still shorter. The total amount of flow we need to shift to equalize the travel times is approximately given by Newton's method:

$$\Delta h = \frac{\sum_{(i,j) \in \sigma_1^\zeta} t_{ij} - \sum_{(i,j) \in \sigma_2^\zeta} t_{ij}}{\sum_{(i,j) \in \sigma_1^\zeta} t'_{ij} + \sum_{(i,j) \in \sigma_2^\zeta} t'_{ij}}. \quad (7.87)$$

We must also determine whether such a shift is feasible (would shifting this much flow force an x_{ij}^r value to become negative?) and, unlike the algorithms earlier in this chapter, how much of this flow shift comes from each of the relevant origins in Z_ζ .

To preserve feasibility, for any relevant origin r , we must subtract the same amount from x_{ij}^r for each link in the longer segment, and add the same amount to each link in the shorter segment. Call this amount Δh^r . The non-negativity constraints require $\Delta h^r \leq \min_{(i,j) \in \sigma_1^\zeta} \{x_{ij}^r\}$; let $\overline{\Delta h}^r$ denote the right-hand side of this inequality, which must hold for every relevant origin.

If

$$\Delta h \leq \sum_{r \in Z_\zeta} \overline{\Delta h}^r \quad (7.88)$$

then the desired shift is feasible, and we choose the origin-specific shifts Δh^r to be proportional to their maximum values $\overline{\Delta h}^r$ to help maintain proportionality. If this shift is not feasible, then we shift as much as we can by setting $\Delta h^r = \overline{\Delta h}^r$ for each relevant origin.

An example of such a shift is shown in Figure 7.22, continuing the example from before. The left side of the figure shows the state of the network prior to the flow shift. The top panel shows the current origin-specific link flows;

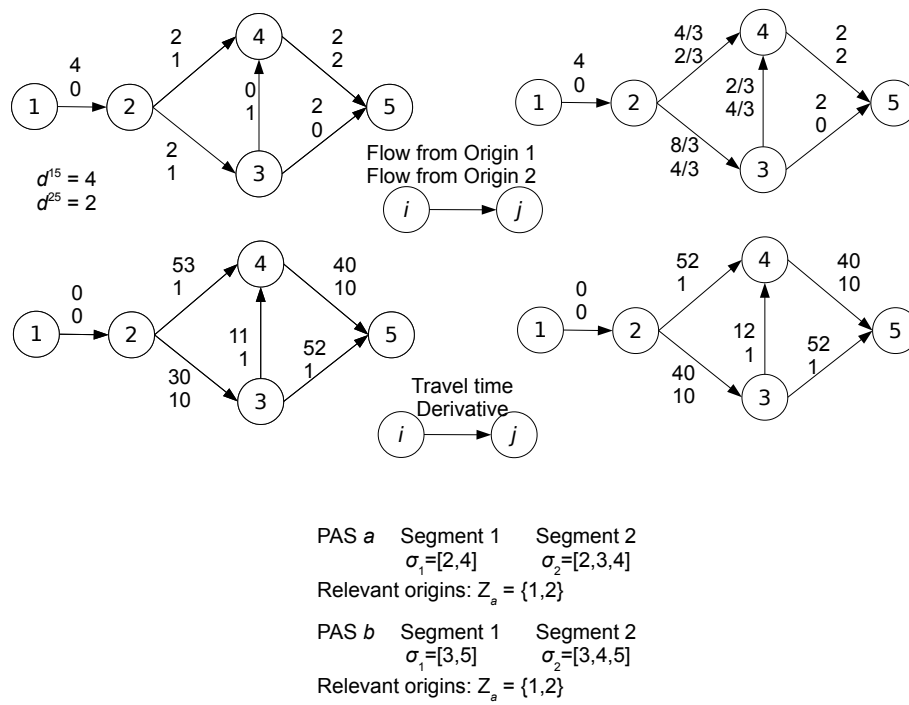


Figure 7.22: Example of flow shifts using TAPAS; left panel shows initial flows and times, right panel shows flows and times after a shift for PAS *a*.

the middle panel the current travel times and travel time derivatives; and the bottom panel shows the structure of both PASs. Starting with PAS a , we first calculate the desired total shift from equation (7.87):

$$\Delta h = \frac{53 - (30 + 11)}{2 + (10 + 1)} = 1. \quad (7.89)$$

For origin 1, we can subtract at most 2 units of flow from segment 1, and for origin 2, we can subtract at most 1. Removing any more would result in negative origin flows on link (2, 4). We thus split Δh in proportion to these maximum allowable values, yielding

$$\Delta h^1 = 2/3 \quad \Delta h^2 = 1/3 \quad (7.90)$$

and producing the solution shown in the right half of Figure 7.22. Since the link performance functions are linear, Newton's method is exact, and travel times are equal on the two segments of PAS a .

Moving to the second PAS, we see that it is at equilibrium as well, and no flow shift is done: the numerator of equation (7.87) is zero. In fact, the entire network is now at equilibrium, but the origin-based link flows do not represent a proportional solution. Proportionality adjustments are discussed below.

The second kind of flow shift involves removing flow from cycles. In TAPAS, there may be occasions where cycles are found among the links with positive flow ($x_{ij}^r > 0$). Such cycles can be detected using the topological ordering algorithm described in Section 2.2.

In such cases, we can subtract flow from every link in the cycle, maintaining feasibility and reducing the value of the Beckmann function. (See Exercise 26). Let \bar{X} denote the minimum value of x_{ij}^r among the links in such a cycle. After subtracting this amount from every x_{ij}^r in the cycle, the solution is closer to equilibrium and the cycle of positive-flow links no longer exists.

Figure 7.23 shows an example of how this might happen. This network has only a single origin, and two PASs. Applying the flow shift formula, we move 1 vehicle from segment [1, 2] to [1, 3, 2], and 1 vehicle from segment [2, 4] to segment [2, 3, 4]. This produces the flow solution in the lower-left of the figure, which contains a cycle of flow involving links (2, 3) and (3, 2). If we subtract 1 unit of flow from both of those links, we have the flow solution in the lower-right. This solution is feasible and has a lower value of the Beckmann function, as you can verify.

Proportionality adjustments

TAPAS uses proportionality adjustments to increase the entropy of the path flow solution. Note that the path flow solution is not explicitly stored, since the number of used paths can grow exponentially with network size. Rather, a path flow solution is implied by the bush, using the procedure described in Section 7.5.1, and the definitions of α and h in equations (7.85) and (7.86). That is, we calculate approach proportions α_{ij}^r using the formula

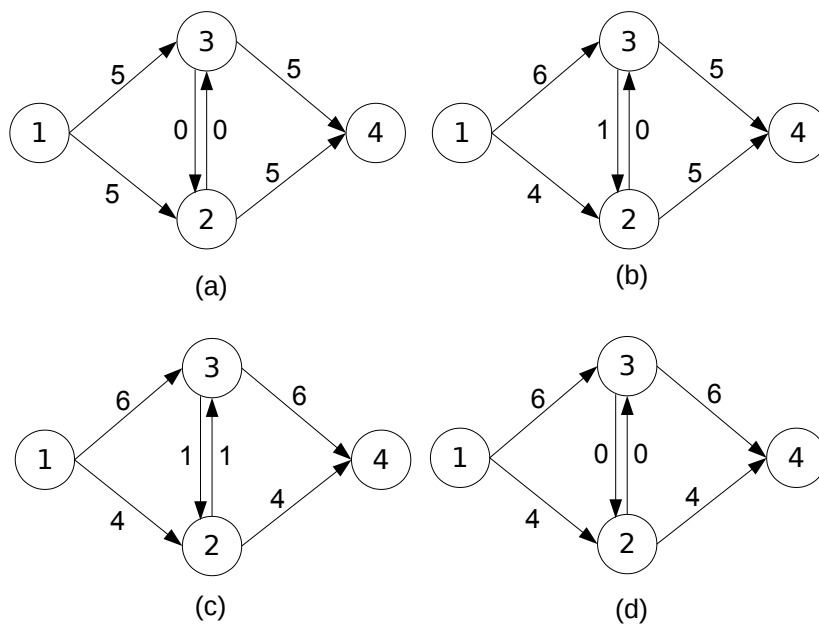


Figure 7.23: (a) Initial flow solution; (b) After a flow shift at a PAS ending at 3; (c) After a flow shift at a PAS ending at 4; (d) After removing a cycle of flow.

In a proportionality adjustment, we shift flows between segments in the PAS *without* changing the total flow on each link, so some origins will shift flow from the links in σ_1^ζ to those in σ_2^ζ , while other origins will shift flow from σ_2^ζ to σ_1^ζ . Following the previous section, we will use Δh^r to denote the amount of flow shifted from each link in σ_1^ζ to each link in σ_2^ζ , using negative numbers to indicate flow shifting from σ_2^ζ to σ_1^ζ . To maintain total link flows at their current levels, we will require

$$\sum_{r \in Z_\zeta} \Delta h^r = 0. \quad (7.91)$$

To describe the problem more formally, let b denote the node at the downstream end of the PAS, and let x_b denote the total flow through this node as in Equation (7.46). We can calculate the flow on the segments σ_1^ζ and σ_2^ζ for each relevant origin r with the formulas

$$g^r(\sigma_1^\zeta) = x_b^r \prod_{(i,j) \in \sigma_1^\zeta} \frac{x_{ij}^r}{x_j^r} \quad (7.92)$$

$$g^r(\sigma_2^\zeta) = x_b^r \prod_{(i,j) \in \sigma_2^\zeta} \frac{x_{ij}^r}{x_j^r}, \quad (7.93)$$

$$(7.94)$$

assuming positive flow through all nodes in the segment ($x_j^r > 0$).¹⁵ If proportionality were satisfied, we would have

$$\frac{g^r(\sigma_1^\zeta)}{g^r(\sigma_1^\zeta) + g^r(\sigma_2^\zeta)} = \frac{\sum_{r' \in Z_\zeta} g^{r'}(\sigma_1^\zeta)}{\sum_{r' \in Z_\zeta} (g^{r'}(\sigma_1^\zeta) + g^{r'}(\sigma_2^\zeta))} \quad (7.95)$$

for all relevant origins.

After applying segment shifts of size Δh^r , the new segment flows will be given by

$$g^r(\sigma_1^\zeta) = x_b^r \prod_{(i,j) \in \pi_1} \frac{x_{ij}^r - \Delta h^r}{x_j^r - \Delta h^r[r \neq j]} \quad (7.96)$$

$$g^r(\sigma_2^\zeta) = x_b^r \prod_{(i,j) \in \pi_2} \frac{x_{ij}^r + \Delta h^r}{x_j^r + \Delta h^r[r \neq j]}, \quad (7.97)$$

$$(7.98)$$

using brackets for an indicator function. We aim to find Δh^r values satisfying constraint (7.91) and (7.95), where the segment flows are computed with equations (7.96) and (7.97).

¹⁵What would happen if this were not true?

Solving this optimization problem exactly is a bit difficult because equations (7.85) and (7.86) are nonlinear. A good approximation method is developed in Exercise 27. A simpler heuristic is to adapt the “proportionality between origins” technique from Section 7.5.1 and approximate the (nonlinear) formulas (7.96) and (7.97) by the (linear) formulas

$$g^r(\sigma_1^\zeta) \approx g_0^r(\sigma_1^\zeta) - \Delta h^r \quad (7.99)$$

$$g^r(\sigma_2^\zeta) \approx g_0^r(\sigma_2^\zeta) + \Delta h^r \quad (7.100)$$

where $g_0^r(\sigma_1^\zeta)$ and $g_0^r(\sigma_2^\zeta)$ are the *current* segment flows (with zero shift).

Substituting equations (7.99) and (7.100) into (7.95) and simplifying, we obtain

$$\Delta h^r = g_0^r(\sigma_1^\zeta) - (g_0^r(\sigma_1^\zeta) + g_0^r(\sigma_2^\zeta)) \frac{\sum_{r' \in Z_\zeta} g_0^{r'}(\sigma_1^\zeta)}{\sum_{r' \in Z_\zeta} (g_0^{r'}(\sigma_1^\zeta) + g_0^{r'}(\sigma_2^\zeta))} \quad (7.101)$$

as a flow shift to heuristically move toward proportionality.

This heuristic is in fact exact if the PAS is “isolated” in the sense that flow does not enter or leave the segments in the middle, so $\alpha_{ij}^r = 1$ for all links in σ_1^ζ and σ_2^ζ except for the last links of each segment.

To illustrate how this procedure works, consider the example in Figure 7.24. In this example, there are multiple destinations in addition to multiple origins. The figure shows the path flow solution implied by the link flow solution at the left. We emphasize that TAPAS does not maintain the path flow solution explicitly, and the path flows are constructed from the link flows using equation (7.85). For instance, the flow on path $[1, 3, 4]$ from Origin 1 is calculated as $3 \times \frac{1}{3} \times \frac{1}{1} = 1$. This solution does not satisfy proportionality. All the vehicles from Origin 1 passing between nodes 2 and 4 use segment $[2, 4]$, while all of those from Origin 2 passing between these nodes use segment $[2, 3, 4]$. Since the total flow on the two segments are equal (two vehicles on each), flow from both origins should split equally between the two segments.

Applying equation (7.101) gives the shifts shown in Figure 7.25. This figure also shows the new origin-specific link flows and implied path flows. In this case, the link flows on the segments of PAS *a* now satisfy proportionality. This is a case where the PAS is isolated, because no vehicles entered and left the segments in the middle, and the heuristic formula (7.101) is exact.

To show how the formula is inexact for a non-isolated PAS, consider the modification of this example shown in Figure 7.26. The only change is that node 3 is now a destination for Origin 1, with a demand of 1, and that as a result the flow from Origin 1 on link (1,3) is increased by one vehicle. Repeating the same process as above, and applying (7.101), we again one swap one vehicle between each pair of segments. This results in the situation in Figure 7.27. The origin-specific link flows are the same as in Figure 7.25, except for the additional vehicle from Origin 1 on link (1,3). But the implied path flows are quite different! This is because the additional vehicle shifted onto link (3,4) was

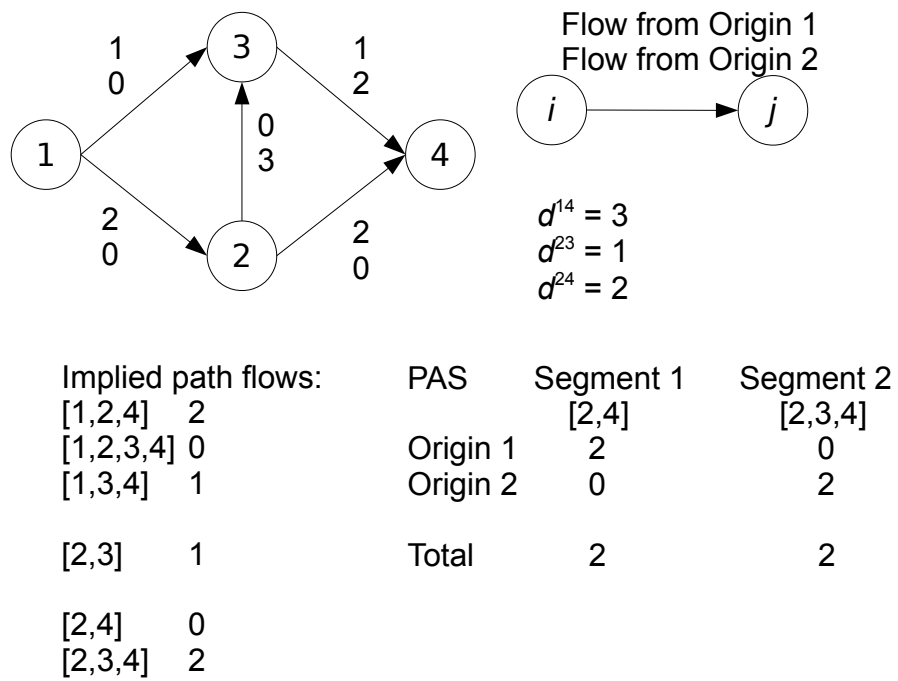
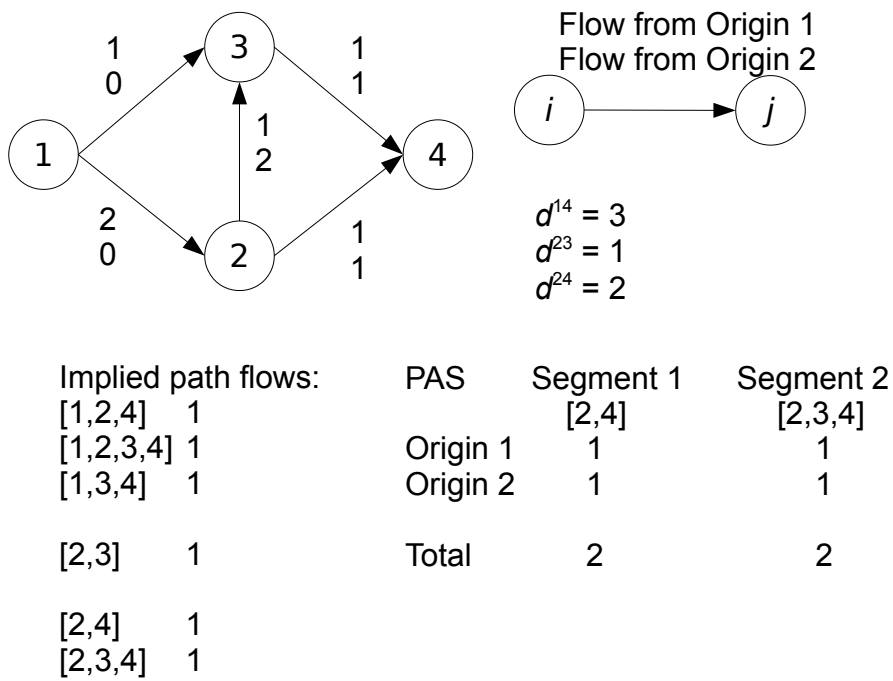


Figure 7.24: Example of a proportionality adjustment in TAPAS.

Figure 7.25: Updated flows after a proportionality adjustment for PAS a .

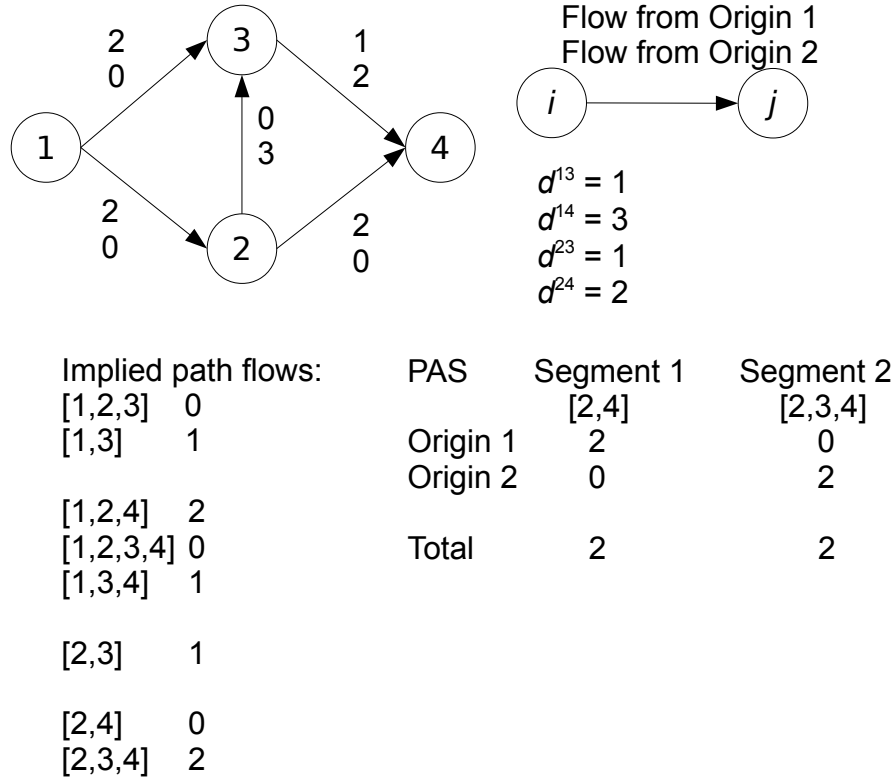


Figure 7.26: Example of a proportionality adjustment for a non-isolated PAS.

“split” between incoming links (1,3) and (2,3), according to equation (7.85), rather than allocated solely to (2,3). (Again, TAPAS does not store the flows on individual paths, and must calculate them implicitly using this formula.) As a result, proportionality is still not satisfied: between segments [2,4] and [2,3,4], Origin 1 splits in the ratio of 3:2, whereas Origin 2 splits in the ratio 1:1. This is closer to proportionality from before, but not exact. Repeated applications of the heuristic shift formula will converge to a proportional solution, in this case.

7.6 Exercises

1. [32] One critique of the BPR link performance function is that it allows link flows to exceed capacity. An alternative link performance “function” is $t_{ij} = t_{ij}^0 / (u_{ij} - x_{ij})$ if $x_{ij} < u_{ij}$, and ∞ otherwise, where t_{ij}^0 and u_{ij} are the free-flow time and capacity of link (i, j) . First show that $t_{ij} \rightarrow \infty$ as $x_{ij} \rightarrow u_{ij}$. How would using this kind of link performance function affect the solution algorithms discussed in this chapter?

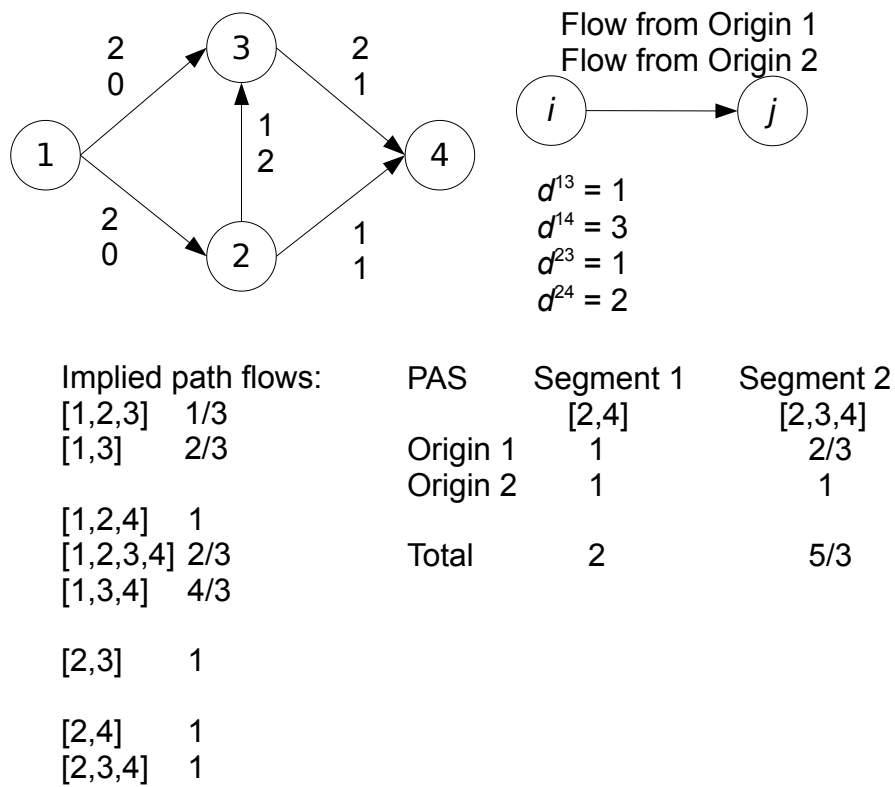


Figure 7.27: The heuristic formula is not always exact for a non-isolated PAS.

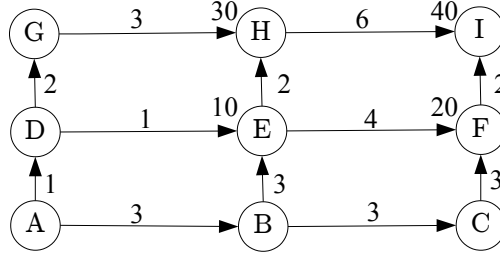


Figure 7.28: Network for Exercise 6.

2. [33] Show that the relative gap γ_1 and average excess cost are always nonnegative, and equal to zero if and only if the link or path flows satisfy the principle of user equilibrium.
3. [61] Some relative gap definitions require a lower bound \underline{f} on the value of the Beckmann function at optimality. Let \mathbf{x} denote the current solution, $f(\mathbf{x})$ the value of the Beckmann function at the current solution, and $TSTT(\mathbf{x})$ and $SPTT(\mathbf{x})$ the total system travel time and shortest path travel time at the current solution, respectively. Show that $f(\mathbf{x}) + SPTT(\mathbf{x}) - TSTT(\mathbf{x})$ is a lower bound on the Beckmann function at user equilibrium.
4. [10] What is the value of the lower bound $\underline{f} = f(\mathbf{x}) + SPTT(\mathbf{x}) - TSTT(\mathbf{x})$ if \mathbf{x} satisfies the principle of user equilibrium?
5. [23] Let (\mathbf{h}, \mathbf{x}) and (\mathbf{g}, \mathbf{y}) be two feasible solutions to the Beckmann formulation (7.3)–(7.6), and let $\lambda \in [0, 1]$. Show that $(\lambda\mathbf{h} + (1-\lambda)\mathbf{g}, \lambda\mathbf{x} + (1-\lambda)\mathbf{y})$ is also feasible, directly from the constraints (without appealing to convexity.)
6. [35] In the network in Figure 7.28, all trips originate at node A. The links are labeled with the current travel times, and the nodes are labeled with the number of trips whose destination is that node.
 - (a) Find the shortest paths from node A to all other nodes, and report the cost and backnode labels upon termination.
 - (b) What would be the target link flow solution $\hat{\mathbf{x}}$ in the method of successive averages or the Frank-Wolfe algorithm?
7. [73] All-or-nothing assignments $\hat{\mathbf{x}}$ play a major role in link-based algorithms. A naïve way to calculate these is to start with zero flows on each link; then find the shortest path from each origin r to each destination s ; then add d^rs to each link in this path. This method may require adding up to $|Z|^2$ terms for each link, in case every shortest path uses the same link.

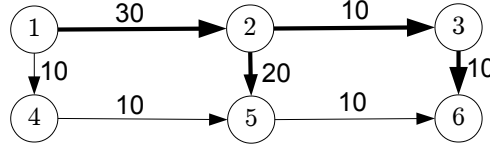


Figure 7.29: Network for Exercise 8, boldface links indicate previous $\hat{\mathbf{x}}$ target.

Formulate a more efficient algorithm which requires solving one shortest path problem per origin, and which requires adding no more than $|Z|$ terms for each link. (*Hint: Do not wait until the end to calculate $\hat{\mathbf{x}}$ and find a way to build \mathbf{x} as you go.*)

8. [42] Consider the network in Figure 7.29, with a single origin and two destinations. Each link has the link performance function $10 + x^2$, and the boldface links indicate the links used in the previous $\hat{\mathbf{x}}$ target. Report the new target link flows $\hat{\mathbf{x}}$, the step size λ , and the new resulting link flows, according to (a) Frank-Wolfe and (b) conjugate Frank-Wolfe.
9. [47] Consider the network in Figure 7.30, where 8 vehicles travel from node 1 to node 4. Each link is labeled with its delay function. For each of the algorithms listed below, report the resulting link flows, average excess cost, and value of the Beckmann function.
 - (a) Perform three iterations of the method of successive averages.
 - (b) Perform three iterations of the Frank-Wolfe algorithm.
 - (c) Perform three iterations of conjugate Frank-Wolfe.
 - (d) Perform three iterations of projected gradient.
 - (e) Perform three iterations of gradient projection.
 - (f) Perform three iterations of Algorithm B (for each iteration, do one flow update and one bush update)
 - (g) Perform three iterations of origin-based assignment (for each iteration, do one flow update and one bush update)
 - (h) Perform three iterations of linear user cost equilibrium (for each iteration, do one flow update and one bush update)
 - (i) Compare and discuss the performance of these algorithms.
10. [48] Consider the network in Figure 7.31. The cost function on the light links is $3 + (x_a/200)^2$, and the delay function on the thick links is $5 + (x_a/100)^2$. 1000 vehicles are traveling from node 1 to 9, and 1000 vehicles from node 4 to node 9. For each of the algorithms listed below, report the resulting link flows, average excess cost, and value of the Beckmann function.

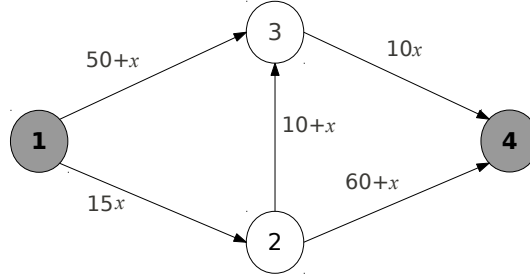


Figure 7.30: Network for Exercise 9.

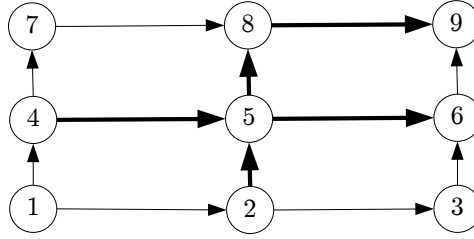


Figure 7.31: Network for Exercise 10

- (a) Perform three iterations of the method of successive averages.
 - (b) Perform three iterations of the Frank-Wolfe algorithm.
 - (c) Perform three iterations of conjugate Frank-Wolfe.
 - (d) Perform three iterations of projected gradient.
 - (e) Perform three iterations of gradient projection.
 - (f) Perform three iterations of Algorithm B (for each iteration, do one flow update and one bush update)
 - (g) Perform three iterations of origin-based assignment (for each iteration, do one flow update and one bush update)
 - (h) Perform three iterations of linear user cost equilibrium (for each iteration, do one flow update and one bush update)
 - (i) Compare and discuss the performance of these algorithms.
11. [43] The method of successive averages, as presented in the text, uses the step size $\lambda_i = 1/(i + 1)$ at iteration i . Other choices of step size can be used, and Exercise 12 shows that the algorithm converges whenever $\lambda_i \in [0, 1]$, $\sum \lambda_i = \infty$ and $\sum \lambda_i^2 < \infty$. Which of the following step size choices guarantee convergence?
- (a) $\lambda_i = 1/(i + 2)$

- (b) $\lambda_i = 4/(i+2)$
- (c) $\lambda_i = 1/i^2$
- (d) $\lambda_i = 1/(\log i)$
- (e) $\lambda_i = 1/\sqrt{i}$
- (f) $\lambda_i = 1/i^{2/3}$

12. [65] (*Proof of convergence for the method of successive averages.*) Consider the method of successive averages applied to the vector of link flows. This produces a sequence of link flow vectors $\mathbf{x}^1, \mathbf{x}^2, \mathbf{x}^3, \dots$ where \mathbf{x}^i is the vector of link flows at iteration i . We can also write down the sequence f^1, f^2, f^3, \dots of the values taken by the Beckmann function for $\mathbf{x}^1, \mathbf{x}^2$, etc. To show that this algorithm converges to the optimal solution, we have to show that either $\mathbf{x}^i \rightarrow \hat{\mathbf{x}}$ or $f^i \rightarrow \hat{f}$ as $i \rightarrow \infty$, where $\hat{\mathbf{x}}$ is the user equilibrium solution and \hat{f} the associated value of the Beckmann function. This exercise walks through one proof of this fact, for any version of the method of successive averages for which $\lambda_i \in [0, 1]$, $\sum \lambda_i = \infty$ and $\sum \lambda_i^2 < \infty$.

- (a) Assuming that the link performance functions are differentiable, show that for any feasible \mathbf{x} and \mathbf{y} there exists $\theta \in [0, 1]$ such that

$$f(\mathbf{y}) = f(\mathbf{x}) + \sum_{(i,j) \in A} t_{ij}(x_{ij})(y_{ij} - x_{ij}) + \frac{1}{2} \sum_{(i,j) \in A} t'_{ij}((1-\theta)x_{ij} + \theta y_{ij})(y_{ij} - x_{ij})^2. \quad (7.102)$$

- (b) Setting $\mathbf{x} = \mathbf{x}^i$ and $\mathbf{y} = \mathbf{x}^{i+1}$, recast equation (7.102) into an expression for the difference in the values of the Beckmann function between two consecutive iterations of MSA, in terms of λ_i and \mathbf{x}_i^* .
 - (c) Sum the resulting equation over an infinite number of iterations to obtain a formula for the limiting value \hat{f} of the sequence f^1, f^2, \dots
 - (d) Use the facts that $\sum \lambda_i = \infty$, $\sum \lambda_i^2 < \infty$, and that \hat{f} has a finite value to show that the limiting values of $SPTT(\mathbf{x}^i)$ and $TSTT(\mathbf{x}^i)$ must be equal, implying that the limit point is a user equilibrium.
13. [34] The derivation leading to (7.16) assumed that the solution to the restricted VI was not at the endpoints $\lambda = 0$ or $\lambda = 1$. Show that if you are solving (7.16) using either the bisection method from Section 4.4.2, or Newton's method (with a "projection" step ensuring $\lambda \in [0, 1]$), you will obtain the correct solution to the restricted VI even if it is at an endpoint.
14. [55] (*Linking Frank-Wolfe to optimization.*) At some point in the Frank-Wolfe algorithm, assume that the current link flows are \mathbf{x} and the target

link flows \mathbf{x}^* have just be found, and we need to find new flows $\mathbf{x}'(\lambda) = \lambda\mathbf{x}^* + (1 - \lambda)\mathbf{x}$ for some $\lambda \in [0, 1]$. Let $z(\lambda)$ be the value of the Beckmann function at $\mathbf{x}'(\lambda)$.

- (a) Using the multi-variable chain rule, we can show that z is differentiable and $z'(\lambda)$ is the dot product of the gradient of the Beckmann function evaluated at $\mathbf{x}'(\lambda)$ and the direction $\mathbf{x}^* - \mathbf{x}$. Calculate the gradient of the Beckmann function and use this to write out a formula for $z'(\lambda)$.
 - (b) Is z a convex function of λ ?
 - (c) Show that $z'(0) = 0$ only if \mathbf{x} is an equilibrium, and that otherwise $z'(0) < 0$.
 - (d) Assume that the solution of the restricted variational inequality in the Frank-Wolfe algorithm is for an “interior” point $\lambda^* \in (0, 1)$. Show that $z'(\lambda^*) = 0$.
 - (e) Combine the previous answers to show that the Beckmann function never increases after an iteration of the Frank-Wolfe algorithm (and always decreases strictly if not at an equilibrium).
15. [74] (*Proof of convergence for Frank-Wolfe.*) Exercise 14 shows that the sequence of Beckmann function values f^1, f^2, \dots from subsequent iterations of Frank-Wolfe is nonincreasing. Starting from this point, show that this sequence has a limit, and that the resulting limit corresponds to the global minimum of the Beckmann function (demonstrating convergence to equilibrium.) Your solution may require knowledge of real analysis.
16. [11] When is the Beckmann function quadratic?
17. [33] Identify conjugate directions for the following quadratic programs:
- (a) $f(x_1, x_2) = x_1^2 + x_2^2$
 - (b) $f(x_1, x_2) = x_1^2 + x_2^2 + \frac{2}{3}x_1x_2$
 - (c) $f(x_1, x_2) = 2x_1^2 + 3x_2^2 + \frac{1}{9}x_1x_2$
18. [38] The *biconjugate* Frank-Wolfe method chooses a target vector \mathbf{x}^* so that the search direction \mathbf{x}^* is conjugate to the last two search directions, rather than just the last one. Let \mathbf{x}^{AON} reflect the all-or-nothing solution at the current point, \mathbf{x}_{-1}^* the target vector used at the last iteration, and \mathbf{x}_{-2}^* the target vector used two iterations ago. Also let λ_{-1} be the step size used for the last iteration, and define

$$\mu = - \frac{\sum_{(i,j) \in A} (\lambda_{-1}(x_{-1}^*)_{ij} + (1 - \lambda_{-1})(x_{-2}^*)_{ij} - x_{ij})(x_{ij}^{\text{AON}} - x_{ij})t'_{ij}}{\sum_{(i,j) \in A} (\lambda_{-1}(x_{-1}^*)_{ij} + (1 - \lambda_{-1})(x_{-2}^*)_{ij} - x_{ij})((x_{-2}^*)_{ij} - (x_{-1}^*)_{ij})t'_{ij}} \quad (7.103)$$

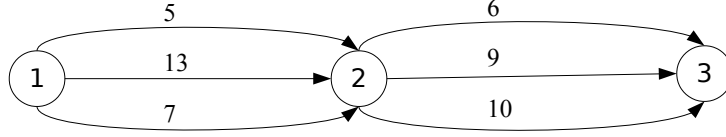


Figure 7.32: Bush for Exercise 20 with current flows.

and

$$\nu = \frac{\mu\lambda_{-1}}{1 - \lambda_{-1}} - \frac{\sum_{(i,j) \in A} ((x_{-1}^*)_{ij} - x_{ij})(x_{ij}^{AON} - x_{ij})t'_{ij}}{\sum_{(i,j) \in A} ((x_{-1}^*)_{ij} - x_{ij})^2 t'_{ij}}. \quad (7.104)$$

Show that the formula

$$\mathbf{x}^* = \frac{1}{1 + \mu + \nu} (\mathbf{x}^{AON} + \nu \mathbf{x}_{-1}^* + \mu \mathbf{x}_{-2}^*) \quad (7.105)$$

gives both a feasible target point \mathbf{x}^* , and one conjugate to the two previous search directions based on the Hessian of the Beckmann function at the current solution.

19. [41] Section 7.4.2 includes an example for Algorithm B, and Figure 7.12 shows the bush link flows at the end of a flow shifting operation. Perform a second iteration of flow shifting on the same bush, recalculating L and U labels, and scanning all nodes in reverse topological order. Report the new link flows and travel times. Also report the new bush after eliminating unused links and adding shortcuts. Does your answer depend on whether you use the L or U labels to define “shortcuts”?
20. [43] Figure 7.32 shows a bush with the current link flows labeled. Each link has delay function $5 + 2x^2$. Calculate the L , U , M , D , and α labels for all links and nodes in the bush, and identify the divergence node for each bush node.
 - (a) Perform one step of flow shifting using Algorithm B.
 - (b) Perform one step of flow shifting using OBA (starting from the original flows).
 - (c) Perform one step of flow shifting using LUCE (starting from the original flows).
 - (d) Report the maximum excess cost before and after each of these flow shifts. Which algorithm reduced this gap measure by the most?
21. [73]. This exercises walks through a proof of the formula (7.71) for choosing the threshold ϵ for finding proportional path flows. A set of paths

$\hat{\Pi} = \cup \hat{\Pi}_{rs}$ is *2-consistent* if there are no paths $\pi_1, \pi_2, \pi_3, \pi_4$ satisfying the following conditions: (1) π_1 and π_3 are in $\hat{\Pi}$; (2) at least one of π_2 and π_4 is not in $\hat{\Pi}$; (3) π_1 and π_2 connect the same OD pair; (4) π_3 and π_4 connect the same OD pair; and (5) π_1 and π_3 use the same links as π_2 and π_4 , exactly the same number of times. For instance, in Figure 7.16, the path set $\hat{\Pi} = \{[1, 2], [1, 3, 5, 6, 4, 2], [3, 5, 6, 4], [3, 5, 7, 8, 6, 4]\}$ is not 2-consistent, because we can choose $\pi_1 = [1, 3, 5, 6, 4, 2]$, $\pi_2 = [1, 3, 5, 7, 8, 6, 4, 2]$, $\pi_3 = [3, 5, 6, 4]$, and $\pi_4 = [3, 5, 7, 8, 6, 4]$. Conditions (1)–(4) are clearly satisfied. For condition (5), look at any link in the network, and count the number of times that link is used in π_1 and π_3 , and the number of times it is used in π_2 and π_4 ; every link is either unused in both pairs, used in exactly one path in both pairs, or used in both paths in each pair.

Now assume that $g_a \leq \epsilon \leq g_r/2$, as in (7.71), and choose $\hat{\Pi}$ to be all paths whose travel time is within ϵ of the shortest for its OD pair. Show that there are no four paths π_1, \dots, π_4 satisfying all of the conditions in the previous paragraph. (Hint: argue by contradiction, and apply each of the conditions, using the assumptions about the acceptance and rejection gaps to bound each path's travel time relative to the shortest path travel time.)

22. [38]. The proof of Theorem 6.4 started from the entropy-maximizing Lagrangian (6.31), which Lagrangianized both the link flow constraints (with multipliers β_{ij}) and the OD matrix constraints (with multipliers γ_{rs}). Alternatively, we can Lagrangianize only the link flow constraints, and replace (h_π/d^{rs}) with h_π (why?), giving the equation

$$\hat{\mathcal{L}}(\mathbf{x}, \boldsymbol{\beta}) = \sum_{\pi \in \Pi} h_\pi \log h_\pi + \sum_{(i,j) \in A} \beta_{ij} \left(\hat{x}_{ij} - \sum_{\pi \in \Pi} \delta_{ij}^\pi h_\pi \right). \quad (7.106)$$

Show that the gradient of this alternative Lagrangian with respect to $\boldsymbol{\beta}$ has components given by (7.83). That is, show that

$$\frac{\partial \hat{\mathcal{L}}}{\partial \beta_{ij}} = x_{ij} - \hat{x}_{ij}$$

where x_{ij} is computed from the current path flows h_π .

23. [59]. The dual algorithm step (7.83) can be compactly written as $\boldsymbol{\beta} \leftarrow \boldsymbol{\beta} + \alpha \Delta \boldsymbol{\beta}$, where $\Delta \beta_{ij} = x_{ij} - \hat{x}_{ij}$. Let $f(\alpha)$ denote the value of the alternative Lagrangian (7.106) after a step of size α is taken, and the new $\boldsymbol{\beta}$ and \mathbf{h} values are calculated. Newton's method can be used to find an α value which approximately minimizes $f(\alpha)$, maximizing entropy in the direction $\Delta \boldsymbol{\beta}$. The Newton step is $\alpha = -f'(0)/f''(0)$.

(a) Show that $f'(0) = \|\nabla_{\boldsymbol{\beta}} \hat{\mathcal{L}}\|^2$. (See Exercise 22.)

(b) Show that $f''(0) = \nabla_{\boldsymbol{\beta}}^T \hat{\mathcal{L}} H_{\boldsymbol{\beta}} \hat{\mathcal{L}} \nabla_{\boldsymbol{\beta}} \hat{\mathcal{L}}$, where $H_{\boldsymbol{\beta}}$ is the Hessian of the alternate Lagrangian with respect to $\boldsymbol{\beta}$.

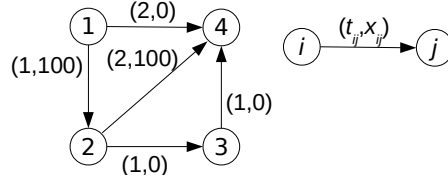


Figure 7.33: Network for Exercise 24.

(c) Show that

$$\frac{f'(0)}{f''(0)} = \frac{\sum_{(i,j) \in A} (x_{ij} - \hat{x}_{ij})^2}{\sum_{(i,j) \in A} \sum_{(k,\ell) \in A} (x_{ij} - \hat{x}_{ij})(x_{k\ell} - \hat{x}_{k\ell}) \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \delta_{k\ell}^\pi}.$$

24. [30]. In the discussion surrounding Figure 7.19, we argued that satisfying the equilibrium conditions around a “spanning” set of PASs (one for each block) was sufficient for establishing equilibrium on the entire network. Consider the network in Figure 7.33, where the demand from origin 1 to destination 4 is 100 vehicles, and there are two PASs: one between segments $[1, 4]$ and $[1, 2, 3, 4]$, and another between segments $[2, 4]$ and $[2, 3, 4]$. These are spanning, in the sense that by shifting flows between these two PASs we can obtain any feasible path flow solution from any other. They also satisfy the equilibrium conditions: for the first PAS, because there is no flow on either segment¹⁶; for the second, because the travel times are equal on the two segments. Yet the network is not at equilibrium, since $[1, 4]$ is the only shortest path and it is unused. Explain this apparent inconsistency.
25. [62]. Develop one or more algorithms to find “short” segments when generating a new PAS. These methods should require a number of steps that grows at most linearly with network size.
26. [22]. Show that the cycle-removing procedure described in the TAPAS algorithm maintains feasibility of the solution (flow conservation at each node, and non-negativity of link flows), and that the Beckmann function decreases strictly (assuming link performance functions are positive).
27. [68]. This exercise develops a technique for approximately solving equations (7.91) and (7.95), better than the heuristic given in the text.
- (a) Define $\Delta h^r(\rho)$ to be the amount of flow that needs to be shifted from origin r ’s flows on σ_1 to σ_2 , to adjust the proportion $g^r(\sigma_1)^\zeta / (g^r(\sigma_1)^\zeta + g^r(\sigma_2)^\zeta)$ to be exactly ρ . A negative value of this function indicates

¹⁶There is flow on link $(1,2)$, but not on the entire segment $[1, 2, 3, 4]$.

shifting flow in the reverse direction, from σ_2 to σ_1 . Show that this function is defined over $[0, 1]$, and its range is $[\underline{\Delta}, \overline{\Delta}]$, where $\underline{\Delta} = -\min_{(i,j) \in \sigma_2} x_{ij}^r$ and $\overline{\Delta} = \min_{(i,j) \in \sigma_1} x_{ij}^r$. Furthermore show that $\Delta h^r(0) = \underline{\Delta}$ and $\Delta h^r(1) = \overline{\Delta}$.

- (b) Show that equation (7.95) is satisfied if the flow shift $\Delta h^r(\rho)$ is applied to all relevant origins $r \in Z_\zeta$.
- (c) Therefore, it is enough to find a value of ρ for which

$$U(\rho) \equiv \sum_{r \in Z_\zeta} \Delta h^r(\rho) = 0,$$

in order to satisfy (7.91). The function U is continuous and defined on the interval $[0, 1]$. Show that $U(0) \leq 0$ and $U(1) \geq 0$, ensuring that a zero exists in this interval.

- (d) Develop a quadratic approximation for $\Delta h^r(\rho)$ based on three known points: $\Delta h^r(0) = \underline{\Delta}$, $\Delta h^r(1) = \overline{\Delta}$, and $\Delta h^r(\rho_r) = 0$, where ρ_r is the current proportion $g^r(\sigma_1)^\zeta / (g^r(\sigma_1)^\zeta + g^r(\sigma_2)^\zeta)$.
- (e) By summing these, develop a quadratic approximation for $U(\rho)$, and give an explicit formula for its root.

Chapter 8

Sensitivity Analysis and Applications

This chapter shows how a sensitivity analysis can be conducted for the traffic assignment problem (TAP), identifying how the equilibrium assignment will change if the problem parameters (such as the OD matrix or link performance functions) are changed. This type of analysis is useful in many ways: it can be used to determine the extent to which errors or uncertainty in the input data create errors in the output data. It can be used as a component in so-called “bilevel” optimization problems, where we seek to optimize some objective function while enforcing that the traffic flows remain at equilibrium. This occurs most often in the network design problem, where one must determine how to improve network links to reduce total costs, and in the OD matrix estimation problem, where one attempts to infer the OD matrix from link flows, or improve upon an existing estimate of the OD matrix.

After exploring the sensitivity analysis problem using the familiar Braess network, the first objective in the chapter is calculating derivatives of the equilibrium link flows with respect to elements in the OD matrix. It turns out that this essentially amounts to solving another, easier, traffic assignment problem with different link performance functions and constraints. The remainder of the chapter shows how these derivatives can be used in the network design and OD matrix estimation problems, which are classic transportation examples of bilevel programs.

8.1 Sensitivity Analysis Preliminaries

Figure 8.1 shows the Braess network. When this network was first introduced, the demand between node 1 and node 4 was $d_{14} = 6$, and the equilibrium solution was found to be $x_{13} = x_{23} = x_{24} = 2$ and $x_{12} = x_{34} = 4$, with a travel time of 92 minutes on all three paths. What if, instead, the demand d_{14} took another value? Figure 8.2 presents four plots showing how the equilibrium

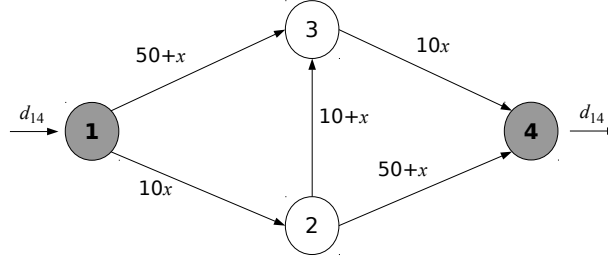


Figure 8.1: Braess network with varying demand from 1 to 4.

solution varies according to the demand level. Panel (a) shows the flows x_{12} and x_{34} , panel (b) shows the flows x_{13} and x_{24} , panel (c) shows the flow x_{23} , and panel (d) shows the shortest path travel time between nodes 1 and 4, at the corresponding equilibrium solution. You can check that when $d_{14} = 6$, the original equilibrium solution is shown in this figure.

Instead of the OD matrix, we also could have changed the link performance functions in the network. Now assume that d_{14} is fixed at its original value of 6, but that the link performance function on link (2,3) can vary. Let $t_{23}(x_{23}) = y + x_{23}$, where y is the free-flow time, resulting in the network shown in Figure 8.3. In the base solution $y = 10$, but conceivably the “free-flow time” could be changed. If the speed limit were increased, y would be lower; if traffic calming were implemented, y would be higher. In an extreme case, if the link were closed entirely you could imagine y takes an extremely large value, large enough that no traveler would use the path. One can also effectively decrease y by providing incentives for traveling on this link (a direct monetary payment, a discount at an affiliated retailer, etc.), and conceivably this incentive could be so large that y is negative. The resulting sensitivity analysis is provided in Figure 8.4

Examining the plots in Figure 8.2 and 8.4, we see that the relationships between the equilibrium solution (link flows and travel times) and the demand or free-flow time are all piecewise linear. Each “piece” of these piecewise linear functions corresponds to a particular subset of the paths being used — for instance, in Figure 8.2, when the demand is lowest, only the middle path is used. When the demand is highest, only the two outer paths are used. When the demand is at a moderate level, all three paths are used. Within each of these regions, the relationship between the demand and the equilibrium solution is linear. These pieces meet at so-called *degenerate* solutions, where the equilibrium solution does not use all of the minimum travel-time paths. (For instance, when $d_{14} = 40/11$ the equilibrium solution requires all drivers to be assigned to the middle path, even though all three have equal travel times.)

In general networks involving nonlinear link performance functions, these relationships cannot be expected to stay linear. However, they are still defined by piecewise functions, with each piece corresponding to a certain set of paths

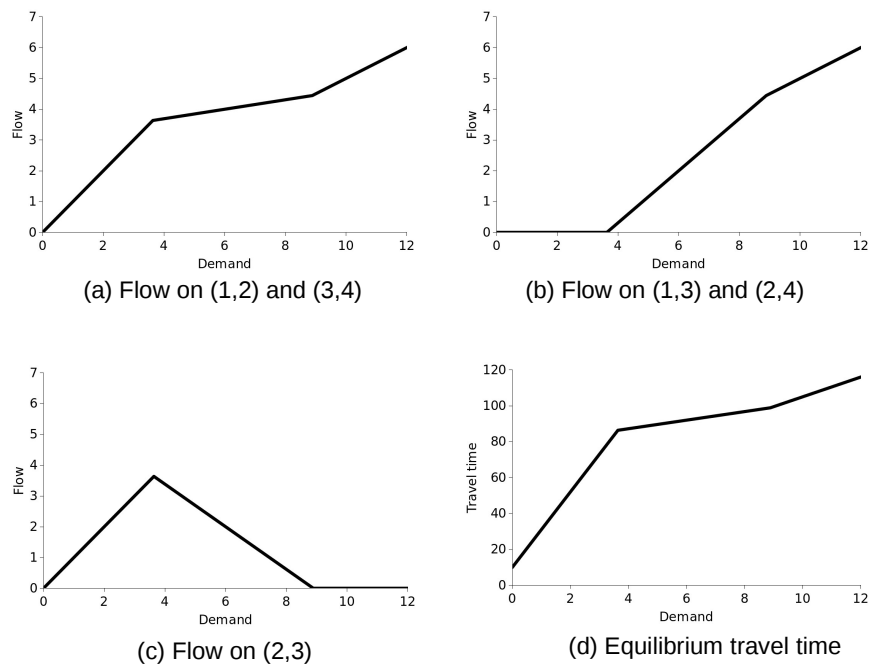


Figure 8.2: Sensitivity analysis of the Braess network to d_{14} .

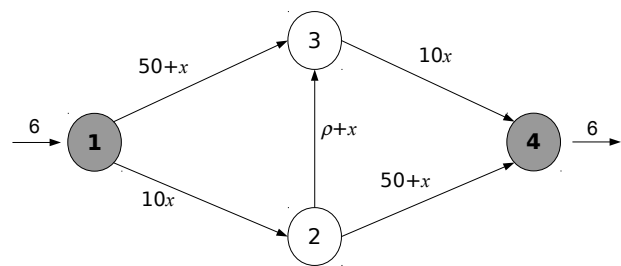
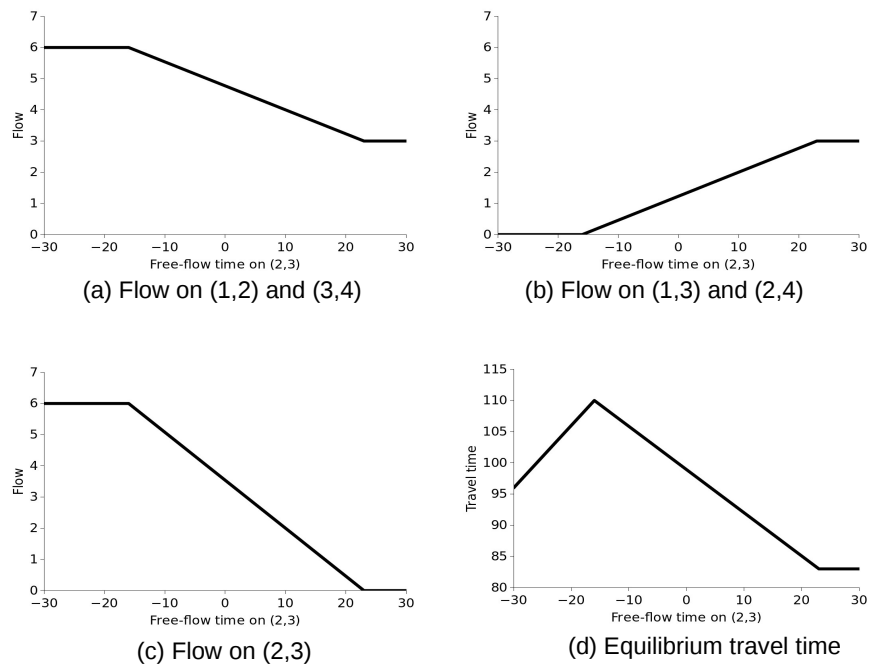


Figure 8.3: Braess network with varying free-flow time on (2,3).

Figure 8.4: Sensitivity analysis of the Braess network to y .

being used, and with the pieces meeting at degenerate solutions. The goal of the sensitivity analyses in this chapter is to identify derivatives of the equilibrium solution (link flows and travel times) at a given point. For these derivatives to be well-defined, we therefore assume that *the point at which our sensitivity analysis occurs is not degenerate*. That is, all minimum-travel time paths have positive flow. This assumption is not too restrictive, because there are only a finite number of degenerate points; for instance, if we pick the demand value at random, the probability of ending up at a degenerate point is zero.

This sensitivity analysis is still local, because the information provided by a derivative grows smaller as we move farther away from the point where the derivative is taken. For a piecewise function, the derivative provides no information whatsoever for pieces other than the one where the derivative was taken.

In this chapter, we show how this kind of sensitivity analysis can be used in two different ways. In the network design problem, this type of sensitivity analysis can be used to determine where network investments are most valuable. In Figure 8.4, the fact that the equilibrium travel time increases when y decreases (around the base solution $y = 10$) highlights the Braess paradox: investing money to improve this link will actually increase travel times throughout the network. If we were to conduct a similar analysis for other links in the network, we would see that the equilibrium travel time would decrease with improvements to the link. In the OD matrix estimation problem, we can use this sensitivity analysis to help calibrate an OD matrix to given conditions.

8.2 Calculating Sensitivities

This section derives sensitivity formulas showing the derivative of the equilibrium link flows and travel times with respect to two parameters: (1) a change in an entry of the OD matrix d^{rs} , and (2) a change to a parameter in the link performance functions (such as the free-flow time or capacity in a BPR function). In this section, assume that we are given some initial OD matrix or link performance functions, and the corresponding equilibrium solution. For our purposes, it will be most convenient if this equilibrium solution is expressed in bush-based form, that is, with vectors $\hat{\mathbf{x}}^r$ showing the flow on each link corresponding to each origin r . In this case, the non-degenerate condition requires that unused links are not part of the equilibrium bushes, that is, if $\hat{x}_{ij}^r = 0$ for any link (i, j) and any origin r , then link (i, j) does not correspond to any shortest path starting from node r — in terms of the distance labels L_i^r , we have $\hat{x}_{ij}^r > 0$ if and only if $L_j^r = L_i^r + t_{ij}$. Let $\mathcal{B}^r = \{(i, j) \in A : \hat{x}_{ij}^r > 0\}$ denote the equilibrium bush for origin r .

The non-degeneracy assumption is important, because one can show that *if the change to the OD matrix or link performance functions is small and the original equilibrium solution is non-degenerate, all of the equilibrium bushes remain unchanged*. Equivalently, even after drivers shift flows to find the new equilibrium, the set of used paths will remain the same as it was before. Fur-

thermore, one can show that the equilibrium solution is differentiable, and the derivatives of the equilibrium link flows or travel times with respect to values in the OD matrix or link performance function parameters can be interpreted as the sensitivities of the equilibrium solution.

There are several ways to calculate the values of these derivatives: historically, the first researchers used matrix-based formulas, and subsequent researchers generalized these formulas using results from the theory of variational inequalities. We adopt a different approach, using the bush-based solution representation, because it leads to an easy solution method and is fairly straightforward. This approach is based on the fact that the equilibrium solution (travel times t_{ij} and bush flows x_{ij}^r) must satisfy the following equations for each origin r :

$$L_j^r - L_i^r - t_{ij}(\hat{x}_{ij}^r) = 0 \quad \forall (i, j) \in \mathcal{B}^r \quad (8.1)$$

$$L_r^r = 0 \quad (8.2)$$

$$\sum_{(h,i) \in \Gamma^{-1}(i)} \hat{x}_{hi}^r - \sum_{(i,j) \in \Gamma(i)} \hat{x}_{ij}^r = d^{ri} \quad \forall i \in N \setminus \{r\} \quad (8.3)$$

$$\sum_{(h,r) \in \Gamma^{-1}(r)} \hat{x}_{hr}^r - \sum_{(r,j) \in \Gamma(r)} \hat{x}_{rj}^r = - \sum_{s \in Z} d^{rs} \quad (8.4)$$

$$\hat{x}_{ij}^r = 0 \quad \forall (i, j) \notin \mathcal{B}^r \quad (8.5)$$

Equations (8.1)–(8.2) reflect the equilibrium condition, and equations (8.3)–(8.4) represent flow conservation. The number of equations for each origin is no more than the sum of the number of links and nodes in the network.

Furthermore, these conditions must remain true even as the problem data (OD matrix and link performance functions) are perturbed. Since derivatives of the equilibrium solution exist under the non-degeneracy assumption, we can differentiate equations (8.1)–(8.4) to identify the relationships which must hold true among these derivatives. For brevity, in this chapter we use ξ_{ij}^r to denote the derivative of x_{ij}^r , and Λ_i to denote the derivative of L_i . These derivatives are taken with respect to either an OD matrix entry or a link performance function parameters, as described separately below.

8.2.1 Changes to the OD matrix

Assume first that we change a single entry in the OD matrix corresponding to origin \hat{r} and destination \hat{s} , so $\xi_{ij}^r = dx_{ij}^r / dd^{\hat{r}\hat{s}}$ and $\Lambda_i = dL_i / dd^{\hat{r}\hat{s}}$. Then differentiating each of equations (8.1)–(8.4) with respect to $d^{\hat{r}\hat{s}}$ gives the following

equations for each $r \in Z$:

$$\Lambda_j^r - \Lambda_i^r - t'_{ij} \sum_{r' \in Z} \xi_{ij}^{r'} = 0 \quad \forall (i, j) \in \mathcal{B}^r \quad (8.6)$$

$$\Lambda_r^r = 0 \quad (8.7)$$

$$\sum_{(h,i) \in \Gamma^{-1}(i)} \xi_{hi}^r - \sum_{(i,j) \in \Gamma(i)} \xi_{ij}^r = \begin{cases} 1 & \text{if } r = \hat{r} \text{ and } i = \hat{s} \\ 0 & \text{otherwise} \end{cases} \quad \forall i \in N \setminus r \quad (8.8)$$

$$\sum_{(h,r) \in \Gamma^{-1}(r)} \xi_{hr}^r - \sum_{(r,j) \in \Gamma(r)} \xi_{rj}^r = \begin{cases} -1 & \text{if } r = \hat{r} \\ 0 & \text{otherwise} \end{cases} \quad (8.9)$$

$$\xi_{ij}^r = 0 \quad \forall (i, j) \notin \mathcal{B}^r \quad (8.10)$$

where t'_{ij} is the derivative of the link performance function, evaluated at the current equilibrium solution $\hat{\mathbf{x}}$ (and thus treated as a constant in these equations). Equations (8.6) enforce the fact that the equilibrium bushes must remain the same. That is, the shortest path labels L_i and travel times t_{ij} must change in such a way that every link on the bush is part of a minimum travel time path to its head node. Equations (8.8) and (8.9) enforce flow conservation. For all bushes except for \hat{r} , the total flow from the origin to each destination is the same, so flow is allowed to redistribute among the bush links, but the flows starting or ending at a node cannot change. For the bush corresponding to \hat{r} , a unit increase in demand from \hat{r} to \hat{s} must be reflected by an additional vehicle leaving \hat{r} and an additional vehicle arriving at \hat{s} .

All together, the system of equations (8.6)–(8.9) involves variables Λ_i^r for each origin r and node i , and ξ_{ij}^r for each origin r and link (i, j) . Furthermore, for each origin, it contains an equation for each link and each node. Therefore, this linear system of equations can be solved to obtain the sensitivity values.¹

However, there is an easier way to solve for Λ_i^r and ξ_{ij}^r . Using the techniques in Chapter 4, you can show that the equations (8.6)–(8.9) are exactly the optimality conditions to the following minimization problem:

$$\begin{aligned} \min_{\xi^r, \Lambda^r} \quad & \frac{1}{2} \sum_{(i,j) \in A} t'_{ij} \left(\sum_{r \in Z} \xi_{ij}^r \right)^2 \\ & + \sum_{r \in Z} \sum_{i \in N} \Lambda_i^r \left(\sum_{(h,i) \in \Gamma^{-1}(i)} \xi_{hi}^r - \sum_{(i,j) \in \Gamma(i)} \xi_{ij}^r - \Delta_{ri} \right) \end{aligned} \quad (8.11)$$

$$\text{s.t.} \quad \Lambda_r^r = 0 \quad \forall r \in Z \quad (8.12)$$

$$\xi_{ij}^r = 0 \quad \forall (i, j) \notin \mathcal{B}^r \quad (8.13)$$

¹A careful reader will note that one of the flow conservation equations for each origin is redundant, but this is of no consequence to what follows.

where Δ_{ri} represents the right-hand side of equation (8.8) or (8.9), that is, $\Delta_{\hat{r}\hat{s}} = 1$, $\Delta_{\hat{r}\hat{r}} = -1$, and $\Delta_{ri} = 0$ otherwise.

This optimization problem can be put in a more convenient form by interpreting Λ_i^r as the Lagrange multiplier for the flow conservation equation corresponding to node i in bush r , and the second term in (8.11) as the Lagrangianization of this equation. Furthermore, defining $\xi_{ij} = \sum_{r \in Z} \xi_{ij}^r$, the optimization problem can be recast in the following equivalent form:

$$\min_{\xi^r} \int_0^{\xi_{ij}} t'_{ij} \xi \, d\xi \quad (8.14)$$

$$\text{s.t.} \quad \sum_{(h,i) \in \Gamma^{-1}(i)} \xi_{hi}^r - \sum_{(i,j) \in \Gamma(i)} \xi_{ij}^r = \Delta_{ri} \quad \forall i \in N, r \in Z \quad (8.15)$$

$$\xi_{ij}^r = 0 \quad \forall (i,j) \notin \mathcal{B}^r \quad (8.16)$$

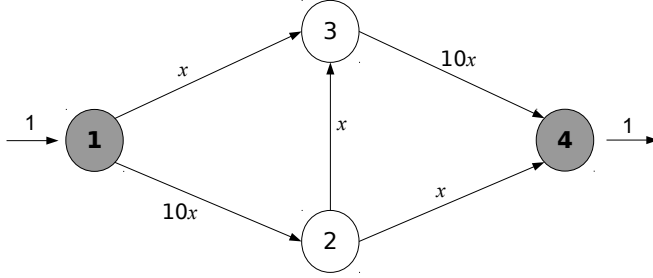
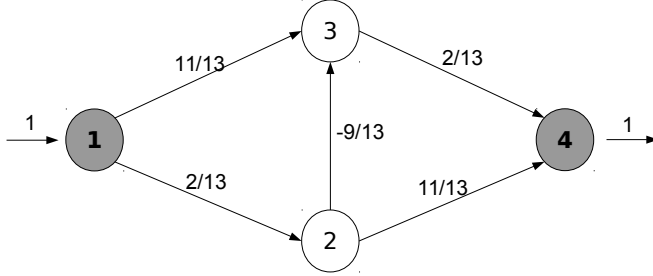
This is essentially a traffic assignment problem (TAP) in bush-based form (see Chapter 6), with the following changes:

- The original link performance functions have been replaced by *linear* link performance functions with slope equal to the derivative of the original link performance function at the original equilibrium solution. (Remember that t'_{ij} is a constant, the value of the link performance function derivative at the equilibrium solution.)
- The equilibrium bushes for each origin are *fixed* at the bushes for the original equilibrium solution.
- The only entry in the OD matrix is one unit of demand from \hat{r} to \hat{s} .
- There are no non-negativity conditions. This is because the solution variables ξ_{ij} represent changes in the original link flows, and it is possible for these changes to be negative as well as positive (cf. Figure 8.2).

If you have access to an implementation of a bush-based algorithm for solving TAP, it is easy to modify the program to take account of these distinctions, and to find the link flow sensitivities ξ_{ij}^r . From here, the values of Λ_i^r can be found by solving a shortest path problem with link travel times $\frac{dt_{ij}}{dx_{ij}} \xi_{ij}$.

As a demonstration, we use the Braess network of Figure 8.1, working around the base demand $d_{14} = 6$ and base equilibrium solution $\hat{x}_{12} = \hat{x}_{34} = 4$, $\hat{x}_{13} = \hat{x}_{23} = \hat{x}_{24} = 2$. At this level of demand, all paths are used and the equilibrium bush contains all of the links in the original network. Furthermore, at this equilibrium solution the derivatives of the link performance functions are $t'_{12} = t'_{34} = 10$ and $t'_{13} = t'_{23} = t'_{24} = 1$. The linear link performance functions based on these derivatives are shown in Figure 8.5. The OD matrix is replaced by a single unit of flow traveling from 1 (our \hat{r}) to 4 (our \hat{s}).

Solving the problem without non-negativity constraints produces the ξ_{ij} values shown in Figure 8.6. Substituting these into the link performance functions in Figure 8.5 shows that the equilibrium is satisfied: all paths have equal travel

Figure 8.5: Modified traffic assignment problem for sensitivity to d_{14} Figure 8.6: Solution to traffic assignment problem for sensitivity to d_{14}

times of $31/13$. This is exactly the slope of the piece of the equilibrium travel time (Figure 8.2d) around $d_{14} = 6$, that is, the equilibrium travel time in the sensitivity problem gives the derivative of the equilibrium travel time in the original network.

8.2.2 Changes to a link performance function

Now assume that we change a parameter y (which may represent the free-flow time, capacity, or any other parameter) in the link performance function corresponding to link (i, j) , so $\xi_{ij}^r = dx_{ij}^r/dy$ and $\Lambda_i = dL_i/dy$. This means that the link performance function t_{ij} now depends on both its flow x_{ij} , and the parameter y , so we write $t_{ij}(x, y)$. We will write $t'_{ij,x}$ to mean the partial derivative with respect to link flow², and $t'_{ij,y}$ to mean the partial derivative with respect to the improvement parameter y . Then differentiating each of equations (8.1)–(8.4)

²This is what we wrote as t'_{ij} in the previous section, when the link performance function only depended on x_{ij} .

with respect to y gives the following equations for each $r \in Z$:

$$\Lambda_j^r - \Lambda_i^r - t'_{ij,x} \sum_{r' \in Z} \xi_{ij}^{r'} - t'_{ij,y} = 0 \quad \forall (i, j) \in \mathcal{B}^r \quad (8.17)$$

$$\Lambda_r^r = 0 \quad (8.18)$$

$$\sum_{(h,i) \in \Gamma^{-1}(i)} \xi_{hi}^r - \sum_{(i,j) \in \Gamma(i)} \xi_{ij}^r = 0 \quad \forall i \in N \setminus \{r\} \quad (8.19)$$

$$\sum_{(h,r) \in \Gamma^{-1}(r)} \xi_{hr}^r - \sum_{(r,j) \in \Gamma(r)} \xi_{rj}^r = 0 \quad (8.20)$$

$$\xi_{ij}^r = 0 \quad \forall (i, j) \notin \mathcal{B}^r \quad (8.21)$$

where as before, $t'_{ij,x}$ and $t'_{ij,y}$ are evaluated at the current, equilibrium solution $\hat{\mathbf{x}}$. Equations (8.17) enforce the fact that the equilibrium bushes must remain the same, taking into account both the change in travel time on (i, j) due to the change in its link performance function as well as changes in all links' travel times from travelers shifting paths. Equations (8.19) and (8.20) enforce flow conservation. These equations are simpler than for the case of a change to the OD matrix, because the total number of vehicles on the network remains the same, and these vehicles can only shift amongst the paths in the bush. There is no change in the flow originating or terminating at any node, and the ξ variables must form a circulation.

As with a change in an OD matrix entry, the system of equations (8.17)–(8.20) is a linear system involving, for each origin, variables for each node and link. Repeating the same steps as before, this system of equations can be seen as the optimality conditions for the following optimization problem:

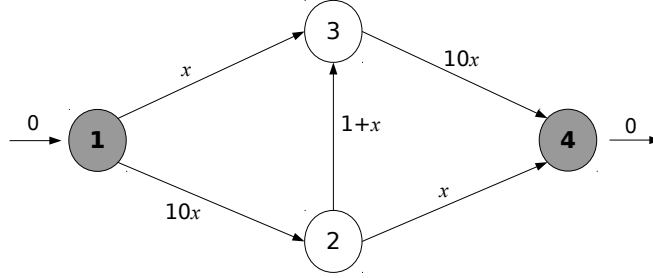
$$\min_{\xi^r} \int_0^{\xi_{ij}} (t'_{ij,x} \xi + t'_{ij,y}) d\xi \quad (8.22)$$

$$\text{s.t.} \quad \sum_{(h,i) \in \Gamma^{-1}(i)} \xi_{hi}^r - \sum_{(i,j) \in \Gamma(i)} \xi_{ij}^r = 0 \quad \forall i \in N, r \in Z \quad (8.23)$$

$$\xi_{ij}^r = 0 \quad \forall (i, j) \notin \mathcal{B}^r \quad (8.24)$$

This is essentially a traffic assignment problem in bush-based form, with the following changes:

- The original link performance functions have been replaced by *affine* link performance functions with slope equal to the derivative of the original link performance function at the original equilibrium solution, and intercept equal to the derivative of the link performance function with respect to the parameter y .
- The equilibrium bushes for each origin are *fixed* at the bushes for the original equilibrium solution.
- All entries in the OD matrix are zero.

Figure 8.7: Modified traffic assignment problem for sensitivity to y_{23}

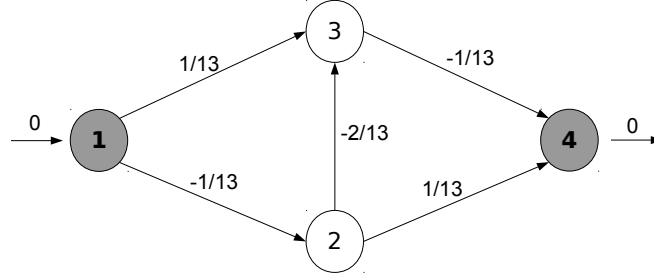
- There are no non-negativity conditions.

As a demonstration, we use the Braess network of Figure 8.3, working around the base free-flow time $y = 10$ and base equilibrium solution $\hat{x}_{12} = \hat{x}_{34} = 4$, $\hat{x}_{13} = \hat{x}_{23} = \hat{x}_{24} = 2$. That is, we replace the link performance function $t_{23}(x_{23}) = 10 + x_{23}$ with the function $t_{23}(x_{23}, y) = y + x_{23}$ and see what happens when y varies. At the base value $y = 10$, all paths are used and the equilibrium bush contains all of the links in the original network. Furthermore, at this equilibrium solution the derivatives of the link performance functions with respect to flows are $t'_{12,x} = t'_{34,x} = 10$ and $t'_{13,x} = t'_{23,x} = t'_{24,x} = 1$. For link (2,3), we add the constant term $t'_{23,y} = 1$. The link performance functions based on these derivatives are shown in Figure 8.7. The OD matrix is set equal to zero, since there is no change in the total demand through the network.

Solving the problem without non-negativity constraints produces the ξ_{ij} values shown in Figure 8.8. Substituting these into the link performance functions in Figure 8.7 shows that the equilibrium is satisfied: all paths have equal travel times of $-9/13$. Again, this is the slope of the piece of the equilibrium travel time (Figure 8.4d) around $y = 10$, that is, the equilibrium travel time in the sensitivity problem gives the derivative of the equilibrium travel time in the original network.

8.3 Network Design Problem

In the network design problem, one must determine how best to spend funds on improving links in the transportation network. This is a challenging optimization problem, and in cases of any practical interest one cannot hope to identify a globally optimal investment policy. This is mainly because *in transportation systems, the planner cannot compel travelers to choose a particular path*. Instead, after any improvement is made to the links, flows will redistribute according to the principle of user equilibrium. The goal is to find the best investment policy, knowing and anticipating how travelers will respond once the network has been

Figure 8.8: Solution to traffic assignment problem for sensitivity to y_{23}

changed.

Specifically, assume that the link performance functions for each link (i, j) now depend on the amount of money y_{ij} invested in that link (perhaps increasing its capacity through widening, or decreasing its free-flow time) as well as on the flow x_{ij} . One example of such a link performance function is

$$t_{ij}(x_{ij}, y_{ij}) = t_{ij}^0 \left(1 + \alpha \left(\frac{x_{ij}}{u_{ij} + K_{ij}y_{ij}} \right)^\beta \right) \quad (8.25)$$

where K_{ij} represents the capacity improvement if a single unit of money is invested on it.

A planning agency may have many different objectives and constraints when determining how to improve a network. This section develops and explores one specific variation of the network design problem, but there are many other variations which have been proposed in the literature. The version presented here is a fairly standard one, which can be extended in a number of different ways. In this variation, the objective of the planning agency is to minimize the total cost, given by the sum of total system travel time $TSTT$ (converted to monetary units by a conversion factor Θ , which also reflects duration of the analysis horizon and discounting) and the costs of the network improvements themselves. The optimization problem is

$$\min_{\mathbf{x}, \mathbf{y}} f(\mathbf{x}, \mathbf{y}) = \Theta \sum_{(i,j) \in A} x_{ij} t_{ij}(x_{ij}, y_{ij}) + \sum_{(i,j) \in A} y_{ij} \quad (8.26)$$

$$\text{s.t.} \quad \mathbf{x} \in \arg \min_{\mathbf{x} \in X} \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x, y_{ij}) dx \quad (8.27)$$

$$y_{ij} \geq 0 \quad \forall (i, j) \in A \quad (8.28)$$

Most of this problem is familiar: the objective (8.26) is to minimize the sum of total system travel time (converted to units of money) and construction cost, and constraint (8.28) requires that money can only be spent on links (not “recovered”

from them with a negative y_{ij} value). Also note that since $\mathbf{y} = \mathbf{0}$ is a feasible solution (corresponding to the “do-nothing” alternative), in the optimal solution to this problem the cost savings (in the form of reduced $TSTT$) must at least be equal to the construction costs, guaranteeing that the optimal investment policy has greater benefit than cost. The key equation here is (8.27), which requires that the link flows \mathbf{x} satisfy the principle of user equilibrium by minimizing the Beckmann function. In other words, *one of the constraints of the network design problem is itself an optimization problem*. This is why the network design problem is called a bilevel program. This type of problem is also known as a mathematical program with equilibrium constraints. This class of problems is extremely difficult to solve, because the feasible region is typically nonconvex.

To see why, consider two feasible solutions $(\mathbf{x}^1, \mathbf{y}^1)$ and $(\mathbf{x}^2, \mathbf{y}^2)$ to the network design problem. The link flows \mathbf{x}^1 are the equilibrium link flows under investment policy \mathbf{y}^1 , and link flows \mathbf{x}^2 are the equilibrium link flows under investment policy \mathbf{y}^2 . If the feasible region were a convex set, then any weighted average of these two solutions would themselves be feasible. Investment policy $\frac{1}{2}\mathbf{y}^1 + \frac{1}{2}\mathbf{y}^2$ still satisfies all the constraints on \mathbf{y} (all link investments are nonnegative). However, the equilibrium link flows under this policy cannot be expected to be the average of \mathbf{x}^1 and \mathbf{x}^2 , because the influence of y_{ij} on t_{ij} can be nonlinear and the sets of paths which are used in \mathbf{x}^1 and \mathbf{x}^2 can be completely different. In other words, the equilibrium link flows after averaging two investment policies need not be the average of the equilibrium link flows under those two policies separately.

Unfortunately, solving optimization problems with nonconvex feasible regions is a very difficult task. Therefore, solution methods for the network design problem are almost entirely heuristic in nature. These heuristics can take many forms; one popular approach is to adapt a metaheuristic method, such as those discussed in Section 4.5. Another approach is to develop a more tailored heuristic based on specific insights about the network design problem. This approach, being more educational, is adopted here. Specifically, we can use the sensitivity analysis from the previous sections to identify derivatives of the objective function f with respect to each link investment y_{ij} , and use this to move in a direction which reduces total cost.

Specifically, notice that constraint (8.27) actually makes \mathbf{x} a function of \mathbf{y} , since the solution to the user equilibrium problem is unique in link flows. That is, the investment policy \mathbf{y} determines the equilibrium link flows \mathbf{x} exactly. So, the objective function can be made a function of \mathbf{y} alone, written $f(\mathbf{x}(\mathbf{y}), \mathbf{y})$. The derivative of this function with respect to an improvement on any link is then

$$\frac{\partial f}{\partial y_{ij}} = \Theta \sum_{(k,\ell) \in A} \frac{\partial f}{\partial x_{k\ell}} \frac{\partial x_{k\ell}}{\partial y_{ij}} + 1 \quad (8.29)$$

or, substituting the derivative of (8.26) with respect to each link flow,

$$\frac{\partial f}{\partial y_{ij}} = \Theta \left\{ \sum_{(k,\ell) \in A} \frac{\partial x_{k\ell}}{\partial y_{ij}} \left(t_{k\ell}(x_{k\ell}, y_{k\ell}) + x_{k\ell} \frac{\partial t_{k\ell}}{\partial x_{k\ell}}(x_{k\ell}, y_{k\ell}) \right) + x_{ij} \frac{\partial t_{ij}}{\partial y_{ij}} \right\} + 1. \quad (8.30)$$

In turn, the partial derivatives $\frac{\partial x_{k\ell}}{\partial y_{ij}}$ can be identified using the technique of Section 8.2.2 as the marginal changes in link flows throughout the network when the link performance function of (i, j) is perturbed.

The vector of all the derivatives (8.30) forms the gradient of f with respect to \mathbf{y} . This gradient is the direction of steepest *ascent*, that is, the direction in which f is increasing fastest. Since we are solving a minimization problem, we should move in the opposite direction. Taking such a step, and ensuring feasibility, gives the updating equation

$$\mathbf{y} \leftarrow [\mathbf{y} - \mu \nabla_{\mathbf{y}} f]^+ \quad (8.31)$$

where μ is a step size to be determined, and the $[\cdot]^+$ operation is applied to each component of the vector. This suggests the following algorithm:

1. Initialize $\mathbf{y} \leftarrow \mathbf{0}$.
2. Calculate the link flows $\mathbf{x}(\mathbf{y})$ by solving the traffic assignment problem with link performance functions $\mathbf{t}(\mathbf{x}, \mathbf{y})$.
3. For each link (i, j) determine $\frac{\partial f}{\partial y_{ij}}$ by solving the sensitivity problem described in Section 8.2.2 and using (8.30).
4. Update \mathbf{y} using (8.31) for a suitable step size μ .
5. Test for convergence, and return to step 2 if not converged.

Two questions are how μ should be chosen in step 4, and how convergence should be tested in step 5. The difficulty in step 4 is that the derivatives provided by a sensitivity analysis are only local, and in particular if μ is large enough that (8.31) changes the set of used paths, the derivative information is meaningless. However, if μ is small enough one will see a decrease in the objective function if at all possible. So, one could start by testing a sequence of μ values (say, 1, 1/2, 1/4, ...), evaluating the resulting \mathbf{y} values, \mathbf{x} values, and f , stopping as soon as f decreases from its current value. (Note that this is a fairly computationally intensive process, since the traffic assignment problem must be solved for each μ to get the appropriate \mathbf{x} value.) Other options include using a stricter stopping criterion such as the Armijo rule, which would ensure that the decrease is “sufficiently large”; or using bisection to try to choose the value of μ which minimizes f in analogy to Frank-Wolfe. All of these methods require solving multiple traffic assignment problems at each iteration.

Regarding convergence in step 5, one can either compare the progress made in decreasing f over the last few iterations, or the changes in the investments

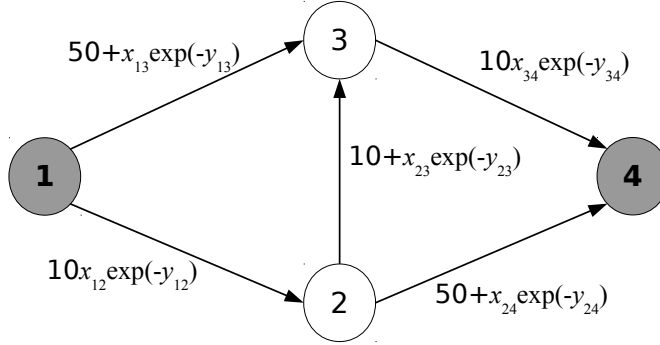


Figure 8.9: Network design problem example.

\mathbf{y} . This choice of stopping criterion is different in nature than those used in Chapter 7 for solution methods to TAP. In that case, we can prove theoretical convergence to the equilibrium solution, and we can design our stopping criteria directly on the equilibrium condition. The method described above, by contrast, is *not* proven to converge to the global optimum, and can get stuck in solutions which are locally optimal but not globally so. (This often happens in nonconvex optimization problems.) By terminating the algorithm when no more progress is made, we are checking that we have found a local optimum, but cannot guarantee that we have found a global one.

As stated above, the network design problem does not have a convex feasible region, and even if one were to eliminate \mathbf{x} by writing \mathbf{x} as a function of \mathbf{y} , the resulting function can be shown to be nonconvex and have multiple local optimal solutions. Therefore, the method described above is not guaranteed to converge to a global optimum solution. This is perhaps a bit disappointing; but, as stated above, at present there is no way to ensure global optimality within a reasonable amount of computation time. Therefore, alternative heuristics can only be compared by the quality of solutions obtained for a given quantity of computational effort.

For an example, consider the Braess network shown in Figure 8.9, where the link performance functions are shown and $\Theta = \frac{1}{20}$. Notice now that the coefficients of x_{ij} in the link performance functions now depend on the amount of money y_{ij} invested as well, with decreasing marginal returns as more money is spent. When $\mathbf{y} = \mathbf{0}$, the solution to the user equilibrium problem $\mathbf{x}(\mathbf{y})$ is the now-familiar solution to the Braess network, which divides flow evenly among all three paths in the network. The total system travel time is 552 vehicle-minutes, so the value of the objective (8.26) function is $\frac{1}{20}(552) + 0 = 27.6$. This completes the first two steps of the algorithm.

For the third step, we must solve five sensitivity problems, one for each link, to determine $\frac{\partial f}{\partial y_{ij}}$. These five problems are shown in Figure 8.10, where we have substituted the initial values $\mathbf{y} = \mathbf{0}$ and the equilibrium link flows \mathbf{x} . In all of

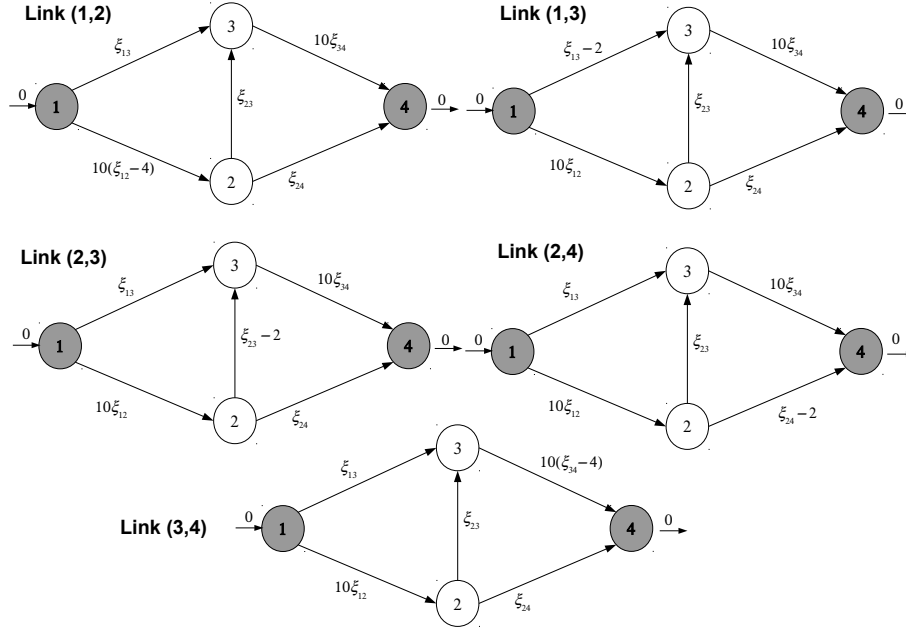


Figure 8.10: Five sensitivity problems for network design.

these, notice that the demand is zero, and the link performance functions are affine. Within each problem, the link being improved has a slightly different link performance function, accounting for the effect of the link improvement. The other links' performance functions only reflect their change due to shifting flows. For instance, in the problem in the upper left, link (1,3) is not improved, so its link performance function is simply $\xi_{13} \frac{\partial t_{13}}{\partial x_{13}} = 10\xi_{12} \exp(-y_{13}) = 10\xi_{12}$ since $y_{13} = 0$. Since link (1,2) is being improved, in addition to the term $\xi_{12} \frac{\partial t_{12}}{\partial x_{12}} = \xi_{12} \exp(-y_{12}) = \xi_{14}$, we add the constant term $\frac{\partial t_{12}}{\partial y_{12}} = -x_{14} \exp(-y_{12}) = -4$ since $x_{12} = 4$ and $y_{12} = 0$.

The solutions (in terms of ξ_{ij}) to the five sensitivity problems are shown in Table 8.1, as can be verified by substituting these ξ values into the networks in Figure 8.10. Substituting these ξ values into equation (8.30), along with the current values of the travel times and link flows, gives the gradient

$$\nabla_{\mathbf{y}} f = \begin{bmatrix} y_{12} \\ y_{13} \\ y_{23} \\ y_{24} \\ y_{34} \end{bmatrix} = \begin{bmatrix} -0.85 \\ 0.49 \\ 1.42 \\ 0.49 \\ -0.85 \end{bmatrix} \quad (8.32)$$

Each component of the gradient shows how the objective of the network de-

Table 8.1: Solutions to Braess sensitivity problems.

Link	Sensitivity problem				
	(1,2)	(1,3)	(2,3)	(2,4)	(3,4)
(1,2)	3.36	-0.17	0.15	0.014	-0.28
(1,3)	-3.36	0.17	-0.15	-0.014	0.28
(2,3)	3.08	-0.15	0.31	-0.15	3.08
(2,4)	0.28	-0.014	-0.15	0.17	-3.36
(3,4)	-0.28	0.014	0.15	-0.17	3.36

sign problem will change if a unit of money is spent improving a particular link. In the derivative formula (8.30), the term in parentheses represents the marginal change in total system travel time, and the addition of unity at the end of the formula represents the increase in total expenditures. If the derivative (8.30) is negative, then the reduction in total system travel time from a marginal improvement in the link will outweigh the investment cost. If it lies between zero and one for a link, then a marginal investment in the link will reduce total system travel time, but the cost of the improvement will outweigh the value of the travel time savings. If it is greater than one, then total system travel time would actually increase if the link is improved, as in the Braess paradox. So, in this example, the gradient (8.32) shows that improvements on links (1,2) and (3,4) will be worthwhile; improvements on links (1,3) and (2,4) would reduce TSTT but not by enough to outweigh construction cost; and an improvement on link (2,3) would actually be counterproductive and worsen congestion.

So, proceeding to step 4 of the network design algorithm, we choose a trial value $\mu = 1$ and apply (8.31) to obtain a candidate solution where $y_{12} = y_{34} = 0.85$ and all other y_{ij} values remain at zero. Re-solving equilibrium with the new link performance functions, the new equilibrium solution is to load all vehicles on the middle path, that is, $x_{12} = x_{23} = x_{34} = 6$ and $x_{13} = x_{24} = 0$. The total system travel time is now 405 vehicle-minutes, so the objective is $\frac{1}{20}(405) + 2(0.85) = 22.0$. This reduces the objective from its current value of 27.6, so we accept the step size $\mu = 1$, and return to step 2.

Notice that solving the network design problem requires solving a very large number of traffic assignment subproblems: once for each iteration to determine \mathbf{x} ; modified sensitivity problems for each link to calculate derivatives; and again multiple times per iteration to identify μ . Solving practical problems can easily require solution of thousands or even millions of traffic assignment problems. In bilevel programs such as network design, having an efficient method for solving traffic assignment problems is critical. Path-based and bush-based algorithms can be efficiently “warm-started,” making them good choices for this application.

8.4 OD Matrix Estimation

Practical application of the traffic assignment problem requires two main inputs: information about the physical network (roadway topology, link performance functions) as well as information about travel demand patterns (the OD matrix). The former is much easier to obtain and validate. Standard link performance functions (such as the BPR function) can be used only knowing the free-flow time and capacity of links, which can be easily estimated; and in principle, this information can be directly inferred from field measurements and traffic sensors. The OD matrix is quite a bit harder to estimate, for several reasons. First, there is no practical way to directly observe the complete OD matrix; at best we can work with a sample of travelers who consent to reveal their travel patterns. Second, at least until recently, it was difficult to obtain this information without travelers explicitly reporting their origins and destinations — traffic sensors tell you what is happening on a specific link, but not where the people are coming from or where they are going. Recently, GPS and Bluetooth technologies have become more commonplace and can in principle provide origins and destinations automatically (a rough estimate could even be obtained from cellular phone traces) — but there are major privacy issues associated with using such data, as well as a number of data inference issues involved in translating these data into an OD matrix. Third, the number of entries in an OD matrix is much larger than the number of links in the network: the number of OD pairs is generally proportional to the *square* of the number of nodes, while the number of links is generally proportional to the number of nodes themselves. It is typical to see practical networks which have tens of thousands of links, but millions of OD pairs.

Therefore, it is natural to find ways to determine an OD matrix which should be used for traffic assignment. One approach, used in the field of travel demand modeling, is based on developing behavioral models of how households make travel choices. Using demographic and other features, there are models to estimate the total number of trips made by households over some period of time, their destinations, their mode choices, and so forth. Another approach is to attempt to infer an OD matrix from traffic counts on specified links in the network. Both methods have advantages and disadvantages: travel demand models provide more insights about underlying behavior (critical when developing long-range forecasts of future demand), and can directly lead to an OD matrix. Unfortunately, the data needed to calibrate demand models is more expensive and cumbersome to work with, often involving recruiting survey participants.

Traffic sensors, on the other hand, collect data automatically, inexpensively, and without privacy issues, but there is a major dimensionality issue. Since the number of OD pairs is much larger than the number of links, one cannot hope to uniquely determine the OD matrix solely from link counts; the problem is massively underdetermined. In a network where all nodes are zones, it is trivial to find an OD matrix which perfectly matches link counts, by creating an OD matrix where all trips travel from one node to an adjacent node. Even though this matrix is entirely unrealistic, there is no deviation whatsoever from

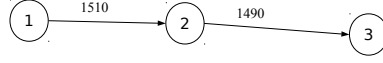


Figure 8.11: Observed link volumes from traffic sensors along a freeway.

the counts. Even worse, due to unavoidable errors in traffic count records (Figure 8.11), sometimes this trivial matrix will match counts much better than a more “realistic” matrix! In this figure, it is likely that the total flow on the freeway is approximately 1500 vehicles from node 1 to node 3; but this does not match the counts as well as 1510 vehicles from 1 to 2, and 1490 from 2 to 3. Simply matching counts does not provide the behavioral insight needed to identify what a “realistic” matrix is.

This section provides a method to reconcile both approaches. The idea is that an initial OD matrix \mathbf{d}^* is already available from a travel demand model. While this matrix is hopefully close to the true value, it also contains sampling and model estimation errors and can never be fully accurate. However, there are sensors on a subset of links $\bar{A} \subset A$, and there is a vector \mathbf{x}^* of traffic volume counts on these links. The intent is to use these traffic counts to try to improve the initial OD matrix \mathbf{d}^* . The following optimization problem expresses this:

$$\min_{\mathbf{d}, \mathbf{x}} f(\mathbf{d}, \mathbf{x}) = \Theta \sum_{(r,s) \in Z^2} (d_{rs} - d_{rs}^*)^2 + (1 - \Theta) \sum_{(i,j) \in \bar{A}} (x_{ij} - x_{ij}^*)^2 \quad (8.33)$$

$$\text{s.t.} \quad \mathbf{x} \in \arg \min_{\mathbf{x} \in X(\mathbf{d})} \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (8.34)$$

$$d_{rs} \geq 0 \quad \forall (r,s) \in Z^2 \quad (8.35)$$

where Θ is a parameter ranging from zero to one and $X(\mathbf{d})$ is the set of feasible link flows when the OD matrix is \mathbf{d} .

The objective function (8.33) is of the least-squares type, and attempts to minimize both the deviation between the final OD matrix \mathbf{d} and the initial estimate \mathbf{d}^* , and the deviation between the equilibrium link flows \mathbf{x} associated with the OD matrix \mathbf{d} , and the actual observations \mathbf{x}^* on the links with sensors. The hope is to match traffic counts reasonably well, while not wandering into completely unrealistic OD matrices. The factor Θ is used to weight the importance of matching the initial OD matrix estimate, and the link flows. It can reflect the relative degree of confidence in \mathbf{d}^* and \mathbf{x}^* ; if the travel demand was obtained from high-quality data and a large sample size, whereas the traffic count data is old and error-prone, a Θ value close to one is appropriate. Conversely, if the travel demand model is less trustworthy but the traffic count data is highly reliable, a lower Θ value is appropriate. In practice, a variety of Θ values can be chosen, and the resulting tradeoffs between matching the initial estimate and link flows can be seen.

As indicated by the constraint (8.34), this optimization problem is also a bilevel program, because the mapping from an OD matrix \mathbf{d} to the resulting

link flows \mathbf{x} involves the equilibrium process. The fact that the optimization problem is bilevel also means that we cannot expect to find the global optimum OD matrix, and that heuristics should be applied. A sensitivity-based heuristic, like the one used for the network design problem, would determine how the link flows would shift if the OD matrix is perturbed, and use this information to find a “descent direction” which would reduce the value of the objective.

Following the same technique as in Section 8.3, the constraint (8.34) defines \mathbf{x} uniquely as a function of \mathbf{d} due to the uniqueness of the link flow solution to the traffic assignment problem. (Also note that the dependence on \mathbf{d} appears through the feasible region, requiring that the minimization take place over $X(\mathbf{d})$.) So, we can rewrite the objective function as a function of \mathbf{d} alone, by defining $\mathbf{x}(\mathbf{d})$ as the equilibrium link flows in terms of the OD matrix and transforming the objective to $f(\mathbf{d}, \mathbf{x}(\mathbf{d}))$. Then, the partial derivative of the objective with respect to any entry d^{rs} in the OD matrix is given by

$$\frac{\partial f}{\partial d^{rs}} = 2\Theta(d_{rs} - d_{rs}^*) + 2(1 - \Theta) \sum_{(i,j) \in \bar{A}} (x_{rs} - x_{rs}^*) \frac{\partial x_{ij}}{\partial d^{rs}} \quad (8.36)$$

where the partial derivatives $\frac{\partial x_{ij}}{\partial d^{rs}}$ are found from the sensitivity formulas in Section 8.2.1.

The vector of all the derivatives (8.36) forms the gradient of f with respect to \mathbf{d} . So, taking a step in the opposite direction, and ensuring that the values in the OD matrix remain non-negative provide the following update rule:

$$\mathbf{d} \leftarrow [\mathbf{d} - \mu \nabla_{\mathbf{d}} f]^+ \quad (8.37)$$

where μ is a step size to be determined, and the $[\cdot]^+$ operation is applied to each component of the vector. This leads to the following algorithm for OD matrix estimation:

1. Initialize $\mathbf{d} \leftarrow \mathbf{d}^*$.
2. Calculate the link flows $\mathbf{x}(\mathbf{d})$ by solving the traffic assignment problem with OD matrix \mathbf{d} .
3. For each OD pair (r, s) determine $\frac{\partial f}{\partial d^{rs}}$ by solving the sensitivity problem in Section 8.2.1 and using (8.36).
4. Update \mathbf{d} using (8.37) for a suitable step size μ .
5. Test for convergence, and return to step 2 if not converged.

The comments about the step size μ and convergence criteria from the network design problem (Section 8.3) apply equally as well here: μ can be determined using an iterative line search procedure, choosing smaller values until the objective function decreases, and the algorithm can be terminated when it fails to make additional substantial progress in reducing the objective.

This procedure is demonstrated using the network in Figure 8.12. In this network, traffic counts are available on three of the links in the network, and

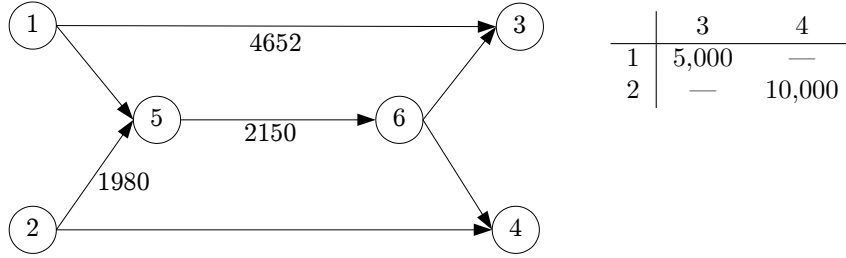


Figure 8.12: Example network for OD matrix estimation. (Observed link volumes and initial OD matrix shown.)

Table 8.2: Solutions to OD matrix estimation sensitivity problems.

Link	Sensitivity problem	
	d_{13}	d_{24}
(1,3)	0.733	0.067
(1,5)	0.267	-0.067
(2,4)	0.067	0.733
(2,5)	-0.067	0.267
(5,6)	0.200	0.200
(6,3)	0.267	-0.067
(6,4)	-0.067	0.267

an initial OD matrix is available based on a travel demand model. The link performance function on every link is $t_{ij} = 10 + x_{ij}/100$, and Θ is given as $\frac{1}{10}$. We begin by initializing the OD matrix to the initial matrix \mathbf{d}^* , and solving a traffic assignment problem.

As the reader can verify, the equilibrium link flows on the three links with available counts are $x_{13} = 4733\frac{1}{3}$, $x_{25} = 1933\frac{1}{3}$, and $x_{56} = 2200$. The first sum in the objective (8.33) gives the fit of the OD matrix to the initial matrix (zero) and the the second sum gives the fit of the equilibrium link flows to the observed link flows (11292). When weighted with $\Theta = 0.1$, this gives an initial objective value of 10164. Next, we solve two sensitivity problems, one for d_{13} and one for d_{24} . (If there was reason to believe there were trips between 1 and 4, and between 2 and 3, we would solve sensitivity problems for those OD pairs as well.) These are shown in Figure 8.13. The link performance functions are the same in both problems; the only difference is in the demand. The solutions to these sensitivity problems is shown in Table 8.2.

We now have all the information needed to calculate the gradient of the objective function, using equation (8.36). Substituting the values of x_{ij} , ξ_{ij} , the

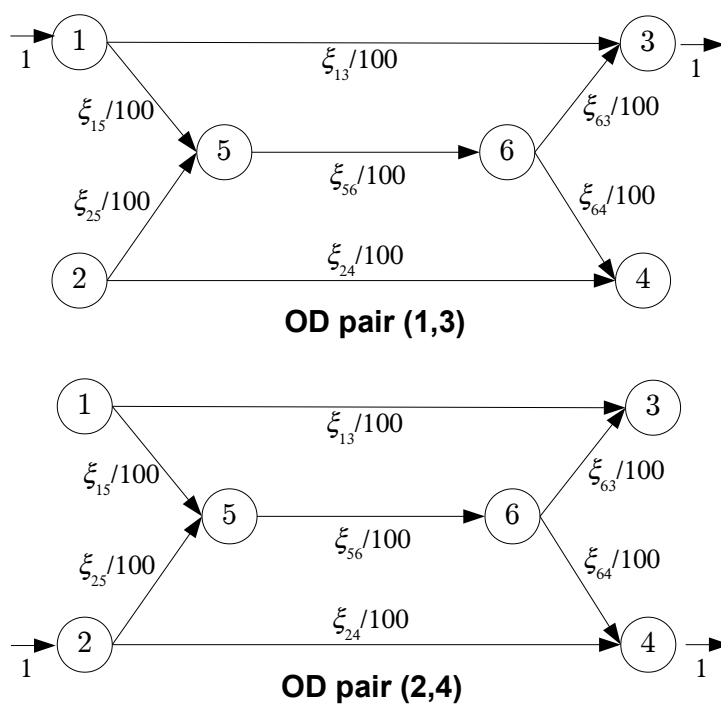


Figure 8.13: Two sensitivity problems for OD matrix estimation.

observed link flows, and the initial OD matrix, we calculate

$$\nabla_{\mathbf{d}} f = \begin{bmatrix} \frac{\partial f}{\partial d_{13}} \\ \frac{\partial f}{\partial d_{24}} \end{bmatrix} = \begin{bmatrix} 131 \\ 5.36 \end{bmatrix} \quad (8.38)$$

Taking a trial step in this direction with $\mu = 1$, the updating rule (8.37) gives the candidate OD matrix $d_{13} = 4869$, $d_{24} = 9995$. The resulting equilibrium link flows include $x_{13} = 4637$, $x_{25} = 1941$, and $x_{56} = 2150$, so the “fit” of the equilibrium link flows and traffic counts has improved from 11293 to 2293. The “fit” of the OD matrix has worsened from 0 to 17179, but with the weight $\Theta = 0.1$, the overall objective still decreases from 10164 to 3782. Therefore, we accept the step size $\mu = 1$, and return to step 3 to continue updating the OD matrix.

8.5 Exercises

1. [63] Given a nondegenerate equilibrium solution in a network with a single origin and continuous link performance functions, show that the equilibrium bushes remain unchanged in a small neighborhood of the current OD matrix.
2. [65] Given a nondegenerate equilibrium solution in a network with a single origin and differentiable link performance functions, show that the derivatives dx_{ij}/dd^{rs} exist at the current equilibrium solution.
3. [56] Extend the results in Exercises 1 and 2 to networks with multiple origins.
4. [84] Show that the entropy-maximizing path flow solution is a continuous function of the OD matrix.
5. [26] Verify that the optimality conditions for (8.11)–(8.13) include the conditions (8.6)–(8.10).
6. [33] In the modified Braess network of Figure 8.14, find the sensitivity of each link’s flow to the demand d^{14} .
7. [35] In the modified Braess network of Figure 8.14, find the sensitivity of each link’s flow to the link performance function parameter y^{23} . What value of this parameter minimizes the equilibrium travel time? Suggest what sort of real-world action would correspond to adjusting this parameter to this optimal value.
8. [44] In the network design problem for Figure 8.15, give the gradient of the objective function at the initial solution $\mathbf{y} = \mathbf{0}$. Assume that $\Theta = \frac{1}{20}$.

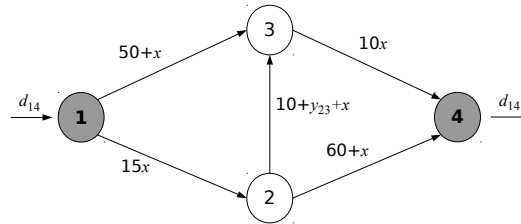


Figure 8.14: Network for use in Exercises 6 and 7.

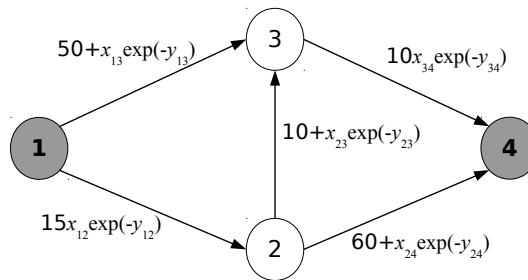


Figure 8.15: Network for use in Exercises 8 and 9.

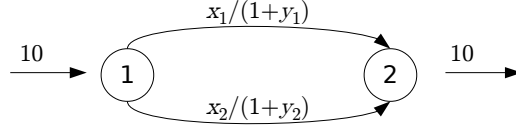


Figure 8.16: Network for use in Exercise 10.

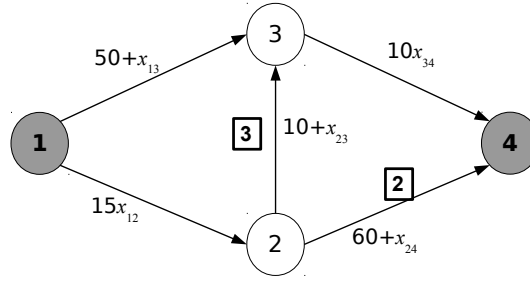


Figure 8.17: Network and observed link flows (in boxes) for Exercise 13.

9. [48] Continue Exercise 8 by performing three iterations of the algorithm given in the text. What is the resulting total system travel time and construction cost?
10. [25] Write out the network design optimization problem for the network in Figure 8.16, with $\Theta = 1$. Show that the feasible region for this problem is not convex by constructing a counterexample.
11. [79] Design a heuristic for the network design problem, based on simulated annealing as discussed in Section 4.5.1. Compare the performance of this heuristic with the algorithm given in the text for several networks.
12. [79] Design a heuristic for the network design problem, based on genetic algorithms as discussed in Section 4.5.2. Compare the performance of this heuristic with the algorithm given in the text for several networks.
13. [47] In the OD matrix estimation problem of Figure 8.17, give the gradient of the objective function at the initial solution $d_{14}^* = 6$, $d_{24}^* = 4$. What is the value of the objective function if $\Theta = \frac{1}{2}$?
14. [49] Continue Exercise 13 by performing three iterations of the algorithm given in the text. What is the resulting OD matrix and objective function value?
15. [59] Generate five additional OD matrices corresponding to the network

in Figure 8.17, with different values of Θ . Which of these OD matrices seems most reasonable to you, and why?

16. [79] Design a heuristic for the OD matrix estimation problem, based on simulated annealing as discussed in Section 4.5.1. Compare the performance of this heuristic with the algorithm given in the text for several networks.
17. [79] Design a heuristic for the OD matrix estimation problem, based on genetic algorithms as discussed in Section 4.5.2. Compare the performance of this heuristic with the algorithm given in the text for several networks.
18. [76] Perform the following “validation” exercise: create a small network with a given OD matrix, and find the equilibrium solution. Then, given the equilibrium link flows, try to compute your original OD matrix using the algorithm given in the text. Do you get your original OD matrix back?

Chapter 9

Extensions of Static Assignment

The basic traffic assignment problem (TAP) was defined in Chapter 6 as follows: we are given a network $G = (N, A)$, link performance functions $t_{ij}(x_{ij})$, and the demand values d^{rs} between each origin and destination. The objective is to find a feasible vector of path flows (or link flows) which satisfy the principle of user equilibrium, that is, that every path with positive flow has the least travel time among all paths connecting that origin and destination. We formulated this as a VI (find $\hat{\mathbf{h}} \in H$ such that $\mathbf{c}(\hat{\mathbf{h}}) \cdot (\hat{\mathbf{h}} - \mathbf{h}) \leq 0$ for all $\mathbf{h} \in H$) and as the solution to the following convex optimization problem:

$$\min_{\mathbf{x}, \mathbf{h}} \quad \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (9.1)$$

$$\text{s.t.} \quad x_{ij} = \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \quad \forall (i, j) \in A \quad (9.2)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (9.3)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (9.4)$$

This formulation remains the most commonly used version of traffic assignment in practice today. However, it is not difficult to see how some of the assumptions may not be reasonable. This chapter shows extensions of the basic TAP which relax these assumptions. This is the typical course of research: the first models developed make a number of simplifying assumptions, in order to capture the basic underlying behavior. Then, once the basic behavior is understood, researchers develop progressively more sophisticated and realistic models which relax these assumptions.

This chapter details three such extensions. Section 9.1 relaxes the assumption that the OD matrix is known and fixed, leading to an *elastic demand* formulation. Section 9.2 relaxes the assumption that the travel time on a link

depends only on the flow on that link (and not on any other link flows, even at intersections). Section 9.3 relaxes the assumption that travelers have accurate knowledge and perception of all travel times in a network, leading to the important class of *stochastic user equilibrium* models.

For simplicity, all of these variations are treated independently of each other. That is, the OD matrix is assumed known and fixed in all sections *except* Section 9.1, and so forth. This is done primarily to keep the focus on the relevant concept of each section, but also to guard the reader against the temptation to assume that a model which relaxes *all* of these assumptions simultaneously is necessarily better than one which does not. While realism is an important characteristic of a model, it is not the only relevant factor when choosing a mathematical model to describe an engineering problem. Other important concerns are computation speed, the existence of enough high-quality data to calibrate and validate the model, transparency, making sure the sensitivity of the model is appropriate to the level of error in input data, ease of explanation to decision makers, and so on. All of these factors should be taken into account when choosing a model, and you can actually do worse off by choosing a more “realistic” model when you don’t have adequate data for calibration — the result may even give the impression of “false precision” when in reality your conclusions cannot be justified.

9.1 Elastic Demand

The assumption that the OD matrix is known and fixed can strain credibility, particularly when considering long time horizons (20–30 years) or when projects are major enough to influence travel decisions at all levels, not just route choice. For instance, consider the “induced demand” phenomenon where major expansion of roadway capacity ends up increasing the amount of demand. This is partly due to changes in route choice (which the basic TAP accounts for), but is also due to changes in other kinds of travel choices, such as departure time, mode, destination, or trip frequency.

Therefore, it is desirable to develop a model which can relax the assumption of an exogenous OD matrix known *a priori*. This section describes how TAP can be extended to accommodate this relaxation. This is called the *elastic demand* formulation of TAP. The elastic demand model is at once more and less useful than basic TAP: more useful because it can provide a more accurate view of the impacts of transportation projects; less useful because it is harder to calibrate.

9.1.1 Demand functions

The new idea in the traffic assignment problem with elastic demand is the *demand function*, which relates the demand for travel between an origin and destination to the travel time between these zones. Specifically, let D^{rs} be a function relating the demand between r and s to the travel time κ^{rs} on the

shortest path between these zones.¹ Generally, D^{rs} is a nonincreasing function — as the travel time between r and s increases, the demand for travel between these nodes is lower. It will also be highly useful to assume that D^{rs} is invertible as well, which will require it to be strictly decreasing. The inverse demand function will give the travel time between r and s corresponding to a given demand level.

As an example, let $D^{rs}(\kappa^{rs}) = 200 \exp(-\kappa^{rs}/100)$. If the travel time between r and s is $\kappa^{rs} = 10$ by the shortest path, then the demand is $d^{rs} = 181$. If the travel time was 20 minutes between these zones, the demand will be 163.7, which is lower as fewer drivers choose to travel between r and s . The inverse demand function is $D_{rs}^{-1}(d^{rs}) = 100 \log(200/d^{rs})$, and can be used to calculate the κ^{rs} value corresponding to a given d^{rs} : when the demand is 181, the shortest path travel time is $\kappa^{rs} = D_{rs}^{-1}(181) = 10$, and so on.² Or, if $D^{rs} = 50 - \frac{1}{2}\kappa^{rs}$, then $D_{rs}^{-1} = 100 - 2d^{rs}$ and a travel time of 10 minutes corresponds to a demand of 45 vehicles as can be seen by substituting either number into the corresponding equation.

An attentive reader may have noticed a potential issue with the demand function $D^{rs} = 50 - \frac{1}{2}\kappa^{rs}$, namely that the demand would be negative if $\kappa^{rs} > 100$. In reality, the demand would simply equal zero if the travel time exceeded 100. We could patch this by redefining D^{rs} as $[50 - \frac{1}{2}\kappa^{rs}]^+$, but then D^{rs} is no longer invertible. Instead, we can allow D^{rs} to take negative values, but replace the relation $d^{rs} = D^{rs}(\kappa^{rs})$ with $d^{rs} = [D^{rs}(\kappa^{rs})]^+$. This allows us to have D^{rs} be strictly decreasing (and thus invertible), but still allow the travel demand to be zero when costs are sufficiently high. While this “trick” may seem a bit trivial (or at least not very useful), it will eventually allow us to formulate the elastic demand equilibrium problem as a variational inequality and convex program, as shown below.

The demand function can be used to define a *consumer surplus* CS representing the benefits of mobility in a region, defined by

$$CS = \sum_{(r,s) \in Z^2} \left(\int_0^{d^{rs}} D_{rs}^{-1}(y) dy - d^{rs} \kappa^{rs} \right) \quad (9.5)$$

(We are using y for the dummy variable of integration instead of d because $\int D^{-1}(d) dd$ is notationally awkward.) The interpretation of this formula is as follows. Each driver has a certain travel time threshold: if the travel time is greater than this threshold, the trip will not be made, and if the travel time is less than this threshold, the trip will be made. Different drivers have different

¹An alternative is to have D^{rs} be a function of the *average* travel time on the used paths between r and s , not the shortest. At equilibrium it doesn't matter because all used paths have the same travel time as the shortest, but in the process of finding an equilibrium the alternative definition of the demand function can be helpful. For our purposes, though, definition in terms of the shortest path time is more useful because it facilitates a link-based formulation.

²In the text, we typically indicate OD pairs with a superscript, as in d^{rs} , and link variables with a subscript, as in x_{ij} . In elastic demand, we will often need to refer to inverse functions, and writing $(D^{rs})^{-1}$ is clumsy. For this reason, OD pairs may also be denoted with a subscript, as in D_{rs}^{-1} . This is purely for notational convenience and carries no significance.

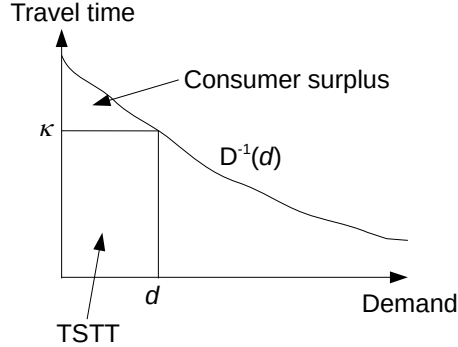


Figure 9.1: Relationship between inverse demand function, consumer surplus, and total system travel time.

thresholds, and the demand function represents the aggregation of these thresholds: when the travel time is κ , $D(\kappa)$ represents the number of travelers whose threshold is κ or higher. If my threshold is, say, 15 minutes and the travel time is 10 minutes, the difference (5 minutes) can be thought of as the “benefit” of travel to me: the trip is worth 15 minutes of my time, but I was able to travel for only 10. Adding this up for all travelers provides the total benefits of travel, which is what CS represents. Figure 9.1 shows the connection between this concept and equation (9.5): assume that drivers are numbered in decreasing order of their threshold values. Then $D^{-1}(1)$ gives the threshold value for the first driver, $D^{-1}(2)$ gives the threshold value for the second driver, and so forth. At equilibrium all drivers experience a travel time of κ , so the benefit to the first driver is $D^{-1}(1) - \kappa$, the benefit to the second driver is $D^{-1}(2) - \kappa$, and so forth. Adding over all drivers gives equation (9.5).

9.1.2 Gartner’s transformation

Before moving to variational inequality and optimization formulations of the elastic demand problem, we’ll take a short digression and show how the elastic demand problem can be cleverly transformed into a traditional equilibrium problem with fixed demand. This transformation works if each demand function D^{rs} is bounded above. Repeat the following for each OD pair (r, s) . Let \bar{d}_{rs} be such an upper bound for OD pair (r, s) . Create a new link (r, s) directly connecting origin r to destination s , and make its link performance function $t_{rs}(x_{rs}) = D_{rs}^{-1}(\bar{d}_{rs} - x_{rs})$ where x_{rs} is the flow on that new link. If D_{rs} is decreasing, then $D_{rs}^{-1}(\bar{d}_{rs} - x_{rs})$ is increasing in x_{rs} so this is a valid link performance function.

Now, solve a *fixed* demand problem where the demand from each origin to each destination is \bar{d}_{rs} . An equilibrium on this network corresponds to an elastic

demand equilibrium on the original network as follows: the flows on the links common to both networks represent flows on the actual traffic network; the flow on the new direct connection links represent drivers who choose not to travel due to excess congestion. Think of \bar{d}_{rs} as the total number of people who might possibly travel from r to s ; those that actually complete their trips travel on the original links and those who choose not to travel choose the direct connection link. At equilibrium, all used paths connecting r to s (including the direct connection link) have the same travel time κ_{rs} ; therefore, the flow on the direct connection link x_{rs} must be such that $D_{rs}^{-1}(\bar{d}_{rs} - x_{rs}) = \kappa_{rs}$, or equivalently $x_{rs} = \bar{d}_{rs} - D_{rs}(\kappa_{rs})$, which is exactly the number of drivers who choose not to travel when the equilibrium times are κ_{rs} . That is, the demand $d^{rs} = \bar{d}_{rs} - x_{rs}$.

The downside of this approach is that it requires creating a large number of new links. In a typical transportation network, the number of links is proportional to the number of nodes and of the same order of magnitude (so a network of 1,000 nodes may have 3–4,000 links). However, the number of OD pairs is roughly proportional to the *square* of the number of nodes, since every node could potentially be both an origin and a destination. So, a network with 1,000 nodes could have roughly 1,000,000 OD pairs. Implementing the Gartner transformation requires creating a new link for every one of these OD pairs, which would result in 99.9% of the network links being the artificial arcs for the transformation!

9.1.3 Variational inequality formulation

Because the elastic demand problem can be expressed as a version of the regular traffic assignment problem through the Gartner transformation, we immediately have a variational inequality formulation of the elastic demand equilibrium problem. Partition the vectors of link flows and travel times into regular and direct-connect (Gartner transformation) links, using \mathbf{x} to represent regular link flows, \mathbf{x}^{\rightarrow} flows on direct-connect links, and \mathbf{t} and \mathbf{t}^{\rightarrow} similarly. Then the variational inequality is

$$\begin{bmatrix} \mathbf{t}(\bar{\mathbf{x}}) \\ \mathbf{t}^{\rightarrow}(\bar{\mathbf{x}}^{\rightarrow}) \end{bmatrix} \cdot \begin{bmatrix} \bar{\mathbf{x}} - \mathbf{x} \\ \bar{\mathbf{x}}^{\rightarrow} - \mathbf{x}^{\rightarrow} \end{bmatrix} \leq 0 \quad (9.6)$$

Since $x_{rs} = \bar{d}_{rs} - d_{rs}$ and $t_{rs}(x_{rs}) = D_{rs}^{-1}(\bar{d}_{rs} - x_{rs})$, the variational inequality can be written in terms of the link flows and OD demands as

$$\begin{bmatrix} \mathbf{t}(\bar{\mathbf{x}}) \\ -\mathbf{D}^{-1}(\bar{\mathbf{d}}) \end{bmatrix} \cdot \begin{bmatrix} \bar{\mathbf{x}} - \mathbf{x} \\ \bar{\mathbf{d}} - \mathbf{d} \end{bmatrix} \leq 0 \quad (9.7)$$

which is the customary form.

9.1.4 Optimization formulation

The Gartner transformation can also lead directly to a convex programming formulation of the elastic demand problem, in a similar way as the variational inequality was derived in the previous section. However, it is also instructive to derive the convex programming formulation from first principles.

As discussed above, the demand is related to the demand function by $d^{rs} = [D^{rs}(\kappa^{rs})]^+$. Put another way, *the demand must always be at least as much as the demand function; further, if the demand is greater than zero then it must equal the demand function*. Thinking laterally, you might notice this is similar to the principle of user equilibrium: the travel time on any path must always be at least as large as the shortest path travel time; further, if the demand is positive then the path travel time must equal the shortest path travel time. When deriving the Beckmann function, we showed that the latter statements could be expressed by $c^\pi \geq \kappa^{rs}$ and $h^\pi(c^\pi - \kappa^{rs}) = 0$ (together with the nonnegativity condition $h^\pi \geq 0$). The same “trick” applies for the relationship between demand and the demand function: $d^{rs} \geq D^{rs}(\kappa^{rs})$, $d^{rs}(d^{rs} - D^{rs}(\kappa^{rs})) = 0$, and the nonnegativity condition $d^{rs} \geq 0$.

It will turn out to be easier to express the latter conditions in terms of the inverse demand functions D^{-1} , rather than the “forward” functions D , because the convex objective function we will derive will be based on the Beckmann function. The Beckmann function involves link performance functions (with units of time). Since the inverse demand functions also are measured in units of time, it will be easier to combine them with the link performance functions than the regular demand functions (which have units of vehicles). Expressed in terms of the inverse demand functions, the conditions above become $\kappa^{rs} \geq D_{rs}^{-1}(d^{rs})$, $d^{rs}(D_{rs}^{-1}(d^{rs}) - \kappa^{rs}) = 0$, and $d^{rs} \geq 0$.

So, this is the question before us. What optimization problem has the following as its optimality conditions?

$$c^\pi \geq \kappa^{rs} \quad \forall (r, s) \in Z^2, \pi \in \Pi^{rs} \quad (9.8)$$

$$h^\pi(c^\pi - \kappa^{rs}) = 0 \quad \forall (r, s) \in Z^2, \pi \in \Pi^{rs} \quad (9.9)$$

$$\kappa^{rs} \geq D_{rs}^{-1}(d^{rs}) \quad \forall (r, s) \in Z^2 \quad (9.10)$$

$$d^{rs}(D_{rs}^{-1}(d^{rs}) - \kappa^{rs}) = 0 \quad \forall (r, s) \in Z^2 \quad (9.11)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (9.12)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (9.13)$$

$$d^{rs} \geq 0 \quad \forall (r, s) \in Z^2 \quad (9.14)$$

The Beckmann formulation is a good place to start, since it already includes (9.8), (9.9), (9.12), and (9.13). So let's start by conjecturing that the Lagrangian takes the form

$$\mathcal{L}(\mathbf{h}, \mathbf{d}, \boldsymbol{\kappa}) = \sum_{(i,j) \in A} \int_0^{\sum_{\pi \in \Pi} \delta_{ij}^\pi h^\pi} t_{ij}(x) dx + \sum_{(r,s) \in Z^2} \kappa^{rs} \left(d^{rs} - \sum_{\pi \in \Pi^{rs}} h^\pi \right) + F(\mathbf{d}) \quad (9.15)$$

where $F(\mathbf{d})$ is some function involving the OD matrix. (Note also that \mathcal{L} is now a function of \mathbf{d} in addition to \mathbf{h} and $\boldsymbol{\kappa}$, since the demand is a decision variable.) You can check that the optimality conditions related to $\boldsymbol{\kappa}$ and \mathbf{h} are already included in the list of optimality conditions above. Assuming that there

is a nonnegativity constraint on the demand, (9.14) follows immediately as well. What's left is to show that the conditions $\partial\mathcal{L}/\partial d^{rs} \geq 0$ and $d^{rs}(\partial\mathcal{L}/\partial d^{rs}) = 0$ correspond to (9.10) and (9.11).

Calculating from (9.15), we have

$$\frac{\partial\mathcal{L}}{\partial d^{rs}} = \kappa^{rs} + \frac{\partial F}{\partial d^{rs}},$$

so if $\frac{\partial F}{\partial d^{rs}} = -D_{rs}^{-1}(d^{rs})$, we are done (both equations will be true). Integrating, $F(\mathbf{d}) = -\sum_{(r,s) \in Z^2} \int_0^{d^{rs}} D_{rs}^{-1}(\omega) d\omega$ gives us what we need. De-Lagrangianizing the “no vehicle left behind” constraint, we obtain the optimization problem associated with the elastic demand problem:

$$\min_{\mathbf{x}, \mathbf{h}, \mathbf{d}} \quad \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx - \sum_{(r,s) \in Z^2} \int_0^{d^{rs}} D_{rs}^{-1}(\omega) d\omega \quad (9.16)$$

$$\text{s.t.} \quad x_{ij} = \sum_{\pi \in \Pi} h^\pi \delta_{ij}^\pi \quad \forall (i,j) \in A \quad (9.17)$$

$$\sum_{\pi \in \Pi^{rs}} h^\pi = d^{rs} \quad \forall (r,s) \in Z^2 \quad (9.18)$$

$$h^\pi \geq 0 \quad \forall \pi \in \Pi \quad (9.19)$$

$$d^{rs} \geq 0 \quad \forall (r,s) \in Z^2 \quad (9.20)$$

9.1.5 Solution method

This section shows how the Frank-Wolfe algorithm can be used to solve the optimization problem (9.16)–(9.20). This is certainly not the only choice, and it is worthwhile for you to think about how other algorithms from Chapter 7 could also be used instead of Frank-Wolfe. The implementation of this algorithm is quite similar to how Frank-Wolfe works for the basic traffic assignment problem, with three changes. First, since the OD matrix is a decision variable along with the link flows, we must keep track of both \mathbf{d} as well as \mathbf{x} ; therefore, in addition to the target link flows \mathbf{x}^* we will have a target OD matrix \mathbf{d}^* , and in addition when we update the link flows we must update the OD matrix as well. Each of these pairs of flows and OD matrices should be consistent with each other, in that the link flows \mathbf{x} must be a feasible network loading when the demand is \mathbf{d} (both before and after updating), and similarly \mathbf{x}^* must correspond to \mathbf{d}^* . Luckily, this will not be difficult.

The second change to the Frank-Wolfe algorithm is how the restricted variational inequality is solved. Instead of solving $\mathbf{t}(\bar{\mathbf{x}}') \cdot (\mathbf{x}^* - \mathbf{x}) \leq 0$ for $\bar{\mathbf{x}}' = \lambda \mathbf{x}^* + (1 - \lambda)\mathbf{x}$, $\lambda \in [0, 1]$, we must solve the variational inequality (9.7) in $\bar{\mathbf{x}}'$ and $\bar{\mathbf{d}}'$. As before, the usual solution involves an “interior” $\lambda \in (0, 1)$, in which

case we must solve the equation

$$\sum_{(i,j) \in A} t_{ij}(\lambda x_{ij}^* + (1 - \lambda)x_{ij})(x_{ij}^* - x_{ij}) - \sum_{(r,s) \in Z^2} D_{rs}^{-1}(\lambda d_{rs}^* + (1 - \lambda)d_{rs})(d_{rs}^* - d_{rs}) = 0 \quad (9.21)$$

in λ . This is simply the variational inequality (9.7) written out in terms of its components, substituting $\lambda \mathbf{x}^* + (1 - \lambda)\mathbf{x}$ for $\bar{\mathbf{x}}'$ and $\lambda \mathbf{d}^* + (1 - \lambda)\mathbf{d}$ for $\bar{\mathbf{d}}'$.

Third, the stopping criterion (relative gap or average excess cost) needs to be augmented with a measure of how well the OD matrix matches the values from the demand functions. A simple measure is the *total misplaced flow* defined as $TMF = \sum_{(r,s) \in Z^2} |d^{rs} - [D^{rs}(\kappa^{rs})]^+|$. The total misplaced flow is always nonnegative, and is zero only if all of the entries in the OD matrix are equal to the values given by the demand function (or zero if the demand function is negative). We should keep track of both total misplaced flow and one of the equilibrium convergence measures (relative gap or average excess cost), and only terminate the algorithm when both of these are sufficiently small.

Implementing these changes, the Frank-Wolfe algorithm for elastic demand is as follows:

1. Choose some initial OD matrix \mathbf{d} and initial link flows \mathbf{x} corresponding to that OD matrix.
2. Find the shortest path between each origin and destination, and calculate convergence measures (total misplaced flow, and either relative gap or average excess cost). If both are sufficiently small, stop.
3. Improve the solution:
 - (a) Calculate a target OD matrix \mathbf{d}^* using the demand functions: $d_{rs}^* = [D^{rs}(\kappa^{rs})]^+$ for all OD pairs (r, s) .
 - (b) Using the target matrix \mathbf{d}^* , find the link flows if everybody were traveling on the shortest paths found in step 2, store these in \mathbf{x}^* .
 - (c) Solve the restricted variational inequality by finding λ such that (9.21) is true.
 - (d) Update the OD matrix and link flows: replace \mathbf{x} with $\lambda \mathbf{x}^* + (1 - \lambda)\mathbf{x}$ and replace \mathbf{d} with $\lambda \mathbf{d}^* + (1 - \lambda)\mathbf{d}$.
4. Return to step 1.

This algorithm can also be linked to the convex programming formulation described above. Given a current solution (\mathbf{x}, \mathbf{d}) , it can be shown that the derivative of the Beckmann function in the direction towards $(\mathbf{x}^*, \mathbf{d}^*)$ is nonpositive (and strictly negative if the current solution does not solve the elastic demand problem), and that the solution of the restricted variational inequality (9.21) minimizes the objective function along the line joining (\mathbf{x}, \mathbf{d}) to $(\mathbf{x}^*, \mathbf{d}^*)$. The algebra is a bit tedious and is left as an exercise at the end of the chapter.

9.1.6 Example

Here we solve the small example of Figure 9.2 with the Frank-Wolfe algorithm, using the average excess cost to measure how close we are to equilibrium. The demand function is $D(\kappa) = 50 - \kappa$, so its inverse function is $D^{-1}(d) = 50 - d$.

Initialization. Arbitrarily set $d = 50$, then arbitrarily load all 50 vehicles onto the two links; say $\mathbf{x} = [50 \ 0]$.

Iteration 1. The link travel times are now $\mathbf{t} = [60 \ 20]$, so the shortest path travel time $\kappa = 20$ and the demand function indicates that the demand should be $D = 50 - 20 = 30$. The average excess cost is $(60 \times 50 - 20 \times 50)/50 = 40$, and the total misplaced flow is $|50 - 30| = 20$. The target demand is what the demand function indicates $d^* = 30$, and this flow should all be loaded on the bottom path, so $\mathbf{x}^* = [0 \ 30]$. Solving the equation

$$(10 + 50(1 - \lambda))(-50) + (20 + 30\lambda)(30) - (50 - (30\lambda + 50(1 - \lambda)))(-20) = 0$$

we obtain $\lambda = 12/19 \approx 0.632$ so the new demand and flows are $d = 37.36$ and $\mathbf{x} = [18.42 \ 18.95]$.

Iteration 2. The link travel times are now $\mathbf{t} = [28.42 \ 38.95]$, so $\kappa = 28.42$, $D = 21.58$, $AEC = 5.19$, and $TMF = 15.78$. Both convergence measures have decreased from the first iteration, particularly the average excess cost. Thus, $d^* = 21.58$, $\mathbf{x}^* = [21.58 \ 0]$, and we solve

$$(10 + 21.58\lambda + 18.42(1 - \lambda))(21.58 - 18.42) + (20 + 18.95(1 - \lambda))(0 - 18.95) - (50 - (21.58\lambda + 37.36(1 - \lambda)))(21.58 - 37.36) = 0$$

so $\lambda = 0.726$ and the new demand and flows are $d = 25.9$ and $\mathbf{x} = [20.71 \ 5.19]$

Iteration 3 . The link travel times are now $\mathbf{t} = [30.71 \ 25.19]$, so $\kappa = 25.19$, $D = 24.09$, $AEC = 4.42$, and $TMF = 1.80$. Assuming that these are small enough to terminate, we are done.

9.2 Link Interactions

This section develops another extension of TAP, in which we relax the assumption that the travel time on a link depends only on the flow on that link. There are a few reasons why this kind of extension may be useful:

Junction interactions: At an intersection, the delay on a particular approach often depends on flows from competing approaches. As an example, consider a freeway onramp which has to yield to mainline traffic at a merge.

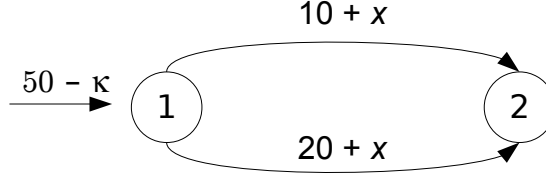


Figure 9.2: Example network for elastic demand equilibrium.

Because merging traffic must find an acceptably large gap in the main lanes, the travel time on the onramp depends on the flow on the main lanes as well as the flow on the onramp. Similar arguments hold at arterial junctions controlled by two-way or four-way stops, at signalized intersections with permissive phases (e.g., left-turning traffic yielding to gaps in oncoming flow), or at actuated intersections where the green times are determined in real-time based on available flow. The basic TAP cannot model the link interactions characterizing these types of links.

Overtaking traffic: On rural highways, overtaking slow-moving vehicles often requires finding a (fairly large) gap in oncoming flow. If the oncoming flow is small, the effect of slow-moving vehicles on average travel time is negligible. However, as the oncoming flow becomes larger and larger, the ability to overtake is diminished and traffic speeds will tend to be determined by the slowest-moving vehicle on the highway. Since traffic moving in different directions on the same highway is modeled with different links, a link interaction model is needed to capture this effect.

Multiclass flow: Consider a network model where there are two types of vehicles (say, passenger cars and semi trucks, or passenger cars and buses). Presumably these vehicles may choose routes differently or even have a different roadway network available to them (heavy vehicles are prohibited from some streets, and buses must drive along a fixed route). This type of situation can be modeled by creating a “two-layer” network, with the two layers representing the links available to each class. However, where these links represent the same physical roadway, the link performance functions should be connected to each other (truck volume influences passenger car speed and vice versa) even if they are not identical (truck speed need not be the same as passenger car speed). Link interaction models therefore allow us to model multiclass flow as well.

However, there are a few twists to the story, some of which are explored below. Section 9.2.1 presents a mathematical formulation of the link interactions

model, but shows that a convex programming formulation is not possible except in some rather unlikely cases. Section 9.2.2 explores the properties of the link interactions model, in particular addressing the issue of uniqueness — even though link flow solutions to TAP are unique under relatively mild assumptions, this is not generally true when there are link interactions. Section 9.2.3 gives us two solution methods for the link interactions model, the diagonalization method and simplicial decomposition. Diagonalization is easier to implement, but simplicial decomposition is generally more powerful.

9.2.1 Formulation

In the basic TAP, the link performance function for link (i, j) was a function of x_{ij} alone, that is, we could write $t_{ij}(x_{ij})$. Now, t_{ij} may depend on the flow on multiple links. For full generality, our notation will allow t_{ij} to depend on the flows on *any or all* other links in the network: the travel time is given by the function $t_{ij}(x_1, x_2, \dots, x_{ij}, \dots, x_m)$ or, more compactly, $t_{ij}(\mathbf{x})$ using vector notation. We assume these are given to us. Everything else is the same as in vanilla TAP: origin-destination demand is fixed, and we seek an equilibrium solution where all used paths have equal and minimal travel time.

Now, how to formulate the equilibrium principle? It's not hard to see that the variational inequality for TAP works equally well here:

$$\mathbf{c}(\bar{\mathbf{h}}) \cdot (\bar{\mathbf{h}} - \mathbf{h}) \leq 0 \quad \forall \mathbf{h} \in H \quad (9.22)$$

where the only difference is that the link performance functions used to calculate path travel times C are now of the form $t_{ij}(\mathbf{x})$ rather than $t_{ij}(x_{ij})$. But this is of no consequence. Path flows $\bar{\mathbf{h}}$ solve the variational inequality if and only if

$$\mathbf{c}(\bar{\mathbf{h}}) \cdot \bar{\mathbf{h}} \leq \mathbf{c}(\bar{\mathbf{h}}) \cdot \mathbf{h} \quad (9.23)$$

for any other feasible path flows \mathbf{h} . That is, if the travel times were fixed at their current values, then it is impossible to reduce the total system travel time by changing any drivers' route choices. This is only possible if all used paths have equal and minimal travel time. Similarly, the link-flow variational inequality

$$\mathbf{t}(\bar{\mathbf{x}}) \cdot (\bar{\mathbf{x}} - \mathbf{x}) \leq 0, \quad (9.24)$$

where \mathbf{x} is any feasible link flow, also represents the equilibrium problem with link interactions.

The ease of translating the variational inequality formulation for the case of link interactions may give us hope that a convex programming formulation exists as well. The feasible region is the same, all we need is to find an appropriate objective function. Unfortunately, this turns out to be a dead end. For example, the obvious approach is to amend the Beckmann function in some way, for instance, changing

$$\sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad \text{to} \quad \sum_{(i,j) \in A} \int_0^{\mathbf{x}} t_{ij}(\mathbf{y}) d\mathbf{y} \quad (9.25)$$

where the simple integral in the Beckmann function is replaced with a line integral between the origin $\mathbf{0}$ and the current flows \mathbf{x} . Unfortunately, this line integral is in general not well-defined, since its value depends on the path taken between the origin and \mathbf{x} .

The one exception is if the vector of travel times $\mathbf{t}(\mathbf{x})$ is a gradient map (that is, it is a conservative vector field). In this case, the fundamental theorem of line integrals implies that the value of this integral is *independent* of the path taken between $\mathbf{0}$ and \mathbf{x} . For $\mathbf{t}(\mathbf{x})$ to be a gradient map, its Jacobian must be symmetric. That is, for every pair of links (i, j) and (k, ℓ) , we need the following condition to be true:

$$\frac{\partial t_{ij}}{\partial x_{k\ell}} = \frac{\partial t_{k\ell}}{\partial x_{ij}} \quad (9.26)$$

That is, *regardless of the current flow vector \mathbf{x}* , the marginal impact of another vehicle added to link (i, j) on the travel time of (k, ℓ) must equal the marginal impact of another vehicle added to link (k, ℓ) on the travel time of (i, j) . This condition is very strong. Comparing with the motivating examples used to justify studying link interactions, the symmetry condition is not usually satisfied: the impact of an additional unit of flow on the mainline on the onramp travel time is much greater than the impact of an additional unit of onramp flow on mainline travel time. The impact of semi truck flow on passenger car travel time is probably greater than the impact of passenger car flow on truck travel time at the margin. Symmetry may perhaps hold in the case of overtaking on a rural highway, but even then it is far from clear. So, when modeling link interactions we cannot hope for condition (9.26) to hold. If it does so, consider it a happy accident: the function (9.25) is then an appropriate convex optimization problem.

9.2.2 Properties

This section explores the properties of the link interaction equilibrium problem defined by the variational inequality (9.22). The first question concerns existence of an equilibrium. Because (9.22) is essentially the same variational inequality derived for TAP, the arguments used to derive existence of an equilibrium (based on Brouwer's theorem) carry over directly and we have the same result:

Proposition 9.1. *If each link performance function $t_{ij}(\mathbf{x})$ is continuous in the vector of link flows \mathbf{x} , then at least one solution exists satisfying the principle of user equilibrium.*

However, uniqueness turns out to be trickier. Consider a network of two parallel links where the demand is 6 vehicles, and the link performance functions are $t_1 = x_1 + 2x_2$ and $t_2 = 2x_1 + x_2$. Setting the two links' travel times equal to each other and using $x_1 + x_2 = 6$, it is easy to see that one equilibrium is the solution $x_1 = x_2 = 3$ when the travel time on both links is 9. However, this is not the only equilibrium solution: if all of the drivers were to choose link 1,

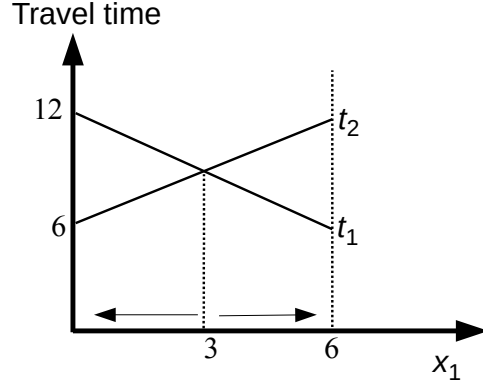


Figure 9.3: Change in path travel times as x_1 varies, “artificial” two-link network.

then $x_1 = 6$, $x_2 = 0$, $t_1 = 6$, and $t_2 = 12$. The top link is the only used path, but it has the least travel time so this solution also satisfies the principle of user equilibrium. Likewise, if $x_1 = 0$ and $x_2 = 6$, then $t_1 = 12$ and $t_2 = 6$ and again the only used path has the least travel time. Therefore, this network has three equilibrium solutions; compare with Figure 9.3.

To make this situation less artificial, we can change the link performance functions to represent a more realistic scenario. Assume that the rate of demand is 1800 vehicles per hour, and that link 1 has a constant travel time of 300 seconds independent of the flow on either link. Link 2 is shorter with a free-flow time of 120 seconds, but must yield to link 1 using gap acceptance principles. In traffic operations, gap acceptance is often modeled with two parameters: the *critical gap* t_c , and the *follow-up gap* t_f . The critical gap is the smallest headway required in the main stream for a vehicle to enter. Given that the gap is large enough for one vehicle to enter the stream, the follow-up gap is the incremental amount of time needed for each additional vehicle to enter. For this example, let t_c be 4 seconds and t_f be 2 seconds. Then, assuming that flows on both links 1 and 2 can be modeled as Poisson arrivals, the travel time on link 2 can be derived as

$$t_2(x_1, x_2) = \frac{1}{u} + \frac{\Lambda}{4} \left[\frac{x_2}{u} - 1 + \sqrt{\left(\frac{x_2}{u} - 1 \right)^2 + \frac{8x_2}{u^2\Lambda}} \right] \quad (9.27)$$

where Λ is the length of the analysis period and u is the capacity of link 2 defined by

$$u = \frac{x_2 \exp(-x_1 t_c)}{1 - \exp(-x_1 t_f)} \quad (9.28)$$

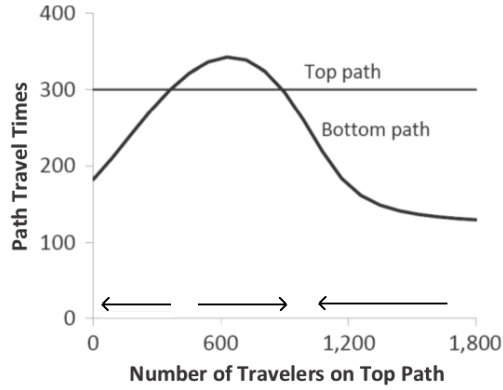


Figure 9.4: Change in path travel times as x_1 varies, “realistic” two-link network.

Figure 9.4 shows the travel times on the two paths as x_1 varies. Again, there are three equilibria: (1) $x_1 = 0$, $x_2 = 1800$, where $t_1 = 300$ and $t_2 = 182$; (2) $x_1 = 362$, $x_2 = 1438$, where $t_1 = t_2 = 300$; and (3) $x_1 = 892$, $x_2 = 908$, where $t_1 = t_2 = 300$.

So, even in realistic examples we cannot expect equilibrium to be unique when there are link interactions. The practical significance is that it raises doubt about which equilibrium solution should be used for project evaluation or ranking. For instance, consider a candidate project which would improve the free-flow time on link 1 from 300 to a smaller value; this would correspond to lowering the horizontal line in Figure 9.4. If we are at one of the equilibria where the travel times are equal, such a change will indeed reduce the travel times experienced by drivers. However, if we are at the equilibrium where the top path is unused, such a change will have no impact whatsoever.

While a complete study of the methods used to distinguish among multiple equilibria is beyond the scope of this section, a simple stability criterion is explained here: an equilibrium solution is *stable* if small perturbations to the solution would incentivize drivers to move back towards that initial equilibrium — that is, if we reassign a few drivers to different paths, the path travel times will change in such a way that those drivers would want to move back to their original paths. By contrast, an *unstable* equilibrium does not have this property: if a few drivers are assigned to different paths, the path travel times will change in such a way that even more drivers would want to switch paths, and so on until another equilibrium is found.

In the simple two-link network we’ve been looking at, stability can be identified using graphs such as those in Figures 9.3 and 9.4. The arrows on the bottom axis indicate the direction of the path switching which would occur for a given value of x_1 . When the arrow is pointing to the left, $t_1 > t_2$ so travelers want to switch away from path 1 to path 2, resulting in a decrease in x_1 (a move

further to the left on the graph). When $t_2 > t_1$, travelers want to switch away from path 2 to path 1, resulting in an increase in x_1 , indicated by an arrow pointing to the right. At the equilibrium solutions, there is no pressure to move in any feasible direction. So, for the first example, the only stable equilibria are the “extreme” solutions with all travelers on either the top or bottom link. The equilibrium with both paths used is unstable in the sense that any shift from one path to another amplifies the difference in travel times and encourages even more travelers to shift in that direction. In the second example, the first and third equilibria are stable, but the second is unstable.

Based on these two examples, an intuitive explanation for the presence of stability with link interactions can be provided. For the regular traffic assignment problem with increasing link performance functions, shifting flow away from a path π and onto another path π' always decreases the travel time on π and increases the travel time on π' . Therefore, if the paths have different travel times, flow will shift in a way that always tends to equalize the travel times on the two paths. Even where there are link interactions, the same will hold true *if the travel time on a path is predominantly determined by the flow on that path*. However, when the link interactions are very strong, the travel time on a path may depend more strongly by the flow on a different path. In the first example, notice that each link's travel time is influenced more by the *other* link's flow than its own. In the second example, for certain ranges of flow the travel time on the merge path is influenced more by the flow on the priority path. In such cases, there is no guarantee that moving flow from a higher-cost path to a lower-cost path will tend to equalize their travel times. In the first example, we have an extreme case where moving flow to a path *decreases* its travel time while *increasing* the travel time of the path the flow moved away from!

To make this idea more precise, the following section introduces the mathematical concept of strict monotonicity.

Strict Monotonicity

Let $\mathbf{f}(\mathbf{x})$ be a vector-valued function whose domain and range are vectors of the same dimension. For instance, $\mathbf{t}(\mathbf{x})$ maps the vector of link flows to the vector of link travel times; the dimension of both of these is the number of links in the network. We say that \mathbf{f} is *strictly monotone* if for any two distinct vectors \mathbf{x} and \mathbf{y} in its domain, the dot product of $\mathbf{f}(\mathbf{x}) - \mathbf{f}(\mathbf{y})$ and $\mathbf{x} - \mathbf{y}$ is strictly positive.

For example, let $\mathbf{f}(\mathbf{x})$ be defined by $f_1(x_1, x_2) = 2x_1$ and $f_2(x_1, x_2) = 2x_2$. Then for any distinct vectors \mathbf{x} and \mathbf{y} , we have

$$\begin{aligned} (\mathbf{f}(\mathbf{x}) - \mathbf{f}(\mathbf{y})) \cdot (\mathbf{x} - \mathbf{y}) &= ([2x_1 \quad 2x_2] - [2y_1 \quad 2y_2]) \cdot ([x_1 \quad x_2] - [y_1 \quad y_2]) \\ &= 2((x_1 - y_1)^2 + (x_2 - y_2)^2) \end{aligned}$$

Since $\mathbf{x} \neq \mathbf{y}$, the right-hand side is always greater than zero, so f is monotone. As another example, let the function $\mathbf{g}(\mathbf{x})$ be defined by $g_1(x_1, x_2) = 2x_2$

and $g_2(x_1, x_2) = 2x_1$. If we choose $\mathbf{x} = [0 \ 1]$ and $\mathbf{y} = [1 \ 0]$, then

$$(\mathbf{g}(\mathbf{x}) - \mathbf{g}(\mathbf{y})) \cdot (\mathbf{x} - \mathbf{y}) = ([2 \ -2]) \cdot ([-1 \ 1]) = -4$$

so \mathbf{g} is not strictly monotone. Note that proving strict monotonicity requires a general argument valid for *any* distinct vectors \mathbf{x} and \mathbf{y} ; showing that a function is not strictly monotone only requires a single counterexample.

Warning! It is very common for students to think that the link performance functions are strictly monotone if they are strictly increasing functions of the flow on each link. This is not true: in the first example in this section, all link performance functions are increasing in each flow variable but if we compare $\mathbf{x} = [0 \ 6]$ and $\mathbf{y} = [5 \ 1]$, we have

$$(\mathbf{t}(\mathbf{x}) - \mathbf{t}(\mathbf{y})) \cdot (\mathbf{x} - \mathbf{y}) = ([12 \ 6] - [7 \ 11]) \cdot ([0 \ 6] - [5 \ 1]) = -50$$

so these link performance functions are not strictly monotone. Roughly speaking, strict monotonicity requires the diagonal terms of the Jacobian of \mathbf{t} to be large compared to the off-diagonal terms. The precise version of this “roughly speaking” fact is the following:

Proposition 9.2. *If \mathbf{f} is a continuously differentiable function whose domain is convex, then \mathbf{f} is strictly monotone if and only if its Jacobian is positive definite at all points in the domain.*

With this definition of monotonicity in hand, we can provide the uniqueness result we’ve been searching for:

Proposition 9.3. *Consider an instance of the traffic assignment problem with link interactions. If the link performance functions $\mathbf{t}(\mathbf{x})$ are continuous and strictly monotone, then there is exactly one user equilibrium solution.*

Proof. Since $\mathbf{t}(\mathbf{x})$ is continuous, we are guaranteed existence of at least one equilibrium solution from Brouwer’s theorem; let $\hat{\mathbf{x}}$ be such an equilibrium and let $\tilde{\mathbf{x}}$ be any other feasible link flow solution. We need to show that $\tilde{\mathbf{x}}$ cannot be an equilibrium. Arguing by contradiction, assume that $\tilde{\mathbf{x}}$ is in fact an equilibrium. Then it would solve the variational inequality (9.24), so

$$\mathbf{t}(\tilde{\mathbf{x}}) \cdot (\tilde{\mathbf{x}} - \hat{\mathbf{x}}) \leq 0$$

Adding a clever form of zero to the left hand side, this would imply

$$(\mathbf{t}(\tilde{\mathbf{x}}) - \mathbf{t}(\hat{\mathbf{x}})) \cdot (\tilde{\mathbf{x}} - \hat{\mathbf{x}}) + \mathbf{t}(\hat{\mathbf{x}}) \cdot (\tilde{\mathbf{x}} - \hat{\mathbf{x}}) \leq 0 \quad (9.29)$$

But since the link performance functions are strictly monotone, the first term on the left-hand side is strictly positive. Furthermore, since $\hat{\mathbf{x}}$ is an equilibrium the variational inequality (9.24) is true, so $\mathbf{t}(\hat{\mathbf{x}}) \cdot (\hat{\mathbf{x}} - \tilde{\mathbf{x}}) \leq 0$, which implies that the second term on the right-hand side is nonnegative. Therefore, the left-hand side of (9.29) is strictly positive, which is a contradiction. Therefore $\tilde{\mathbf{x}}$ cannot satisfy the principle of user equilibrium. \square

9.2.3 Algorithms

This section presents two algorithms for the traffic assignment problem with link interactions. If the link performance functions are strictly monotone, it can be shown that both of these algorithms converge to the unique equilibrium solution. Otherwise, it is possible that these algorithms may not converge, although they will typically do so if they start sufficiently close to an equilibrium. In any case, these algorithms may be acceptable heuristics even when strict monotonicity does not hold.

Diagonalization

The diagonalization method is a variation of Frank-Wolfe, which differs only in how the step size λ is found. Recall that the Frank-Wolfe step size is found by solving the equation

$$\sum_{(i,j) \in A} t_{ij}(\lambda x_{ij}^* + (1 - \lambda)x_{ij})(x_{ij}^* - x_{ij}) = 0 \quad (9.30)$$

since this minimizes the Beckmann function along the line segment connecting \mathbf{x} to \mathbf{x}^* . Since there is no corresponding objective function when there are asymmetric link interactions, it is not clear that a similar approach will necessarily work. (And in any case, t_{ij} is no longer a function of x_{ij} alone, so the formula as stated will not work.)

To make this formula logical, construct a temporary link performance function $\tilde{t}_{ij}(x_{ij})$ which only depends on its own flow. This is done by assuming that the flow on all other links is constant: $\tilde{t}_{ij}(x_{ij}) = t_{ij}(x_1, \dots, x_{ij}, \dots, x_m)$. For example, if $t_1(x_1, x_2, x_3) = x_1 + x_2^2 + x_3^3$ and the current solution is $x_1 = 1$, $x_2 = 2$, and $x_3 = 3$, then $\tilde{t}_1(x_1) = 31 + x_1$, since this is what we would get if x_2 and x_3 were set to constants at their current values of 2 and 3, respectively.

The step size λ is then found by adapting the Frank-Wolfe formula, using $\tilde{t}_{ij}(x_{ij})$ in place of $t_{ij}(\mathbf{x})$. That is, in the diagonalization method λ solves

$$\sum_{(i,j) \in A} \tilde{t}_{ij}(\lambda x_{ij}^* + (1 - \lambda)x_{ij})(x_{ij}^* - x_{ij}) = 0 \quad (9.31)$$

At each iteration, new \tilde{t} functions are calculated based on the current solution. The complete algorithm is as follows:

1. Find the shortest path between each origin and destination, and calculate the relative gap (unless it is the first iteration). If the relative gap is sufficiently small, stop.
2. Shift travelers onto shortest paths:
 - (a) Find the link flows if everybody were traveling on the shortest paths found in step 1, store these in \mathbf{x}^* .

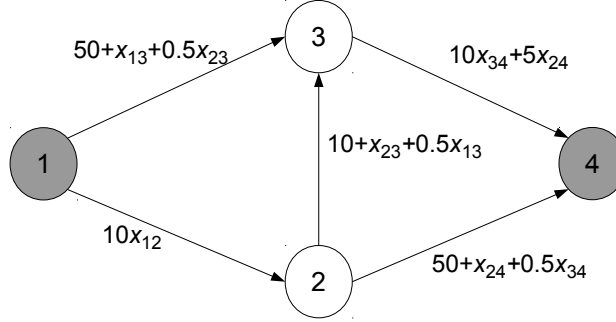


Figure 9.5: Braess network with link interactions.

- (b) If this is the first iteration, set $\mathbf{x} \leftarrow \mathbf{x}^*$ and move to step 3. Otherwise, continue with step c.
 - (c) Using the current solution \mathbf{x} , form the diagonalized link performance functions $\tilde{t}_{ij}(x_{ij})$ for each link.
 - (d) Find λ which solves equation (9.31).
 - (e) Update $\mathbf{x} \leftarrow \lambda \mathbf{x}^* + (1 - \lambda)\mathbf{x}$.
3. Calculate the new link travel times and the relative gap. Increase the iteration counter i by one and return to step 1.

As an example, consider the modified Braess network shown in Figure 9.5. At each merge node, the travel time on each of the incoming links depends on the flow on *both* links which merge together. The link flow vectors are indexed $\mathbf{x} = [x_{12} \ x_{13} \ x_{23} \ x_{24} \ x_{34}]$.

Iteration 1. In the first iteration, load all travelers onto shortest paths, so $\mathbf{x} = \mathbf{x}^* = [6 \ 0 \ 6 \ 0 \ 6]$, $\mathbf{t} = [60 \ 53 \ 16 \ 53 \ 60]$ and the average excess cost is 23.

Iteration 2. The all-or-nothing loading on shortest paths given \mathbf{t} is $\mathbf{x}^* = [0 \ 6 \ 0 \ 0 \ 6]$. The diagonalized link performance functions are obtained by assuming the flows on all other links are constant at \mathbf{x} : $\hat{t}_{12} = 10x_{12}$, $\hat{t}_{13} = 53 + x_{13}$, $\hat{t}_{23} = 10 + x_{23}$, $\hat{t}_{24} = 53 + x_{24}$, and $\hat{t}_{34} = 10x_{34}$. So we solve the equation (9.31) for λ :

$$(53 + 6\lambda)6 + 10(6(1 - \lambda))(-6) + (10 + 6(1 - \lambda))(-6) = 0$$

omitting terms where $x_{ij} = x_{ij}^*$ because they are zero in (9.31). The solution is $\lambda = 23/72$, so we update $\mathbf{x} = [4\frac{1}{12} \ 1\frac{11}{12} \ 4\frac{1}{12} \ 0 \ 6]$ and $\mathbf{t} = [40\frac{5}{6} \ 53\frac{23}{24} \ 15\frac{1}{24} \ 53 \ 60]$ (using the regular cost functions, not the diagonalized ones.) The average excess cost is now 21.43.

Iteration 3. The all-or-nothing assignment is $\mathbf{x}^* = [6 \ 0 \ 0 \ 6 \ 0]$ and the diagonalized link performance functions are $\hat{t}_{12} = 10x_{12}$, $\hat{t}_{13} = 52\frac{1}{24} + x_{13}$, $\hat{t}_{23} = 10\frac{23}{24} + x_{23}$, $\hat{t}_{24} = 53 + x_{24}$, and $\hat{t}_{34} = 10x_{34}$. The solution to (9.31) is $\lambda = 0.284$, which gives $\mathbf{x} = [4.63 \ 1.37 \ 2.92 \ 1.70 \ 4.30]$, $\mathbf{t} = [46.3 \ 52.8 \ 13.6 \ 53.9 \ 51.5]$, so the average excess cost is 6.44.

and so on until convergence is reached.

Simplicial decomposition

An alternative to diagonalization is the simplicial decomposition algorithm. This algorithm is introduced at this point (rather than in Chapter 7) for several reasons. First, historically it was the first provably convergent algorithm for the equilibrium problem with link interactions. Second, although it is an improvement on Frank-Wolfe, for the basic TAP it is outperformed by the path-based and bush-based algorithms presented in that chapter. However, like those algorithms, it overcomes the “zig-zagging” difficulty that Frank-Wolfe runs into (cf. Figure 7.4).

The price of this additional flexibility is that more computer memory is needed. Frank-Wolfe and MSA are exceptionally economical in that they only require two vectors to be stored: the current link flows \mathbf{x} and the target link flows \mathbf{x}^* . In simplicial decomposition, we will “remember” all of the target link flows found in earlier iterations, and exploit this longer-term memory by allowing “combination” moves towards several of these previous targets simultaneously. In the algorithm, the set \mathcal{X} is used to store all target link flows found thus far.

A second notion in simplicial decomposition is that of a “restricted equilibrium.” Given a set $\mathcal{X} = \{\mathbf{x}_1^*, \mathbf{x}_2^*, \dots, \mathbf{x}_k^*\}$ and a current link flow solution \mathbf{x} , we say that \mathbf{x} is a restricted equilibrium if it solves the variational inequality

$$\mathbf{t}(\mathbf{x}) \cdot (\mathbf{x} - \mathbf{x}') \leq 0 \quad \forall \mathbf{x}' \in X(\mathbf{x}, \mathcal{X}) \quad (9.32)$$

where $X(\mathbf{x}, \mathcal{X})$ means the set of link flow vectors which are obtained by a convex combination of \mathbf{x} and any of the target vectors in \mathcal{X} .³ Equivalently, \mathbf{x} is a restricted equilibrium if none of the targets in \mathcal{X} lead to improving directions in the sense that the total system travel time would be reduced by moving to some $\mathbf{x}_i^* \in \mathcal{X}$ while fixing the travel times at their current values. That is,

$$\mathbf{t}(\mathbf{x}) \cdot (\mathbf{x} - \mathbf{x}_i^*) \leq 0 \quad \forall \mathbf{x}_i^* \in \mathcal{X} \quad (9.33)$$

At a high level, simplicial decomposition works by iterating between adding new target vectors to \mathcal{X} , and then finding a restricted equilibrium using the current vectors in \mathcal{X} . In this light, Frank-Wolfe can be seen as a special case of simplicial decomposition, where \mathcal{X} only consists of the current target vector (forgetting any from past iterations), because if there is only one target $\mathbf{x}^* \in \mathcal{X}$

³A flow vector \mathbf{x}' is a convex combination of \mathbf{x} and the target vectors in \mathcal{X} if there exist nonnegative constants $\lambda_0, \lambda_1, \dots, \lambda_k$ such that $\mathbf{x}' = \lambda_0 \mathbf{x} + \lambda_1 \mathbf{x}_1^* + \dots + \lambda_k \mathbf{x}_k^*$.

then the restricted equilibrium in $X(\mathbf{x}, \mathcal{X})$ is simply the restricted equilibrium along the line segment connecting \mathbf{x} to \mathbf{x}^* .

In practice, it is too expensive to exactly find a restricted equilibrium at each iteration. Instead, several “inner iteration” steps are taken to move towards a restricted equilibrium with the current set \mathcal{X} before looking to add another target. In each inner iteration, the current solution \mathbf{x} is adjusted to $\mathbf{x} + \mu \Delta \mathbf{x}$, where μ is a step size and $\Delta \mathbf{x}$ is a direction which moves toward restricted equilibrium. Smith (1984) shows that one good choice for this direction is

$$\Delta \mathbf{x} = \frac{\sum_{\mathbf{x}_i^* \in \mathcal{X}} [\mathbf{t}(\mathbf{x}) \cdot (\mathbf{x} - \mathbf{x}_i^*)]^+ (\mathbf{x}_i^* - \mathbf{x})}{\sum_{\mathbf{x}_i^* \in \mathcal{X}} [\mathbf{t}(\mathbf{x}) \cdot (\mathbf{x} - \mathbf{x}_i^*)]^+} \quad (9.34)$$

This rather intimidating-looking formula is actually quite simple. It is nothing more than a weighted average of the directions $\mathbf{x}_i^* - \mathbf{x}$ (potential moves toward each target in \mathcal{X}), where the weight for each potential direction is the extent to which it improves upon the current solution: $[\mathbf{t}(\mathbf{x}) \cdot (\mathbf{x} - \mathbf{x}_i^*)]$ is the reduction in total system travel time obtained by moving from \mathbf{x} to \mathbf{x}_i^* while holding travel times constant. If this term is negative, there is no need to move in that direction, so the weight is simply set to zero. The denominator is simply the sum of the weights, which serves as a normalizing factor.

The step size μ is chosen through trial-and-error. One potential strategy is to iteratively test μ values in some sequence (say, $1, 1/2, 1/4, \dots$) until we have found a solution acceptably closer to restricted equilibrium than \mathbf{x} .⁴ “Acceptably closer” can be calculated using the *Smith gap*

$$\gamma_S = \sum_{\mathbf{x}_i^* \in \mathcal{X}} ([\mathbf{t}(\mathbf{x}) \cdot (\mathbf{x} - \mathbf{x}_i^*)]^+)^2 \quad (9.35)$$

which is similar to the gap measures described in Chapter 7 in that it is zero if and only if \mathbf{x} is a restricted equilibrium, and positive otherwise. It reflects a restricted equilibrium because the summation is only over the vectors in \mathcal{X} : if any of them reflect an “improvement” over the current solution, in that total travel time would be reduced (if link travel times were held fixed), the corresponding term in square brackets is positive, and $\gamma_S > 0$. It can be shown that taking a small enough step in the direction (9.34) will reduce the Smith gap, and a line search or trial-and-error method can be used to determine what an acceptable step size is. Squaring the term in brackets ensures that γ_S is differentiable, which plays a role in deriving the convergence rate of this method.

The Smith gap could thus be converted into a stopping criterion for the (unrestricted) traffic assignment problem, by extending the sum to include every possible all-or-nothing assignment, not just the ones in \mathcal{X} , but this idea will not be explored further in this book.

Putting all of this together, the simplicial decomposition algorithm can be stated as:

⁴Those familiar with nonlinear optimization may see parallels between this and the Armijo rule.

1. Initialize the set $\mathcal{X} \leftarrow \emptyset$
2. Find shortest paths for all OD pairs.
3. Form the all-or-nothing assignment \mathbf{x}^* based on shortest paths.
4. If \mathbf{x}^* is already in \mathcal{X} , stop.
5. Add \mathbf{x}^* to \mathcal{X} .
6. **Subproblem:** Find a restricted equilibrium \mathbf{x} using only the vectors in \mathcal{X} .
 - (a) Find the improvement direction $\Delta \mathbf{x}$ using equation (9.34).
 - (b) Update $\mathbf{x} \leftarrow \mathbf{x} + \mu \Delta \mathbf{x}$, with μ sufficiently small (to reduce γ_S).
 - (c) Update travel times.
 - (d) Return to step 1 of subproblem unless γ_S is small enough.
7. Return to step 2.

Below we apply this algorithm to the example in Figure 9.5, choosing μ via trial and error from the sequence $1/2, 1/4, \dots$ and stopping at the first value that reduces γ_S . Other ways of choosing μ are also possible.

Iteration 1. We set

$$\mathbf{x}_1^* = [6 \quad 0 \quad 6 \quad 0 \quad 6]$$

and $\mathcal{X} = \{\mathbf{x}_1^*\}$. For the subproblem, the only possible solution is $\mathbf{x} = \mathbf{x}_1^*$, which has $\gamma_S = 0$ (it is trivially a restricted equilibrium) and travel times are

$$\mathbf{t} = [60 \quad 53 \quad 16 \quad 53 \quad 60] .$$

Iteration 2. The new all-or-nothing assignment is

$$\mathbf{x}_2^* = [0 \quad 6 \quad 0 \quad 0 \quad 6] ,$$

and $\mathcal{X} = \{\mathbf{x}_1^*, \mathbf{x}_2^*\}$. For the first iteration of the subproblem, notice that $\mathbf{t} \cdot (\mathbf{x} - \mathbf{x}_1^*) = 0$ and $\mathbf{t} \cdot (\mathbf{x} - \mathbf{x}_2^*) = 138$, so Smith's formula (9.34) reduces to

$$\Delta \mathbf{x} = \frac{0}{138}(\mathbf{x}_1^* - \mathbf{x}) + \frac{138}{138}(\mathbf{x}_2^* - \mathbf{x}) = [-6 \quad 6 \quad -6 \quad 0 \quad 0] .$$

Taking a step of size $\mu = 1/2$ gives us

$$\mathbf{x} = [3 \quad 3 \quad 3 \quad 0 \quad 6] .$$

The new travel times are

$$\mathbf{t} = [30 \quad 54\frac{1}{2} \quad 14\frac{1}{2} \quad 53 \quad 60] ,$$

so $\mathbf{t} \cdot \mathbf{x} = 657$, $\mathbf{t} \cdot \mathbf{x}_1^* = 627$, and $\mathbf{t} \cdot \mathbf{x}_2^* = 687$, so

$$\gamma_S = ([657 - 627]^+)^2 + ([657 - 687]^+)^2 = 30^2 + 0^2 = 900 .$$

Assume this is “small enough” to complete the subproblem.

Iteration 3. The new all-or-nothing assignment is

$$\mathbf{x}_3^* = [6 \quad 0 \quad 0 \quad 6 \quad 0]$$

and $\mathcal{X} = \{\mathbf{x}_1^*, \mathbf{x}_2^*, \mathbf{x}_3^*\}$. For the first iteration of the subproblem, calculate $\mathbf{t} \cdot (\mathbf{x} - \mathbf{x}_1^*) = 30$, $\mathbf{t} \cdot (\mathbf{x} - \mathbf{x}_2^*) = -30$, and $\mathbf{t} \cdot (\mathbf{x} - \mathbf{x}_3^*) = 159$. So (9.34) gives

$$\begin{aligned} \Delta \mathbf{x} &= \frac{30}{189}(\mathbf{x}_1^* - \mathbf{x}) + \frac{0}{189}(\mathbf{x}_2^* - \mathbf{x}) + \frac{159}{189}(\mathbf{x}_3^* - \mathbf{x}) \\ &= [3 \quad -3 \quad -2.05 \quad 5.05 \quad -5.05] . \end{aligned}$$

Taking a step of size $\mu = 1/2$ would give

$$\mathbf{x} = [4.5 \quad 1.5 \quad 1.98 \quad 2.52 \quad 3.48]$$

and

$$\mathbf{t} = [45 \quad 52.5 \quad 12.7 \quad 54.3 \quad 47.4]$$

which has $\gamma_S = 1308$.

Assume that this is no longer “small enough” to return to the master problem, so we begin a second subproblem iteration. Smith’s formula (9.34) now gives $\Delta \mathbf{x} = 0.2(\mathbf{x} - \mathbf{x}_2^*) + 0.8(\mathbf{x} - \mathbf{x}_3^*)$, and the trial solution $\mathbf{x} + \frac{1}{2}\Delta \mathbf{x}$ has $\gamma_S = 1195$, which is an improvement. In this case, choosing a smaller μ would work even better; for instance $\mu = 1/4$ would reduce the Smith gap to 557. There is thus a tradeoff between spending more time on finding the “best” value of μ , or spending more time on finding new search directions and vectors for \mathcal{X} . Balancing these is an important question for implementation.

The algorithm can continue from this point or terminate if this average excess cost is small enough.

9.3 Stochastic User Equilibrium

This section describes another extension to the basic TAP. To this point in the text, we have been using the principle of equilibrium to determine link and path flows, requiring all used paths between the same origin and destination to have equal and minimal travel time. We derived this principle by assuming that all drivers choose the least-travel time path between their origin and destination. However, this assumption implicitly requires *drivers to have perfect knowledge of the travel times on all routes in the network*. In reality, we know this is not true: do you know the travel times on literally *all* routes between an origin and destination? And can you accurately distinguish between a route with a travel time of 16 minutes, and one with a travel time of 15 minutes and 59 seconds? Relaxing these assumptions leads us to the *stochastic user equilibrium* (SUE) model.

In SUE, rather than requiring that each driver follow the true shortest path between each origin and destination, we assume that drivers follow the path they *believe* to be shortest, but allow for some perception error between their belief and the actual travel times. An alternative, mathematically equivalent, interpretation (explained below) is that drivers do in fact perceive travel times accurately, but care about factors other than travel time. This section explains the development of the SUE model. The mathematical foundation for the SUE model is in discrete choice concepts, which are briefly reviewed in Section 9.3.1. The specific application of discrete choice to the route choice decision is taken up in Section 9.3.2.

These sections address the “individual” perspective of logit route choice. The next steps to creating the SUE model are an efficient network loading model (a way to find the choices of *all* drivers efficiently), and then finally the equilibrium model which combines the network loading with updates to travel times, to account for the mutual dependence between travel times and route choices. These are undertaken in Sections 9.3.4 and 9.3.5, respectively. For the most part, this discussion assumes a relatively simplistic model for perception errors in travel times; Section 9.3.7 briefly discusses how more general situations can be handled.

9.3.1 Discrete choice modeling

This section provides a brief overview of discrete choice concepts. The application to route choice is in the following section. Discrete choice is a large area of scholarly inquiry in and of itself, and so the discussion is restricted to what is needed in the chapter.

Consider an individual who must make a choice from a set of options. For instance, when purchasing groceries, you must choose one store from a set of alternatives (the grocery stores in your city). When dressing in the morning, you must choose one set of clothes among all of the clothing you own. And, more relevant to transportation, when choosing to drive from one point to another, you must choose one route among all of the possible routes connecting your origin to your destination.

Mathematically, let I be a finite set of alternatives. Each alternative $i \in I$ is associated with a *utility* value U_i representing the amount of happiness or satisfaction you would have if you were to choose option i . We assume that you would choose an alternative $i^* \in \arg \max_i \{U_i\}$ which maximizes the utility you receive. Now, the utility U_i consists of two parts: an *observable utility* V_i , and an *unobservable utility* denoted by the random variable ϵ :

$$U_i = V_i + \epsilon_i \quad (9.36)$$

The difference between observable and unobservable utility can be explained in different ways. One interpretation is that V_i represents the portion of the utility that is due to objective factors visible to the modeler; when choosing a grocery store, that might include the distance from your home, the price,

the variety of items stocked, etc. The unobserved utility consists of subjective factors that the modeler cannot see (or chooses not to include in the model), even though they are real insofar as they affect your choice. In the grocery store example, this might include your opinion on the taste of the store brands, the cleanliness of the store, and so on. Then, by modeling the unobserved utility as a random variable ϵ , we can express choices in terms of probabilities. (The modeler does not know all of the factors affecting your choice, so they can only speak of probabilities of choosing different options based on what is observable.)

A second interpretation is that the observed utility V_i actually represents all of the factors that you care about. However, for various reasons you are incapable of knowing all of these reasons with complete accuracy. (You probably have a general sense of the prices of items at a grocery store, but very few know the exact price of every item in a store's inventory.) Then the random variable ϵ_i represents the *error* between the true utility (V_i) and what you believe the utility to be (U_i). Either interpretation leads to the same mathematical model.

Depending on the distribution we choose for the random variables ϵ_i , different discrete models are obtained. A classic is the *logit* model, which is obtained when the unobserved utilities ϵ_i are assumed to be independent across alternatives, and to have Gumbel distributions with zero mean and identical variance. Under this assumption, the probability of choosing alternative i is given by

$$p_i = \frac{\exp(\theta V_i)}{\sum_{j \in I} \exp(\theta V_j)} \quad (9.37)$$

where θ is a nonnegative parameter reflecting the relative magnitude of the observed utility relative to the unobserved utility. Notice what happens in this formula as θ takes extreme values: if $\theta = 0$, then all terms in the numerator and denominator are 1, and the probability of choosing any alternative is exactly the same. (Interpretation: the unobserved utility ϵ is much more important than the observed utility, so the observed utility has no impact on the decision made. Since the unobserved utility has the same distribution for every alternative, each is equally likely.) Or, if θ grows large, then the denominator of (9.37) will be dominated by whichever terms have the largest observed utility V_i . If there is some alternative i^* for which $V_{i^*} > V_i$ for all $i \neq i^*$ (the observed utility for i^* is strictly greater than any other alternative), then as $\theta \rightarrow \infty$, the probability of choosing i^* approaches 1 and the probability of choosing any other alternative approaches 0. Another important consequence of (9.37) is that the probability of choosing any alternative is always strictly positive; there is *some* chance that the unobserved utility will be large enough that any option could be chosen.

While the logit model is nice in that we have a closed-form expression (9.37) for the probabilities of choosing any alternative, the logit assumptions are very strong — particularly the assumption that the ϵ_i are independent and identically distributed. The exercises at the end of the chapter explore some examples demonstrating how these assumptions can lead to unreasonable results. Another common assumption is that the ϵ are drawn from a multivariate normal distribution, which allows for correlation among the unobserved utilities for different

alternatives. This leads to the *probit* choice model, which is more flexible and arguably realistic. However, unlike the logit model, the probit model does not have a closed-form expression for probabilities like (9.37). Instead, Monte Carlo sampling methods are used to estimate choices.

The majority of this section is focused on logit-based models. While probit models are more general and arguably more realistic, the logit model has two major advantages from the perspective of a book like this. First, computations in logit models can often be done analytically, simplifying explanations and making it possible to give examples you can easily verify. This helps you better understand the main ideas in stochastic user equilibrium and build intuition. Second, logit models admit faster solution algorithms, algorithms which scale relatively well with network size. This is an important practical advantage for logit models. Nevertheless, Section 9.3.7 provides some discussion on probit and other models and what needs to change from the logit discussion below.

9.3.2 Logit route choice

This section specializes the discrete choice framework from the previous section to route choice in networks. Consider a traveler leaving origin r for destination s . They must choose one of the paths π connecting r to s , that is, they must make a choice from the set Π_{rs} . The most straightforward way to generalize the principle of user equilibrium to account for perception errors is to set the observed utility equal to the negative of path travel time, so

$$U^\pi = -c^\pi + \epsilon^\pi \quad (9.38)$$

with the negative sign indicating that *maximizing utility* for drivers means *minimizing travel time*. Assuming that the ϵ^π are independent, identically distributed Gumbel random variables, we can use the logit formula (9.37) to express the probability that path π is chosen:

$$p^\pi = \frac{\exp(-\theta c^\pi)}{\sum_{\pi' \in \Pi_{rs}} \exp(-\theta c_{\pi'})} \quad (9.39)$$

The comments in the previous section apply to the interpretation of this formula. As θ approaches 0, drivers' perception errors are large relative to the path travel times, and each path is chosen with nearly equal probability. (The errors are so large, the choice is essentially random.) As θ grows large, perception errors are small relative to path travel times, and the path with lowest travel time is chosen with higher and higher probability. At any level of θ , there is a strictly positive probability that each path will be taken.

For concreteness, the route choice discussion so far corresponds to the second interpretation of SUE, where the unobserved utility represents perception errors in the utility. The first interpretation would mean that ϵ^π represents factors other than travel time which affect route choice (such as comfort, quality of scenery, etc.). Either of these interpretations is mathematically consistent with the discussion in Section 9.3.1.

The fact that the denominator of (9.39) includes a summation over *all* paths connecting r to s is problematic, from a computation standpoint. The number of paths can grow exponentially with network size. Any use of stochastic user equilibrium in a practical setting, therefore, requires a way to compute link flows without explicitly calculating the sum in (9.39).

This is done by carefully defining which paths are in the choice set for travelers. The notation Π^{rs} in (9.39) in this book means the set of all acyclic paths connecting origin r to destination s . Following Chapter 7, we will use the notation $\hat{\Pi}^{rs}$ to define the set of paths being considered by travelers in SUE; these sets are sometimes called sets of *reasonable paths*. With a suitable definition of this set, using $\hat{\Pi}^{rs}$ in place of Π^{rs} in equation (9.39) leads to tractable computation schemes.

Two possibilities are common: selecting an acyclic subset of links, and choosing $\hat{\Pi}^{rs}$ to contain the paths using these links only; or setting $\hat{\Pi}^{rs}$ to consist of literally *all* paths (even cyclic ones) connecting origin r to destination s . Both of these are discussed next. The key to both of these definitions of $\hat{\Pi}^{rs}$ is that we can determine how many of the travelers passing through a given node came from each of the available incoming links, *without needing to know the specific path they are on*. This is known as the Markov property, and is discussed at more length in the optional Section 9.3.3.

Totally acyclic paths

For a particular origin r and destination s , choose a set of allowable links \hat{A}^{rs} , and let $\hat{\Pi}^{rs}$ consist of all paths starting at r , ending at s , and only containing links from \hat{A}^{rs} . We require two conditions on the set of allowable links:

1. There is at least one path from r to s using allowable links; this ensures that $\hat{\Pi}^{rs}$ is nonempty.
2. The set of allowable links contains no cycle; this ensures that all paths in $\hat{\Pi}^{rs}$ are acyclic.

We say a set of paths $\hat{\Pi}^{rs}$ is *totally acyclic* if it can be generated from an allowable link set satisfying these conditions. This definition is closely related to the idea of a bush from Section 6.2.3. If \hat{A}^{rs} contains a path from r to every destination, it is also a bush; and if furthermore the sets \hat{A}^{rs} are the same for all destinations s , we can do the network loading for all the travelers leaving origin r simultaneously, rather than separately for each destination.

Note that there are collections of acyclic paths which are not totally acyclic. In the network in Figure 9.6, if we choose $\hat{\Pi}^{14} = \{[1, 2, 3, 4], [1, 3, 2, 4]\}$, both paths in this set are acyclic, but the set of allowable links needs to include every link in the network:

$$\hat{A}^{14} = \{(1, 2), (1, 3), (2, 3), (2, 4), (3, 2), (3, 4)\}$$

This set contains the cycle $[2, 3, 2]$, so it is not possible to generate a reasonable path set containing $[1, 2, 3, 4]$ and $[1, 3, 2, 4]$ from an acyclic set of allowable links.

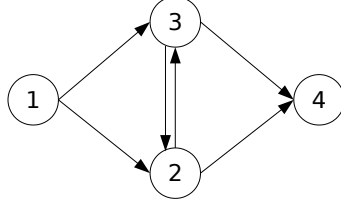


Figure 9.6: A set of acyclic paths need not be totally acyclic.

The advantage of totally acyclic path sets is that we can define a topological order on the nodes (see Section 2.2). With this topological order, we can efficiently make computations using the logit formula (9.39) without having to enumerate all the paths. This procedure is described in the next section.

We next describe two ways to form sets of totally acyclic paths. For each link, define a quantity c_{ij}^0 which is constant and independent of flow — examples include the free-flow travel time or distance on the link. For each origin r and node i , let L_i^r denote the length of the shortest path from r to i , using the quantities c_{ij}^0 as the link costs. Likewise, for each destination s and node i , let ℓ_i^s denote the length of the shortest path from i to s , again using the quantities c_{ij}^0 as the link costs.

Consider the following sets of paths:

1. The set of all paths, for which the head node of each link is further away from the origin than the tail node, based on the quantities c_{ij}^0 . That is, the sets $\hat{\Pi}^r$ containing all paths starting at r and satisfying $L_j^r > L_i^r$ for each link (i, j) in the path.
2. The set of all paths, for which the head node of each link is closer to the destination than the tail node, based on the quantities c_{ij}^0 . That is, the sets $\hat{\Pi}^s$ containing all paths ending at s satisfying $\ell_j^s < \ell_i^s$ for each link (i, j) in the path. (This is like the first one, but oriented toward the destination, rather than the origin.)
3. The set of paths which satisfy both of the above criteria: $\hat{\Pi}^{rs} = \hat{\Pi}^r \cap \hat{\Pi}^s$.

For instance, if c_{ij}^0 reflects the physical length on each link, then for a given origin, $\hat{\Pi}^r$ would consist of all of the paths which start at that origin and always move away from it, never “doubling back.” Likewise, $\hat{\Pi}^s$ would consist of all paths which always move closer to their destination node s , without any diversions that lead it away. The third option has paths which both continually move away from their origin and toward their destination. Exercise 15 asks you to show that all three possibilities for $\hat{\Pi}$ are totally acyclic.

Figure 9.7 illustrates these three definitions. The top of the figure shows a network with node A as origin and node F as destination, and the links

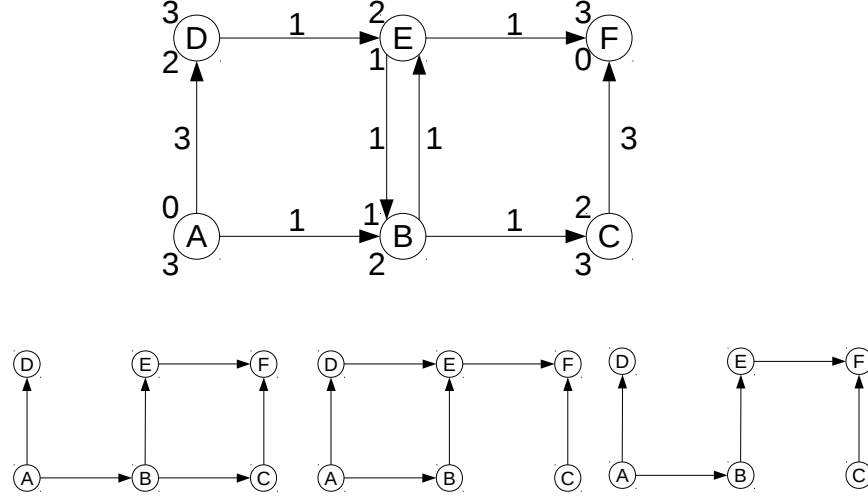


Figure 9.7: Demonstration of three alternative definitions of totally acyclic paths. Top panel shows c_{ij}^0 values on links, shortest path distances L_i from A above nodes, and shortest path distances ℓ_i to F below nodes.

are labeled with their c_{ij}^0 values. The nodes are labeled with their L_i values (above each node) and ℓ_i values below. The bottom of the figure shows the links satisfying each of the three criteria ($L_j > L_i$; $\ell_j < \ell_i$; and both of these simultaneously). The paths in these networks are the allowable paths in the original network. Notice that in all cases, there are no cycles in these links (even though the original network had the cycle [2,5,2]).

Of these principles, the third imposes stricter conditions on which paths are in the reasonable set. The first and second are weaker, and includes some paths which may not seem reasonable to you. For instance, the spiral path in Figure 9.8 satisfies the first condition, since the distance from the origin is always increasing. However, it does not satisfy the third condition, since at times the distance to the destination increases as well.

However, a major advantage of the first two principles is that we can aggregate travelers by origin or destination. With the first principle, the destination of travelers can be ignored for routing purposes — if a path is reasonable for a travel from an origin r to a node i , that path segment is reasonable for travel to any node beyond i as well. This allows us to aggregate travelers by origin (as in Section 6.2.3) and calculate a “one-to-all” path set for each origin, rather than having separate path sets for each OD pair. A similar destination-based aggregation is possible with the second principle.

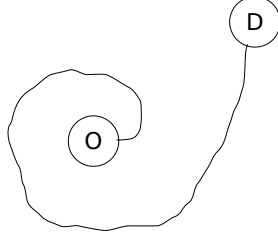


Figure 9.8: Is this path reasonable?

Full cyclic path set

Instead of restricting the path set to create a totally acyclic collection, an alternative is to have $\hat{\Pi}^{rs}$ consist of literally all paths from origin r to destination s , *even including cycles*. This will often mean these sets are infinite. For example, consider the network in Figure 9.6. Under this definition, the (cyclic) path $[1, 2, 3, 2, 4]$ is part of $\hat{\Pi}^{rs}$, as is $[1, 2, 3, 2, 3, 2, 4]$, and so on.

Including cyclic paths, especially paths with arbitrarily many repetitions of cycles, may seem counterintuitive. There are several reasons why this definition of $\hat{\Pi}^{rs}$ is nevertheless useful. One reason is that requiring total acyclicity is in fact quite a strong condition. In Figure 9.6, there is no totally acyclic path set that includes both $[1, 2, 3, 4]$ and $[1, 3, 2, 4]$ as paths — if we want to allow one path as reasonable, then by symmetry the other should be reasonable as well. But any path set including both of those includes both links $(2, 3)$ and $(3, 2)$, which form a cycle.

So, it is desirable to have an alternative to total acyclicity that still does not require path enumeration. As shown in the following section, it is possible to compute the link flows resulting from (9.39) without having to list all the paths, if all cyclic paths are included. The intuition is that all travelers at a given node can be treated identically in terms of which link they move to next: in Figure 9.6, we can split the vehicles arriving at node 2 between links $(2, 3)$ and $(2, 4)$ without having to distinguish whether they came via link $(1, 2)$, as if on the path $[1, 2, 4]$ or $[1, 2, 3, 4]$, or whether they came via link $(3, 2)$, as if on the path $[1, 3, 2, 4]$, or even $[1, 2, 3, 2, 4]$.

Furthermore, some modelers are philosophically uncomfortable with including restrictions like those in the previous section, without evidence that those rules really represent traveler behavior. Determining which sets of paths travelers actually consider (and why) is complicated, and still not well-understood.⁵ One school of thought is that it is therefore better to impose no restrictions at all, essentially taking an “agnostic” position with respect to the sets $\hat{\Pi}^{rs}$, rather than imposing restrictions which may not actually represent real behavior.

⁵Emerging data sources, such as Bluetooth readers, are providing more complete information on observed vehicle trajectories. This may provide more insight on this subject.

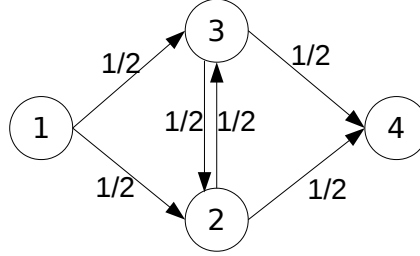


Figure 9.9: Flows from stochastic network loading when all links have unit cost, $\theta = \log 2$, and the full cyclic path set is allowed.

As an example, assume that every link in Figure 9.6 has the same travel time of 1 unit, and that $\theta = \log 2$. Then there are two paths of length 2 ($[1,2,4]$ and $[1,3,4]$), two paths of length 3 ($[1,2,3,4]$ and $[1,3,2,4]$), two paths of length 4 ($[1,2,3,2,4]$ and $[1,3,2,3,4]$), and so on. Therefore the denominator in the logit formula is

$$\sum_{\pi \in \tilde{\Pi}} \exp(-\theta c^\pi) = \frac{1}{2} + \frac{1}{2} + \frac{1}{4} + \frac{1}{4} + \frac{1}{8} + \frac{1}{8} + \cdots = 2 \quad (9.40)$$

and the probability of choosing one of the length-2 paths is $(1/2)/2 = 1/4$, the probability of choosing one of the length-3 paths is $1/8$, and so on. The flows on each link can be calculated by multiplying these path flows by the number of times that path uses a link. For instance, to calculate the flow on link $(1,2)$, observe that it is used by paths $[1,2,4]$, $[1,2,3,4]$, $[1,2,3,2,4]$, $[1,2,3,2,3,4]$, and so on, with respective probabilities $1/4$, $1/8$, $1/16$, $1/32$, etc. Thus the total flow on this link is the sum of these, or $1/2$. The flows on links $(1,3)$, $(2,4)$, and $(3,4)$ are also found to be $1/2$ by the same technique. Calculating the flow on links $(2,3)$ and $(3,2)$ is trickier, because some paths use these multiple times. For example, path $[1,2,3,2,3,4]$ uses link $(2,3)$ twice, so even though the probability of selecting this path is $1/2^5 = 1/32$, it actually contributes twice this ($1/16$) to the flow. It is possible to show that the flow on these links is also $1/2$, giving the final flows in Figure 9.9.

As this example shows, direct calculations involving this path set usually involve summing infinite series. As will show in Section 9.3.4, there is an alternative method that allows us to make these computations without explicitly calculating such sums.

9.3.3 (*) The Markov property and the logit formula

(This optional section gives mathematical reasons why totally acyclic path sets, and complete path sets, both allow for efficient calculation of the logit formula.)

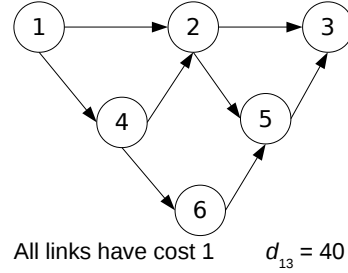


Figure 9.10: Network for demonstrating the Markov property.

Both of these path set definitions allow for the computation of the logit formula to be disaggregated by node and by link, without having to enumerate the paths in the network. The key to this is the *Markov property*. An informal statement of this property is that if we randomly select a traveler passing through a node, and are curious about the probability that their path leaves that node, there is no information provided by knowing how they arrived to that node.

For example, consider the network in Figure 9.10, where the demand is $d_{13} = 40$ vehicles and $\theta = \log 2$. All links have unit cost. Assume first that all paths in this network are allowed. Then the right panel of Figure 9.11 shows the flow on each path, and the left panel shows the flow on each link.

In this network 18 vehicles pass through node 2. Suppose we pick one of them at random, and want to know the probability that the next link in this vehicle's path is (2,3), as opposed to (2,5). From examining the path flows in Figure 9.11, we can see that this probability is $(8 + 4)/(8 + 4 + 4 + 2) = 2/3$ (ignoring the flow on path [1,4,6,5,3] in the denominator, since these trips do not pass through node 2). Now, suppose that these vehicles also reported the segment of their path that led them to node 2 — that is, they also report whether they came via segment [1,2] or segment [1,4,2]. Does this change our answers in any way?

If we know they came from segment [1,2], then they are either on path [1,2,3] or [1,2,5,3], and the probability that they continue on (2,3) is $8/(8+4) = 2/3$. If we know they came from segment [1,4,2], then they are either on path [1,4,2,3] or path [1,4,2,5,3], and the probability that they continue on (2,3) is $4/(4+2)$, which is still $2/3$. So knowing the first segment of their trip does not provide any additional information as to the remaining segment.

The situation changes if we modify the allowable path set to only include three paths, [1,2,3], [1,2,5,3], and [1,4,2,3].⁶ Here Figure 9.12 shows the corresponding path flows and link flows.

Let us ask the same question of the travelers passing through node 2. With-

⁶A natural way this path set might arise is to include paths that are only within a small threshold of the shortest path cost, thus including paths with cost 2 or 3 but excluding paths with cost 4.

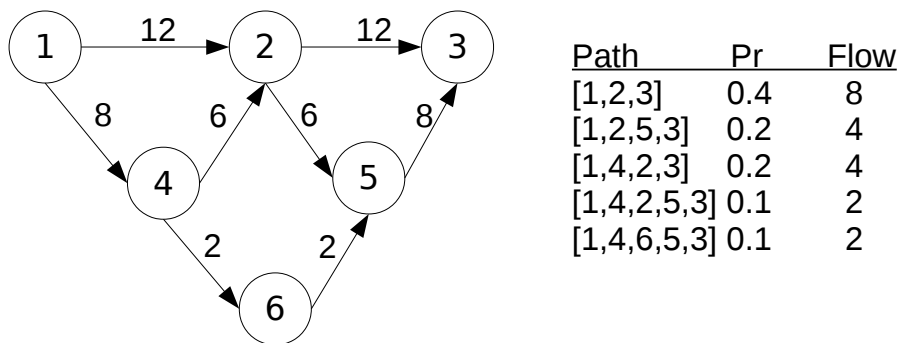


Figure 9.11: Link and path flows when all paths are allowed.

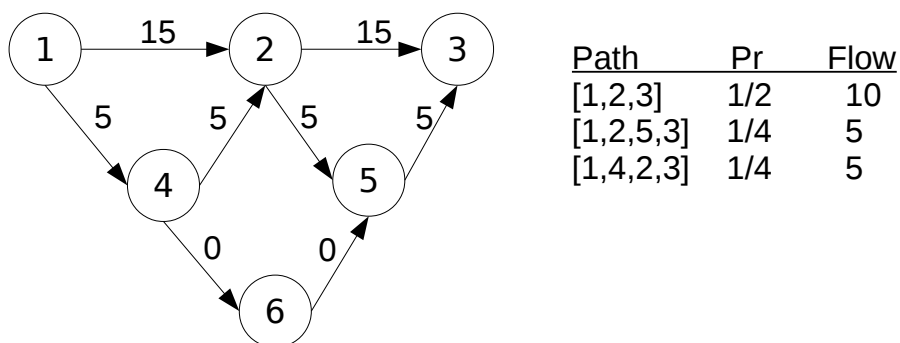


Figure 9.12: Link and path flows when only a subset of paths is allowed.

out knowing anything further, the probability that they continue on link (2,3) is $(10 + 5)/20 = 3/4$. However, if we know they came from [1,2], then the probability that they continue on (2,3) is $10/15 = 2/3$. If we know they came from [1,4,2], then the probability that they continue on (2,3) is 1, and there is no other option! So in this case, knowing the first segment of the path *does* give us additional information about the rest of their journey.

It turns out that the Markov property will be very useful, and will allow us to efficiently evaluate the logit formula without enumerating paths. Informally, we can do computations using just the link flows (the left panels in Figures 9.11 and 9.12) without having to use the path flows (the right panels of these figures) — we can get the path flows on the right from the link flows on the left. In large networks, the link-based representation is much more compact and efficient.

A more formal statement of this property is as follows. To keep the formulas clean, assume that there is a single origin r and destination s ; in a general network, we can apply the same logic separately to each OD pair. In logit assignment, the path set $\hat{\Pi}$ is said to satisfy the Markov property if there exist values P_{ij} for each link such that

$$p^\pi = \frac{\exp(-\theta c^\pi)}{\sum_{\pi' \in \hat{\Pi}} \exp(-\theta c^{\pi'})} = \prod_{(i,j) \in \pi} P_{ij}^{\delta_{ij}^\pi} \quad (9.41)$$

where δ_{ij}^π is the number of times path π uses link (i, j) . That is, that the probability of a traveler selecting any path can be computed by multiplying P_{ij} values across its links. The P_{ij} values can be interpreted as conditional probabilities: given that a traveler is passing through i , they express the probability that their path leaves that node through link (i, j) .

The segment substitution property

To make the connection between totally acyclic and complete path sets and the Markov property, we first show that both path set choices satisfy the *segment substitution property*. Given any path $\pi = [r, i_1, i_2, \dots, s]$, a *segment* σ is any contiguous subset of one or more of its nodes. For example, the path $[1, 2, 3, 2, 4]$ contains $[1, 2, 3]$, $[2, 3]$, $[2]$, and $[2, 3, 2, 4]$. It does not contain the segment $[1, 4]$; even though both of those nodes are in the path, they do not appear consecutively. Note that a segment can consist of a single link, such as $[2, 3]$, or even a single node, such as $[2]$. We use the notation \oplus to indicate joining segments, so $[1, 2, 3, 2, 4] = [1, 2, 3] \oplus [3, 2, 4]$. If two segments are being joined, the end node of the first must match the starting node of the second.

The set of reasonable paths $\hat{\Pi}$ satisfies the segment substitution property if, for any pair of reasonable paths which pass through the same two nodes, the paths formed by exchanging the segments between those nodes are also reasonable. That is, if $\pi = \sigma_1 \oplus \sigma_2 \oplus \sigma_3$ is reasonable, and if there is another reasonable path $\pi' = \sigma'_1 \oplus \sigma'_2 \oplus \sigma'_3$ with σ_2 and σ'_2 starting and ending at the same node, then the paths $\sigma_1 \oplus \sigma'_2 \oplus \sigma_3$ and $\sigma'_1 \oplus \sigma_2 \oplus \sigma'_3$ are also reasonable.

With the allowable path set in Figure 9.12, the segment substitution property is not satisfied. There are paths $[1,2,5,3]$ and $[1,4,2,3]$, both passing through nodes 1 and 2, which can be decomposed in the following way:

$$[1, 2, 5, 3] = [1] \oplus [1, 2] \oplus [2, 5, 3] \quad (9.42)$$

$$[1, 4, 2, 3] = [1] \oplus [1, 4, 2] \oplus [2, 3] \quad (9.43)$$

Notice that the corresponding pairs of segments on the right-hand sides all start and end at the same nodes. We can generate two new paths by “crossing” the middle segments of (9.42) and (9.43): $[1,2,3]$ and $[1,4,2,5,3]$. The first of these is allowable, but the second is not.

By contrast, you can verify that segment substitution is satisfied for the allowable path set in Figure 9.11. No matter which pairs of paths you choose, swapping the segments results in another allowable path.

A reasonable path set cannot have the Markov property unless it satisfies the segment substitution property, as shown in the above example. Without the segment substitution property, we could not know how the vehicles at node i would split without knowing the specific paths they were on. Segment substitution ensures that all travelers passing through node i are considering the same set of outgoing links (and, indeed, the same set of path segments continuing on to the destination).

It is fairly easy to show that both path set definitions considered above — sets of totally acyclic paths, and the set of all paths (even cyclic ones) — satisfy the segment substitution property; see Exercise 16.

Decomposing the logit formula

Given a reasonable path set $\hat{\Pi}$ satisfying the segment substitution property, and any two nodes a and b , let Σ_{ab} denote the set of segments which start and end at these nodes, and are part of a reasonable path. Define the quantity

$$V_{ab} = \sum_{\sigma \in \Sigma_{ab}} \exp(-\theta c^\sigma) \quad (9.44)$$

where $c^\sigma = \sum_{(i,j) \in \sigma} c_{ij}$ is the travel time of a segment. In particular, V_{rs} is a sum over all of the reasonable paths from origin to destination, and corresponds to the denominator of the logit formula. Thus, the probability that a traveler chooses a particular path π is simply $\exp(-\theta c^\pi)/V_{rs}$, and the flow on this path is $d^{rs} \exp(-\theta c^\pi)/V_{rs}$. Also note the special case V_{ii} where the start and end nodes are the same. In this case Σ_{ii} consists only of the single-node segment $[i]$ with zero cost, so $V_{ii} = 1$.

Furthermore, note that the numerator of the logit formula can be factored by segment, so that if $\pi = \sigma_1 \oplus \sigma_2 \oplus \sigma_3$, we have

$$\exp(-\theta c^\pi) = \exp(-\theta(c^{\sigma_1} + c^{\sigma_2} + c^{\sigma_3})) = \exp(-\theta c^{\sigma_1}) \exp(-\theta c^{\sigma_2}) \exp(-\theta c^{\sigma_3}). \quad (9.45)$$

To calculate the flow on a link x_{ij} , we need to add the flow from all of the reasonable paths which use this link. Every reasonable path using link (i, j) takes the form $\sigma_1 \oplus [i, j] \oplus \sigma_2$, where σ_1 goes from the origin r to node i and σ_2 goes from node j to the destination s .⁷ Therefore

$$x_{ij} = \sum_{\pi \in \hat{\Pi}^{rs}} \delta_{ij}^{\pi} h^{\pi} \quad (9.46)$$

$$= \frac{d^{rs}}{V_{rs}} \sum_{\pi \in \hat{\Pi}^{rs}} \delta_{ij}^{\pi} \exp(-\theta c^{\pi}) \quad (9.47)$$

$$= \frac{d^{rs}}{V_{rs}} \sum_{\sigma_1 \in \Sigma_{ri}} \sum_{\sigma_2 \in \Sigma_{js}} \exp(-\theta(c^{\sigma_1} + t_{ij} + c^{\sigma_2})) \quad (9.48)$$

$$= \frac{d^{rs}}{V_{rs}} \sum_{\sigma_1 \in \Sigma_{ri}} \sum_{\sigma_2 \in \Sigma_{js}} \exp(-\theta c^{\sigma_1}) \exp(-\theta t_{ij}) \exp(-\theta c^{\sigma_2}) \quad (9.49)$$

$$= \frac{d^{rs}}{V_{rs}} \left(\sum_{\sigma_1 \in \Sigma_{ri}} \exp(-\theta c^{\sigma_1}) \right) \exp(-\theta t_{ij}) \left(\sum_{\sigma_2 \in \Sigma_{js}} \exp(-\theta c^{\sigma_2}) \right) \quad (9.50)$$

$$= d^{rs} \frac{V_{ri} \exp(-\theta t_{ij}) V_{js}}{V_{rs}} \quad (9.51)$$

The third equality groups the sum over paths according to the starting and ending segment. This equation for x_{ij} will be used extensively in the rest of this section.

Similarly, to find the number of vehicles passing through a particular node i (call this x_i), observe that every path through i can be divided into a segment from r to i , and a segment from i to s . Grouping the paths according to these segments and repeating the algebraic manipulations above gives the formula

$$x_i = d^{rs} \frac{V_{ri} V_{is}}{V_{rs}} \quad (9.52)$$

With equations (9.51) and (9.52) in hand, it is easy to show the Markov property holds in any reasonable path set with the segment substitution property. Define

$$P_{ij} = \frac{x_{ij}}{x_i} = \frac{\exp(-\theta t_{ij}) V_{js}}{V_{is}} \quad (9.53)$$

and then multiply these values together for the links in a path, say, $\pi =$

⁷If the link starts at the origin or ends at the destination, we may have $i = r$ or $j = s$, in which case σ_1 or σ_2 will consist of a single node, $[r]$ or $[s]$.

$[r, i_1, i_2, \dots, i_k, s]$:

$$\prod_{(i,j) \in \pi} P_{ij} = \prod \frac{\exp(-\theta t_{ij}) V_{js}}{V_{is}} \quad (9.54)$$

$$= \exp \left(-\theta \sum_{(i,j) \in \pi} t_{ij} \right) \frac{V_{i_1 s} V_{i_2 s} \cdots V_{i_k s} V_{ss}}{V_{rs} V_{i_1 s} V_{i_2 s} \cdots V_{i_k s}} \quad (9.55)$$

$$= \exp(-\theta c^\pi) \frac{V_{ss}}{V_{rs}} \quad (9.56)$$

But $V_{ss} = 1$ and V_{rs} is simply the denominator in the logit formula, so this product is exactly p^π . Therefore the logit path flow assignment in the set of reasonable paths satisfies the Markov property.

9.3.4 Stochastic network loading

The stochastic network loading problem is to determine the flows on each link x_{ij} , given their travel times t_{ij} , according to a particular discrete choice model. This section describes how to do so for the logit model, using formula (9.39). For stochastic network loading, we assume that these travel times are fixed and constant, and therefore unaffected by the path and link flows we calculate.

In the larger stochastic user equilibrium problem (where travel times can depend on flows), stochastic network loading plays the role of a subproblem in an iterative algorithm. This is analogous to how shortest paths and all-or-nothing loadings are often used as a subproblem in the classical traffic assignment problem: although travel times do depend on link flows, to find an equilibrium we can solve a number of shortest path problems, temporarily fixing the link costs at particular values.

Throughout this section, we are assuming a single origin r and destination s to simplify the notation. If there are many origins and destinations, these procedures should be repeated for each, and the link flows added to obtain the total link flows. (There is no harm in doing so, since we are assuming link travel times are constant, and therefore the different OD pairs do not interact with each other.) Depending on the choice of path set, it may be possible to aggregate all travelers from the same origin or destination, and load them at once. This is more efficient than doing a separate loading for each OD pair, but requires a more limited definition of the reasonable path sets.

In principle, the stochastic network loading procedure is straightforward. Given the link travel times t_{ij} , we can calculate the path travel times c^π as in Section 5.1. The path flows h^π can then be calculated from the logit formula (9.39), from which the link flows can be calculated by addition, again as in Section 5.1.

While conceptually straightforward, this procedure faces the practical difficulty that the logit formula requires summations over the set of all reasonable paths $\hat{\Pi}^{rs}$, which can grow exponentially with network size. In a realistic-sized

network, this renders the above procedure unusable, or at least computationally taxing.

This section describes how the link flows can be calculated *without* explicitly enumerating paths or using the logit formula. This is possible for both of the Π definitions discussed in the previous section: a totally acyclic set of paths, or including the full path set, even including all cycles.

These two procedures have several features in common. As was described in Section 9.3.3, both of these path set definitions have the Markov property. This section derived an important formula, repeated here:

$$x_{ij} = d^{rs} \frac{V_{ri} \exp(-\theta t_{ij}) V_{js}}{V_{rs}}, \quad (9.57)$$

where V_{ab} is the sum of $\exp(-\theta c^\sigma)$ for all reasonable path segments σ starting at node a and ending at node b . This formula is important because it allows us to calculate the flow on each link without having to enumerate all of the paths that use that link.

It is also possible to show (see Exercise 19) that we can replace t_{ij} with $L_i + t_{ij} - L_j$, where \mathbf{L} is a vector of node-specific constants in this formula; this reduces numerical errors in computations. It is common to use shortest path distances at free-flow (hence the use of the notation L). Doing this for the link and segment costs helps avoid numerical issues involved with calculating exponentials of large values. With this re-scaling, we define the *link likelihood* as

$$L_{ij} = \exp(\theta(L_j - L_i - t_{ij})) \quad (9.58)$$

and thus

$$x_{ij} = d^{rs} \frac{V_{ri} L_{ij} V_{js}}{V_{rs}}. \quad (9.59)$$

It remains to describe how the V_{ab} values can be efficiently calculated. This section shows how this can be done for both totally acyclic path sets, and the complete set of all paths (even cyclic ones). In the case of a totally acyclic path set, the relevant V values can be calculated in a single pass over the network in topological order, and then the link flows calculated in a second pass over the network in reverse topological order. In the case of the full cyclic path set, the cycles create dependencies in the link weights and in the link flow formulas which prevent them from being directly evaluated. But we can still calculate them explicitly using matrix techniques.

Totally acyclic paths

The defining feature of a totally acyclic path set is that the collection of links used by reasonable paths has no cycles. This allows us to define a topological order on the nodes, so that each allowable link connects a lower-numbered node to a higher-numbered one. In acyclic networks, it is often easy to perform calculations recursively, in increasing or decreasing topological order. Stochastic network loading is one of these cases, using Dial's method.

Define W_i and W_{ij} to be the *weight* of a node and link, respectively. The node weight W_i is a shorthand for V_{ri} , and W_{ij} is a shorthand for $V_{ri}L_{ij}$. These can be calculated recursively, using the following procedure:

1. Calculate the link likelihoods L_{ij} using equation (9.58) for all allowable links; set $L_{ij} = 0$ for any link not in the allowable set.
2. Set the current node i to be the origin r , and initialize its weight: $W_r = 1$.
3. For all links (i, j) leaving node i , set $W_{ij} = W_i L_{ij}$.
4. If i is the destination, stop. Otherwise, set i to be the next node in topological order.
5. Calculate $W_i = \sum_{(h,i) \in \Gamma^{-1}(i)} W_{hi}$ by summing the weights of the incoming links.
6. Return to step 3.

With the node and link weights in hand, we proceed to calculate the flows on each link x_{ij} and the flow through each node x_i . Using the link weights, we can rewrite equation (9.52) for the flow through each node as

$$x_i = d^{rs} \frac{V_{ri} V_{is}}{V_{rs}} = d^{rs} \frac{W_i V_{is}}{W_s}. \quad (9.60)$$

Combining with equation (9.59) for link flow, we have

$$x_{hi} = d^{rs} \frac{V_{rh} L_{hi} V_{is}}{V_{rs}} = x_i \frac{W_{hi}}{W_i} \quad (9.61)$$

so if we know the flows to some node i , we can calculate the flows to its incoming links (h, i) . So, the link flows can be calculated in reverse topological order:

1. Initialize all flows to zero: $x_i = 0$ for all nodes, and $x_{ij} = 0$ for all links.
2. Set the current node i to be the destination s , and initialize its flow: $x_i = d^{rs}$, since all vehicles must reach the destination.
3. For all links (h, i) entering node i , set $x_{hi} = x_i W_{hi} / W_i$.
4. If i is the origin, stop. Otherwise, set i to be the previous node in topological order.
5. Compute $x_i = \sum_{(i,j) \in \Gamma(i)} x_{ij}$ as the sum of the flows on outgoing links.
6. Return to step 3.

As an example, Dial's method is demonstrated on the network shown in Figure 9.13, where the demand is 2368 vehicles from node 1 to node 5 and $\theta = 1$. For convenience, the results of the calculations in this example are shown in Tables 9.1 and 9.2. As a preliminary step, we calculate the shortest

path from node 1 to all other nodes at free-flow conditions (assuming $\mathbf{x} = \mathbf{0}$) using standard techniques. Details are omitted since the technique is familiar, and the resulting shortest path labels L_i are shown in Table 9.1.

Next, we must identify the allowable links. This example will adopt the first principle from the previous section, where the allowable links (i, j) are those for which $L_i < L_j$. Every link is thus allowable except for $(2, 3)$, because this link connects node 2 (with shortest path label $L_2 = 2$) to node 3 ($L_3 = 1$). Therefore, link $(2, 3)$ will be excluded from the remaining steps of Dial's method, since no vehicles will use this link.

The next step is to calculate the link likelihoods L_{ij} for the allowable links $(i, j) \in A_B$. So, $L_{12} = \exp(2 - 0 - 2) = 1$, $L_{35} = \exp(3 - 1 - 3) = e^{-1} \approx 0.368$, and so forth. Weights are now calculated in forward topological order. The topological ordering for the reasonable bush is 1, 3, 2, 4, and 5 in that order. Notice that the original network has a cycle $[2, 3, 2]$ and so no topological order can exist on the full network. But when restricted to the allowable set, the cycle disappears and node 3 must come before node 2 topologically. For the origin, $W_1 = 1$ by definition. The weights on links leaving node 1 can now be calculated: $W_{12} = W_1 L_{12} = 1$, and $W_{13} = 1 \times 1 = 1$. Proceeding to node 3 (the next in topological order), W_3 is calculated as the sum of the weights on incoming links: $W_3 = W_{13} = 1$ and thus $W_{32} = 1 \times 1 = 1$ and $W_{35} = 1 \times e^{-1} \approx 0.368$. Node 2 is next in topological order, and its weight is the sum of the weights on its incoming links: $W_2 = W_{12} + W_{32} = 2$, so $W_{24} = 2$. Node 4 is next, and we have $W_4 = 2$ and $W_{45} = 2$. Finally, the weight of node 5 is $W_5 = W_{35} + W_{45} = 2.368$.

Link and node flows are now calculated in *reverse* topological order, starting with node 5 and then proceeding to nodes 4, 2, 3, and 1. For node 5, the node flow X_5 is simply the demand destined to this node (since there are no outgoing links), so $X_5 = 2368$. This flow is now distributed among the two incoming links $(3, 5)$ and $(4, 5)$ in proportion to their weights. So, $x_{35} = X_5 W_{35} / (W_{35} + W_{45}) = 2368 \times 0.368 / (2 + 0.368) = 368$ and $x_{45} = 2000$. Proceeding upstream, the node flow at 4 is simply the flow on link $(4, 5)$, the only outgoing link (since no vehicles have node 4 as their destination), and $X_4 = 2000$. There is only one incoming link $(2, 4)$, so $x_{24} = 2000$. (This follows trivially from the formula in step 3c since $W_4 = W_{24}$. Thus $X_2 = 2000$. Since there are two incoming reasonable links to node 2 with equal weight, they receive equal flow (again, following from the formula in 3c), setting $x_{32} = x_{12} = 1000$. Then $X_3 = x_{32} + x_{35} = 1368$, and $x_{13} = X_3 = 1368$. Finally, the node flow at the origin 1 is $X_1 = 1368 + 1000 = 2368$, as it should be.

Dial's method is now complete, having calculated the flows on each link. This method is completely consistent with (9.39), if we were to restrict attention to the reasonable paths in the network. If we were to use this formula directly, we would first enumerate the three reasonable paths $[1, 3, 5]$, $[1, 3, 2, 4, 5]$, and $[1, 2, 4, 5]$ and calculate their costs: $C^{[1,3,5]} = 4$, $C^{[1,3,2,4,5]} = C^{[1,2,4,5]} = 3$. Equation (9.39) then gives

$$p^{[1,3,5]} = \frac{e^{-4}}{e^{-4} + e^{-3} + e^{-3}} = 0.155 \quad p^{[1,3,2,4,5]} = \frac{e^{-3}}{e^{-4} + e^{-3} + e^{-3}} = 0.422$$

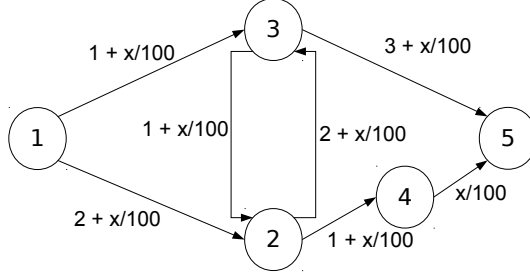


Figure 9.13: Network for demonstration of Dial's method.

Table 9.1: Node calculations in Dial's method example.

Node i	L_i	W_i	X_i
1	0	1	2368
2	2	1	1368
3	1	1	2000
4	3	2	2000
5	3	$2 + e^{-1}$	2368

$$p^{[1,2,4,5]} = \frac{e^{-3}}{e^{-4} + e^{-3} + e^{-3}} = 0.422$$

as the path choice proportions. Multiplying each of these by the total demand (2368) gives path flows

$$h^{[1,3,5]} = 368 \quad h^{[1,3,2,4,5]} = 1000 \quad h^{[1,2,4,5]} = 1000$$

As you can verify, these path flows correspond to the same link flows shown in Table 9.2.

Table 9.2: Link calculations in Dial's method example.

Link (i, j)	L_{ij}	W_{ij}	x_{ij}
(1,2)	1	1	1000
(1,3)	1	1	1368
(2,3)	0	0	0
(2,4)	1	2	2000
(3,2)	1	1	1000
(3,5)	e^{-1}	e^{-1}	368
(4,5)	1	2	2000

Full cyclic path set

We can also efficiently calculate the link flows from logit network loading, if the set of reasonable paths contains all paths, even cyclic ones. Again we will use formula (9.59). If we have an efficient way to compute V_{ab} for all pairs of nodes a and b , we can substitute them into this formula to directly obtain the link flows. Since the set of reasonable paths contains cycles, we cannot calculate these V values inductively on topological order, as was done above. Rather, a different approach is needed.

Let \mathbf{V} be the $n \times n$ matrix whose components are V_{ab} . Recall that V_{ab} is defined as the sum of $\exp(-\theta c^\sigma)$ for all segments σ starting at node a and ending at node b . We will calculate this sum by dividing the sum into segments of the same length.

We begin by calculating the part of V_{ab} which corresponds to the segments of length one (that is, the segments consisting of a single link.) If there are no parallel links in the network, then this is simply L_{ab} : there is at most one such segment, which must be $\sigma = [a, b]$, and if it exists $\exp(-\theta c^\sigma) = \exp(-\theta t_{ab}) = L_{ab}$. If it does not exist, then $L_{ab} = 0$, which is again the part of V_{ab} corresponding to segments of length one (which is empty if no such segment exists.) We can proceed similarly if there are parallel links; see Exercise 20.

Now, let \mathbf{L} be the $n \times n$ matrix whose components are L_{ab} , and form the matrix product \mathbf{L}^2 . Its components are

$$(L^2)_{ab} = \sum_{c \in N} L_{ac} L_{cb} \quad (9.62)$$

by the definition of matrix multiplication.⁸

For a given node c , the product $L_{ac} L_{cb}$ is zero unless there is a link both from a to c , and from c to b . So, we can restrict the sum in (9.62) to be over nodes c for which there are links (a, c) and (c, b) — which is precisely the nodes c for which there is a segment $[a, c, b]$ of length two. Furthermore, for such segments,

$$\exp(-\theta c^{[a,c,b]}) = \exp(-\theta[t_{ac} + t_{cb}]) = \exp(-\theta t_{ac}) \exp(-\theta t_{cb}) = L_{ac} L_{cb}. \quad (9.63)$$

So, the components of \mathbf{L}^2 are exactly the portion of the sums defining V_{ab} for segments of length two!

We demonstrate this using the example from Figure 9.6, recalling that $d^{14} = 1$, $\theta = \log 2$, and that all links have unit cost. In this network, the matrices \mathbf{L} and \mathbf{L}^2 take the form

$$\mathbf{L} = \begin{bmatrix} 0 & 1/2 & 1/2 & 0 \\ 0 & 0 & 1/2 & 1/2 \\ 0 & 1/2 & 0 & 1/2 \\ 0 & 0 & 0 & 0 \end{bmatrix} \quad \mathbf{L}^2 = \begin{bmatrix} 0 & 1/4 & 1/4 & 1/2 \\ 0 & 1/4 & 0 & 1/4 \\ 0 & 0 & 1/4 & 1/4 \\ 0 & 0 & 0 & 0 \end{bmatrix}. \quad (9.64)$$

⁸The parentheses are intentional: $(L^2)_{ab}$ is the component of matrix \mathbf{L}^2 in row a and column b . This is *not* the same as $(L_{ab})^2$, the square of the component of matrix \mathbf{L} in row a and column b .

Looking at the first row of \mathbf{L}^2 , we see that $(L^2)_{12} = 1/4$, because there is one segment which starts at node 1, ends at node 2, and contains two links ($[1,3,2]$) and $\exp(-\theta c^{[1,3,2]}) = 1/4$. Similarly $(L^2)_{13} = 1/4$, and $(L^2)_{14} = 1/2$ because there are two segments starting at node 1, ending at node 2, and containing two links ($[1,2,4]$ and $[1,3,4]$). The sum $\exp(-\theta c^{[1,2,4]}) + \exp(-\theta c^{[1,3,4]})$ is indeed $1/2$. In this matrix, $(L^2)_{23} = 0$, because there are no segments of length two starting at node 2 and ending at node 3.

Proceeding a step further, we have

$$(L^3)_{ab} = \sum_{c \in N} (L^2)_{ac} L_{cb}. \quad (9.65)$$

By the same logic, we see that this sum expresses the component of V_{ab} corresponding to segments of length three. Every segment of length three connecting a to b consists of a segment of length two connecting a to some node c , followed by a link (c, b) . Group the sum of all segments of length three by this final link, and note that $(L^2)_{ac}$ already contains the relevant portion of the sum for the first segment.

Thus, by induction, the components of matrix \mathbf{L}^n contains the portion of the sum defining \mathbf{V} corresponding to segments of length n . Therefore

$$\mathbf{V} = \mathbf{L}^0 + \mathbf{L}^1 + \mathbf{L}^2 + \mathbf{L}^3 + \cdots, \quad (9.66)$$

where the infinite sum is needed since we allow paths with an arbitrary number of cycles. Assuming that this sum exists, we can calculate it as follows:

$$\mathbf{V} = \mathbf{I} + \mathbf{L} + \mathbf{L}^2 + \mathbf{L}^3 + \cdots \quad (9.67)$$

$$= \mathbf{I} + \mathbf{L}(\mathbf{I} + \mathbf{L} + \mathbf{L}^2 + \cdots) \quad (9.68)$$

$$= \mathbf{I} + \mathbf{L}\mathbf{V}. \quad (9.69)$$

Therefore $\mathbf{V} - \mathbf{L}\mathbf{V} = \mathbf{I}$, or

$$\mathbf{V} = (\mathbf{I} - \mathbf{L})^{-1}. \quad (9.70)$$

After calculating the matrix \mathbf{V} with this formula, we can directly read off its components V_{ab} and use them to calculate the link flows using (9.59).

To complete the example, we have

$$\mathbf{V} = \begin{bmatrix} 1 & 1 & 1 & 1 \\ 0 & 4/3 & 2/3 & 1 \\ 0 & 2/3 & 4/3 & 1 \\ 0 & 0 & 0 & 1 \end{bmatrix} \quad (9.71)$$

and, for example, the flow on link (2,3) is given by

$$x_{23} = d^{14} \frac{V_{12} L_{23} V_{34}}{V_{14}} = 1 \frac{1 \cdot \frac{1}{2} \cdot 1}{1} = \frac{1}{2}. \quad (9.72)$$

Repeating this process will give the flow on every link — and unlike in the previous section, does not involve summing an infinite series term-by-term. The formula (9.59) handles all of the paths.

9.3.5 Stochastic user equilibrium

Stochastic network loading methods, described in the previous section, are the analogue of the all-or-nothing assignments we used to find \mathbf{x}^* in the method of successive averages, or Frank-Wolfe, in Chapter 7. Recall that in those methods, we identified \mathbf{x}^* by finding the link flows that would be observed if all the travel times were held constant at the current values. In TAP, all drivers aim to take the shortest path, so \mathbf{x}^* was obtained by loading all flow onto the shortest paths at the current travel times. In logit SUE, drivers do not always take the shortest path, but instead choose paths by (9.39). The methods in the previous section thus calculate an “ \mathbf{x}^* ” in the sense that it reflects the link flows which would arise if travel times were held constant.

The full development of the SUE model requires relaxing the assumption of constant travel times, in the same way as finding a traditional user equilibrium requires more than a single shortest path computation. SUE is very easy to define as a fixed point problem. Let $\mathbf{C}(\mathbf{h})$ denote the vector of path travel times as a function of the vector of path flows (this is the same as before). However, the logit formula directly gives us a complementary function $\mathbf{H}(\mathbf{c})$ giving path flows as a function of path travel times, with components

$$h^\pi = d^{rs} \frac{\exp(-\theta c^\pi)}{\sum_{\pi' \in \Pi^{rs}} \exp(\theta c_{\pi'})} \quad (9.73)$$

where (r, s) is the OD pair corresponding to path π . Therefore, the SUE problem can be expressed as follows: find a feasible path flow vector \mathbf{h}^* such that $\mathbf{h}^* = \mathbf{H}(\mathbf{C}(\mathbf{h}^*))$. This is a standard fixed-point problem. Clearly \mathbf{H} and \mathbf{C} are continuous functions if the link performance functions are continuous, and the feasible path set is compact and convex, so Brouwer’s theorem immediately gives existence of a solution to the SUE problem.

Notice that this was much easier than showing existence of an equilibrium solution to the original traffic assignment problem! For that problem, there was no equivalent of (9.73). Travelers were all using shortest paths, but if there were two or more shortest paths there was no rule for how those ties should be broken. As a result, we had to reformulate the problem as a variational inequality and introduce an auxiliary function based on movement of a point under a force. For the SUE problem, there is no need for such machinations, and we can write down the fixed point problem immediately.

However, fixed-point theorems do not offer much help in terms of actually finding the SUE solution. It turns out that convex programming and variational inequality formulations exist as well, and we will get to them shortly. But first, as a practical note, we mention that our definition of the reasonable path set $\hat{\Pi}$ should be specified *without reference to the final travel times*. The reason is that until we have found the equilibrium solution, we do not know what the link and path travel times will be. If the sets of reasonable paths vary from iteration to iteration, as the flows and travel times change, there may be problems with convergence or solution consistency. This is why the methods described above for generating totally acyclic path sets relied on constants c_{ij}^0

which were independent of flow: constants such as physical length or free-flow times. If using the full cyclic path set, there is no concern; at all iterations every path is reasonable.

Path-flow formulation and the method of successive averages

Consider the following optimization problem, developed by Caroline Fisk:

$$\min_{\mathbf{x}, \mathbf{h}} \quad \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx + \frac{1}{\theta} \sum_{\pi \in \hat{\Pi}} h^\pi \log h^\pi \quad (9.74)$$

$$\text{s.t.} \quad x_{ij} = \sum_{\pi \in \hat{\Pi}} h^\pi \delta_{ij}^\pi \quad \forall (i, j) \in A \quad (9.75)$$

$$\sum_{\pi \in \hat{\Pi}^{rs}} h^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \quad (9.76)$$

$$h^\pi \geq 0 \quad \forall \pi \in \hat{\Pi} \quad (9.77)$$

Comparing with Beckmann's formulation from Chapter 6, we see that it is identical except that the sums on path variables are now over the set of reasonable paths, rather than all paths, and that an additional term is added to the objective function. We show below that this term ensures that the optimal path flows satisfy the logit formula (9.39). This objective function is strictly convex in the *path flows* (see Exercise 23), so the SUE path flow solution is unique. This is a different situation than the classical traffic assignment problem, which has a unique equilibrium solution in link flows, but not in path flows.

First, we show that the non-negativity constraint (9.77) can be ignored without any problem, and in fact that at optimality the flow on each path is *strictly* positive. The $\log h^\pi$ term in the objective is not defined at all if a path flow is negative; if $h^\pi = 0$ we can define $h^\pi \log h^\pi = 0$, since this is the limiting value by l'Hôpital's rule. But the derivative of $h \log h$ becomes infinitely steep as h approaches zero: $\frac{d}{dh}(h \log h) = \log h + 1 \rightarrow -\infty$ as $h \rightarrow 0$. Therefore the objective function cannot be minimized at $h = 0$; no matter what the change in the other terms in the objective would be, it is better to have a very slightly positive h value than a zero one. This is a very useful property, since the optimality conditions are much simpler if there are no non-negativity constraints.

Next, as with Beckmann's formulation, we Lagrangianize the demand constraint (9.76), and substitute (9.75) in place of \mathbf{x} , to obtain a Lagrangian in terms of path flow variables only:

$$\begin{aligned} \mathcal{L}(\mathbf{h}, \boldsymbol{\kappa}) = & \sum_{(i,j) \in A} \int_0^{\sum_{\pi \in \hat{\Pi}} h^\pi \delta_{ij}^\pi} t_{ij}(x) dx + \frac{1}{\theta} \sum_{\pi \in \hat{\Pi}} h^\pi \log h^\pi + \\ & \sum_{(r,s) \in Z^2} \kappa^{rs} \left(d^{rs} - \sum_{\pi \in \hat{\Pi}^{rs}} h^\pi \right) \end{aligned} \quad (9.78)$$

Since there are no non-negativity conditions, it is enough to find the stationary points of the Lagrangian, that is, the points where $\frac{\partial \mathcal{L}}{\partial h} = 0$. Therefore, we require

$$\frac{\partial \mathcal{L}}{\partial h^\pi} = c^\pi + \frac{1}{\theta}(1 + \log h^\pi) - \kappa^{rs} = 0, \quad (9.79)$$

or, solving for h^π ,

$$h^\pi = \exp(\theta \kappa^{rs} - 1) \exp(-\theta c^\pi). \quad (9.80)$$

To find the value of the Lagrange multiplier κ^{rs} , we substitute (9.80) into the demand constraint (9.76), and find that κ^{rs} must be chosen such that

$$\exp(\theta \kappa^{rs} - 1) = \frac{d^{rs}}{\sum_{\pi' \in \Pi^{rs}} \exp(-\theta c^{\pi'})}. \quad (9.81)$$

Substituting into (9.80) gives the logit formula, and therefore the path flows solving Fisk's convex optimization problem are those solving logit stochastic user equilibrium.

This convex optimization problem can be solved using the method of successive averages. In this way, SUE can be solved quite simply, using a familiar method from the basic TAP. The only change is that \mathbf{x}^* is calculated using a method from the previous section, rather than by finding an all-or-nothing assignment loading all flow onto shortest paths. The algorithm is as follows:

1. Choose an initial feasible link assignment \mathbf{x} .
2. Update link travel times based on \mathbf{x} .
3. Calculate target flows \mathbf{x}^* :
 - (a) For each OD pair (r, s) , use a method from the previous section to calculate OD-specific flows \mathbf{x}_{rs}^* .
 - (b) Calculate $\mathbf{x}^* \leftarrow \sum_{(r,s) \in Z^2} \mathbf{x}_{rs}^*$
4. Update $\mathbf{x} \leftarrow \lambda \mathbf{x}^* + (1 - \lambda) \mathbf{x}$ for some $\lambda \in (0, 1]$.
5. If \mathbf{x} and \mathbf{x}^* are sufficiently close, terminate. Otherwise, return to step 2.

(It is often possible to do the calculations in step 3 per origin or per destination, rather than separately for each OD pair.)

Notice that the termination criteria is slightly different than before. With the classical traffic assignment problem, we argued that it was a *bad* idea to compare the current solution to the previous solution and terminate when they are sufficiently small. Here, we are not doing that, but are doing something which looks similar: comparing the current solution \mathbf{x} with the “target” solution \mathbf{x}^* . This works because, unlike in the regular traffic assignment problem, the mapping $\mathbf{H}(\mathbf{C})$ is always continuous and well-defined. This means that \mathbf{x}^* varies slowly with \mathbf{x} : small changes in link flows mean small changes in link and path travel times, which means small changes in path flows from (9.73). With classic traffic assignment, a small change in link flows could mean a shift in the shortest

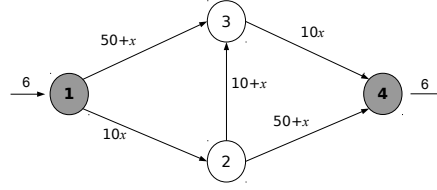


Figure 9.14: Network for demonstrating method of successive averages in stochastic assignment.

path, which would result in a dramatic change in \mathbf{x}^* , moving all flow from an OD pair to that new shortest path. Furthermore, when there are ties there are multiple possible \mathbf{x}^* values. So, in this case we had to introduce auxiliary convergence measures like the relative gap or average excess cost. SUE is simpler in that we can simply compare the current and target solutions — in fact, as defined earlier, the relative gap and average excess cost do not make sense, since the equilibrium principle is no longer defined by all travelers being on the shortest path. For these reasons, the method of successive averages works much better for SUE than for deterministic assignment.

To demonstrate this algorithm, consider the Braess paradox network shown in Figure 9.14, where the demand from node 1 to node 4 is 6 vehicles, $\theta = 0.01$, and all paths are allowable. Assume that we choose the initial solution \mathbf{x} by performing Dial's method on this network, using the *free-flow* travel times. Table 9.3 shows the calculations; in this table, we first calculate the L_i values for nodes (based on shortest paths), then the link likelihoods L_{ij} ; then the node and link weights W_i and W_{ij} , and finally the node and link flows x_i and x_{ij} . This gives initial link flows of $\mathbf{x} = [4.282 \ 1.718 \ 2.563 \ 1.718 \ 4.282]$. Recalculating the link performance functions with these flows gives new link travel times, which we put into Dial's method again, giving the result $\mathbf{x}^* = [3.976 \ 2.024 \ 1.951 \ 2.024 \ 3.976]$. *Notice that \mathbf{x}^* and \mathbf{x} are quite close to each other!* This is a very different situation than when the method of successive averages is applied to the classical traffic assignment problem (compare with the examples in Section 7.2.1), where \mathbf{x}^* was an extreme-point solution quite far away from \mathbf{x} .

So the link flows \mathbf{x} are updated by averaging \mathbf{x}^* into the old \mathbf{x} values, using a weight of $\lambda = 1/2$, producing the results in the rightmost column of Figure 9.14. The process is repeated until convergence.

Disaggregate link-flow formulation and Frank-Wolfe

For the classical TAP, the Frank-Wolfe algorithm was much faster than the method of successive averages, because it chose λ adaptively, to maximize the reduction in the objective function at each iteration. In theory, it is possible to do the same thing with the convex program described above. In practice, it is

Table 9.3: Method of successive averages for stochastic user equilibrium

	Initialization			Iteration 1			
Nodes	L_i	W_i	x_i	L_i	W_i	x_i	
1	0	1	6	0	1	6	
2	0	1	4.282	42.8	1	3.976	
3	10	1.670	4.282	51.7	1.964	3.976	
4	10	2.341	6	94.5	2.964	6	
Links	L_{ij}	W_{ij}	x_{ij}	L_{ij}	W_{ij}	x_{ij}^*	x_{ij}
(1,2)	1	1	4.282	1	1	3.976	4.129
(1,3)	0.670	0.670	1.718	1	1	2.024	1.871
(2,3)	1	1	2.563	0.964	0.964	1.951	2.257
(2,4)	0.670	0.670	1.718	1	1	2.024	1.871
(3,4)	1	1.670	4.282	1	1.964	3.976	4.129

harder because this program makes use of the path-flow variables \mathbf{h} . With the classical traffic assignment problem, we could express solutions and the objective solely in terms of the link flows \mathbf{x} . The addition of the $h \log h$ terms to the SUE objective function renders this impossible. Furthermore, the number of paths grows very quickly with network size (and if we are choosing the reasonable path set to include fully cyclic paths, the number of paths is usually infinite).

The method of successive averages avoids this problem by not actually referring to the objective function at any point — if you review the steps above, you see that the objective function is never calculated. Its role is implicit, guaranteeing that the algorithm will eventually converge, since the direction $\mathbf{x}^* - \mathbf{x}$ is one along which the objective is decreasing. If we can find an efficient way to evaluate the objective function, then we can develop an analogue to the Frank-Wolfe method for classical assignment.

It turns out that we can reformulate the objective function in terms of the destination-aggregated flows on each link, x_{ij}^s (see Section 6.2.3), using the Markov property of the logit loading. There is an equivalent disaggregation by origin (if we reverse our interpretation of the Markov property; see Exercise 18.) Recall the following results from Section 9.3.3:⁹

- There exist values P_{ij}^s for each link (i, j) and destination s , such that

$$h^\pi = d^{rs} \prod_{(i,j) \in \pi} (P_{ij}^s)^{\delta_{ij}^\pi} \quad (9.82)$$

for any path π connecting origin r to destination s .

- The P_{ij}^s values can be interpreted as the conditional probability that a vehicle arriving at node i and destined for node s will have link (i, j) as

⁹Technically, this section only proved them for the case of a single origin-destination pair, but they hold for each destination as well.

the next link in its path. Therefore,

$$P_{ij}^s = \frac{x_{ij}^s}{x_i^s} \quad (9.83)$$

for each destination s , node $i \neq s$, and outgoing link (i, j) .

Using these properties, we derive an equivalent formula for the additional term to the Beckmann function.

Proposition 9.4. *Let \mathbf{h} be a feasible path flow vector satisfying the Markov property, and let x_{ij}^s and x_i^s be the corresponding destination-aggregated link and node flows. Then*

$$\sum_{\pi \in \tilde{\Pi}} h^\pi \log \left(\frac{h^\pi}{d^{rs}} \right) = \sum_{s \in Z} \left(\sum_{(i,j) \in A} x_{ij}^s \log x_{ij}^s - \sum_{\substack{i \in N \\ i \neq s}} x_i^s \log x_i^s \right) \quad (9.84)$$

Proof. We treat each destination s separately; summing over all destinations gives the result. Using (9.82) and (9.83), we have

$$\sum_{\pi \in \tilde{\Pi}^s} h^\pi \left(\frac{h^\pi}{d^{rs}} \right) = \sum_{\pi \in \tilde{\Pi}^s} h^\pi \log \left(\prod_{(i,j) \in A} \left(\frac{x_{ij}^s}{x_i^s} \right)^{\delta_{ij}^\pi} \right) \quad (9.85)$$

$$= \sum_{\pi \in \tilde{\Pi}^s} h^\pi \sum_{(i,j) \in A} \delta_{ij}^\pi \log \frac{x_{ij}^s}{x_i^s} \quad (9.86)$$

$$= \sum_{(i,j) \in A} \left(\sum_{\pi \in \tilde{\Pi}^s} h^\pi \delta_{ij}^\pi \right) \log \frac{x_{ij}^s}{x_i^s} \quad (9.87)$$

$$= \sum_{(i,j) \in A} x_{ij}^s \log \frac{x_{ij}^s}{x_i^s} \quad (9.88)$$

$$= \sum_{(i,j) \in A} x_{ij}^s \log x_{ij}^s - \sum_{(i,j) \in A} x_{ij}^s \log x_i^s \quad (9.89)$$

$$= \sum_{(i,j) \in A} x_{ij}^s \log x_{ij}^s - \sum_{\substack{i \in N \\ i \neq s}} x_i^s \log x_i^s \quad (9.90)$$

by grouping terms in the last sum by the tail node. \square

Furthermore, the left-hand side of equation (9.84) is very similar to the second term in the objective function (9.74). In fact,

$$\sum_{\pi \in \tilde{\Pi}^s} h^\pi \log \left(\frac{h^\pi}{d^{rs}} \right) = \sum_{\pi \in \tilde{\Pi}^s} h^\pi \log h^\pi - \sum_{r \in Z} d^{rs} \log d^{rs} \quad (9.91)$$

using properties of logarithms and the fact that $\sum_{\pi \in \tilde{\Pi}^s} h^\pi = d^{rs}$. They differ only by $\sum_{r \in Z} d^{rs} \log d^{rs}$, which is a constant — and recall from Proposition 4.2

that adding a constant does not affect the optimal solution to an optimization problem.

This means that we can replace the objective function (9.74) with

$$\sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx + \frac{1}{\theta} \sum_{s \in Z} \left(\sum_{(i,j) \in A} x_{ij}^s \log x_{ij}^s - \sum_{\substack{i \in N \\ i \neq s}} x_i^s \log x_i^s \right) \quad (9.92)$$

which does not require path enumeration. We can thus develop the following analogue to the Frank-Wolfe algorithm:

1. Choose an initial, feasible destination-aggregated link assignment \mathbf{x}^s , and the corresponding aggregated link flows \mathbf{x} .
2. Update link travel times based on \mathbf{x} .
3. Calculate target flows:
 - (a) For each destination s , use a method from the previous section to calculate OD-specific flows \mathbf{x}_s^* .
 - (b) Calculate $\mathbf{x}^* \leftarrow \sum_{s \in Z} \mathbf{x}_s^*$
4. Find the value of λ minimizing (9.92) along the line $(1 - \lambda) [\mathbf{x} \quad \mathbf{x}^s] + \lambda [\mathbf{x}^* \quad \mathbf{x}_s^*]$.
5. Update $\mathbf{x} \leftarrow \lambda \mathbf{x}^* + (1 - \lambda) \mathbf{x}$.
6. If \mathbf{x} and \mathbf{x}^* are sufficiently close, terminate. Otherwise, return to step 2.

9.3.6 Relationships between logit loading and most likely path flows (*)

(This optional section draws a connection between the logit-based stochastic equilibrium model and the concept of entropy maximization in most likely path flows.)

The use of the $h \log h$ terms to represent logit assignment in the previous section may have reminded you of the most likely path flows problem in deterministic user equilibrium, discussed in Section 6.2.2. In fact, these two concepts are related to each other quite closely! This optional section explores this relationship, and describes another algorithm for most likely based path flows.

Recall that the most likely path flows were defined as those maximizing

entropy, and solving the following optimization problem:

$$\max_{\mathbf{h}} \quad - \sum_{\pi \in \hat{\Pi}} h_{\pi} \log(h_{\pi}/d) \quad (9.93)$$

$$\text{s.t.} \quad \sum_{\pi \in \Pi} \delta_{ij}^{\pi} h_{\pi} = x_{ij}^* \quad \forall (i, j) \in A \quad (9.94)$$

$$\sum_{\pi \in \hat{\Pi}} h_{\pi} = d \quad (9.95)$$

$$h_{\pi} \geq 0 \quad \forall \pi \in \Pi \quad (9.96)$$

where we have assumed there is a single origin-destination pair for simplicity.

We can transform this into an equivalent optimization that more closely resembles the stochastic user equilibrium optimization problems. First, we replace maximization by minimization by changing the sign of the objective. Second, we can remove d from the objective, because

$$\sum_{\pi \in \hat{\Pi}} h_{\pi} \log(h_{\pi}/d) = \sum_{\pi \in \hat{\Pi}} h_{\pi} \log h_{\pi} - \sum_{\pi \in \hat{\Pi}} h_{\pi} \log d \quad (9.97)$$

$$= \sum_{\pi \in \hat{\Pi}} h_{\pi} \log h_{\pi} - d \log d \quad (9.98)$$

and the constant $d \log d$ can be ignored. This gives

$$\min_{\mathbf{h}} \quad \sum_{\pi \in \hat{\Pi}} h_{\pi} \log h_{\pi} \quad (9.99)$$

$$\text{s.t.} \quad \sum_{\pi \in \Pi} \delta_{ij}^{\pi} h_{\pi} = x_{ij}^* \quad \forall (i, j) \in A \quad (9.100)$$

$$\sum_{\pi \in \hat{\Pi}} h_{\pi} = d \quad (9.101)$$

$$h_{\pi} \geq 0 \quad \forall \pi \in \Pi \quad (9.102)$$

We make one more change: rather than insisting that the path flows match the equilibrium link flows exactly, we change constraint (9.100) to require that the average travel time across all travelers be c^* . That is, we change the constraint to

$$\sum_{\pi \in \hat{\Pi}} h_{\pi} c^{\pi} = c^* d. \quad (9.103)$$

Arguing as in the previous section, we can ignore the non-negativity constraint (9.102), since it will not be binding at optimality. We then Lagrangianize the other two constraints, introducing multipliers θ and κ for (9.103) and (9.101), respectively. The resulting Lagrangian function is

$$\mathcal{L}(\mathbf{h}, \theta, \kappa) = \sum_{\pi \in \hat{\Pi}} h^{\pi} \log h^{\pi} + \theta \left(\sum_{\pi \in \hat{\Pi}} h^{\pi} c^{\pi} - c^* d \right) + \kappa \left(d - \sum_{\pi \in \hat{\Pi}} h^{\pi} \right) \quad (9.104)$$

without non-negativity constraints. The optimality conditions are thus

$$\frac{\partial \mathcal{L}}{\partial h^\pi} = 1 + \log h^\pi + \theta c^\pi - \kappa = 0 \quad \forall \pi \in \hat{\Pi} \quad (9.105)$$

$$\frac{\partial \mathcal{L}}{\partial \theta} = \sum_{\pi \in \hat{\Pi}} h^\pi c^\pi - c^* d = 0 \quad (9.106)$$

$$\frac{\partial \mathcal{L}}{\partial \kappa} = \sum_{\pi \in \hat{\Pi}} h^\pi - d = 0 \quad (9.107)$$

$$(9.108)$$

The last two of these are simply the constraints (9.100) and (9.101). The first can be solved for the path flows, giving

$$h^\pi = \exp(\kappa - 1 - \theta c^\pi). \quad (9.109)$$

To satisfy the demand constraint (9.101), κ must be chosen so that the sum of (9.109) over all paths gives the total demand d . Omitting the algebra, we must have

$$\kappa = \log d(1 - \log \sum_{\pi' \in \Pi} \exp(-\theta c^{\pi'})) \quad (9.110)$$

or

$$h^\pi = d \frac{\exp(-\theta c^\pi)}{\sum_{\pi' \in \Pi} \exp(-\theta c^{\pi'})}. \quad (9.111)$$

But this is just the logit formula! The Lagrange multiplier θ must be chosen to satisfy the remaining constraint on the average path cost.

This derivation shows that the most likely path flows and stochastic user equilibrium problems have a similar underlying structure. If we relax the requirement that all travelers be on shortest paths, and simply constrain the average cost of travel, the most likely path flows coincide with a logit loading, where the parameter θ is the Lagrange multiplier for this constraint. As θ approaches infinity, the average cost of travel approaches its value at the deterministic user equilibrium solution, and the stochastic user equilibrium path flows approach the most likely path flows in the corresponding deterministic problem. This provides another algorithmic approach for solving for most likely path flows, in addition to those discussed in Section 7.5. In practice, this algorithm is difficult to implement, because of numerical issues that arise as θ grows large.

9.3.7 Alternatives to logit loading

The majority of this section has focused on the logit model for stochastic route choice, because it demonstrates the main ideas simply, and leads to computationally-efficient solution techniques. There are several serious criticisms of logit assignment. For instance, the assumption that the error terms ϵ^π are independent across paths is hard to defend if paths overlap significantly.

Common methods for creating totally acyclic route sets can also create unreasonable artifacts, and allowing all cyclic paths can also be unreasonable if the network topology creates many such paths with low travel times. See Exercise 25 for some concrete examples. The main alternative to the logit model is the *probit* model, in which the ϵ terms have a multivariate normal distribution, with a (possibly nondiagonal) covariance matrix to allow for correlation between these terms.

The framework of stochastic user equilibrium can be generalized to other distributions of the error terms, including correlation. The full development of this general framework is beyond the scope of this book, but we provide an overview and summary. The objective function (9.74) must be replaced with

$$z(\mathbf{x}) = - \sum_{(r,s) \in Z^2} d^{rs} E \left[\min_{\pi \in \Pi^{rs}} (c^\pi + \epsilon) \right] + \sum_{(i,j) \in A} x_{ij} t_{ij} - \sum_{(i,j) \in A} \int_0^{x_{ij}} t_{ij}(x) dx \quad (9.112)$$

where the expectation is taken with respect to the “unobserved” random variables ϵ , and t_{ij} in the last two terms are understood to be functions of x_{ij} . It can be shown that this function is convex, and therefore that the SUE solution is unique. However, evaluating this function is harder. The first term in (9.112) involves an expectation over all paths connecting an OD pair. In discrete choice, this is known as the *satisfaction function*, and expresses the expected perceived travel time on a path chosen by a traveler. In the case of the logit model, this expectation can be computed in closed form; for most distributions it cannot, and must be evaluated through Monte Carlo sampling or another approximation.

The method of successive averages can still be used, even without evaluating the objective function. Step 4a needs to be replaced with a stochastic network loading, using the current travel times and whatever distribution of ϵ is chosen. This often requires Monte Carlo sampling as well: for (multiple) samples of the ϵ terms, the shortest paths can be found using one of the standard algorithms, and the resulting flows averaged together to form an estimate of \mathbf{x}^* .

Because we are using an estimate of \mathbf{x}^* , it is possible that the target direction is not exactly right, and that it is not in a direction in which $z(\mathbf{x})$ is decreasing. Nevertheless, as long as it is correct “on average” (i.e., the sampling is done in an unbiased manner), one can show that the method of successive averages will still converge to the stochastic user equilibrium solution.

9.4 Exercises

1. [25] Verify that each of the demand functions D below is strictly decreasing and bounded above (for $\kappa \geq 0$), then find the inverse functions $D^{-1}(d)$.
 - (a) $D(\kappa) = 50 - \kappa$
 - (b) $D(\kappa) = 1000/(\kappa + 1)$
 - (c) $D(\kappa) = 50 - \kappa^2$

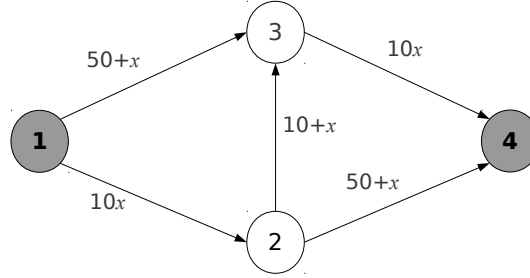


Figure 9.15: Network for Exercise 5.

- (d) $D(\kappa) = (\kappa + 2)/(\kappa + 1)$
2. [13] The total misplaced flow reflects consistency of a solution with (in this case, that the OD matrix should be given by the demand functions). Suggest another measure for how close a particular OD matrix and traffic assignment are to satisfying this consistency condition. Your proposed measure should be a nonnegative and continuous function of values related to the solution (e.g., d^{rs} , x_{ij} , κ^{rs} , etc.), which is zero if and only if the OD matrix is completely consistent with the demand functions. Compare your new measure with the total misplaced flow, and comment on any notable differences.
 3. [31] Verify that the elastic demand objective function (9.16) is convex, given the assumptions made on the demand functions.
 4. [47] Using the network in Figure 9.15, solve the elastic demand equilibrium problem with demand given by $d^{14} = 15 - \kappa^{14}/30$. Perform three iterations of the Frank-Wolfe method, and report the average excess cost and total misplaced flow for the solution.
 5. [48] Using the network in Figure 9.15, solve the elastic demand equilibrium problem with demand given by $d^{14} = 10 - \kappa^{14}/30$.
 - (a) Perform three iterations of the Frank-Wolfe method designed for elastic demand (Section 9.1.5).
 - (b) Transform the problem to an equivalent fixed-demand problem using the Gartner transformation from Section 9.1.2, and perform three iterations of the original Frank-Wolfe method.
 - (c) Compare the performance of these two methods: after three iterations, which is closer to satisfying the equilibrium and demand conditions?
 6. [49] Using the network in Figure 9.16, solve the elastic demand equilibrium problem with demand functions $q^{19} = 1000 - 50(u^{19} - 50)$ and $q^{49} = 1000 - 75(u^{49} - 50)$. The cost function on the light links is $3 + (x_a/200)^2$,

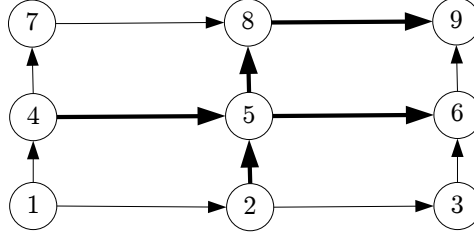


Figure 9.16: Network for Exercise 6.

and the cost function on the thick links is $5 + (x_a/100)^2$. 1000 vehicles are traveling from node 1 to 9, and 1000 vehicles from node 4 to node 9. Perform three iterations of the Frank-Wolfe method and report the link flows and OD matrix.

7. [13] Assume that (\mathbf{x}, \mathbf{d}) and $(\mathbf{x}^*, \mathbf{d}^*)$ are both feasible solutions to an elastic demand equilibrium problem. Show that $(\lambda\mathbf{x} + (1 - \lambda)\mathbf{x}^*, \lambda\mathbf{d} + (1 - \lambda)\mathbf{d}^*)$ is also feasible if $\lambda \in [0, 1]$. This ensures that the Frank-Wolfe solutions are always feasible, assuming we start with \mathbf{x} and \mathbf{d} values which are consistent with each other, and always choose targets in a consistent way.
8. [57] (Calculating derivative formulas.) Let \mathbf{x} and \mathbf{d} be the current, feasible, link flows and OD matrix, and let \mathbf{x}^* and \mathbf{d}^* be any other feasible link flows and OD matrix. Let $\mathbf{x}' = \lambda\mathbf{x}^* + (1 - \lambda)\mathbf{x}$ and $\mathbf{d}' = \lambda\mathbf{d}^* + (1 - \lambda)\mathbf{d}$.
 - (a) Let $f(\mathbf{x}', \mathbf{d}')$ be the objective function for the elastic demand equilibrium problem. Recognizing that \mathbf{x}' and \mathbf{d}' are functions of λ , calculate $df/d\lambda$.
 - (b) For the elastic demand problem, show that $\frac{df}{d\lambda}|_{\lambda=0} \leq 0$ if \mathbf{d}^* and \mathbf{x}^* are chosen in the way given in the text. That is, the objective function is nonincreasing in the direction of the “target.” (You can assume that the demand function values are strictly positive if that would simplify your proof.)
9. [24] Consider a two-link network. For each pair of link performance functions shown below, determine whether or not the symmetry condition (9.26) is satisfied.
 - (a) $t_1 = 4 + 3x_1 + x_2, t_2 = 2 + x_1 + 4x_2$
 - (b) $t_1 = 7 + 3x_1 + 4x_2, t_2 = 12 + 2x_1 + 4x_2$
 - (c) $t_1 = 4 + x_1 + 3x_2, t_2 = 2 + 3x_1 + 2x_2$
 - (d) $t_1 = 3x_1^2 + x_2, t_2 = 4 + x_1 + 4x_2^3$
 - (e) $t_1 = 3x_1^2 + 2x_2^2, t_2 = 2x_1^2 + 3x_2^3$

(f) $t_1 = 50 + x_1$, $t_2 = 10x_2$

10. [34] Determine which of the pairs of link performance functions in the previous exercise are strictly monotone.
11. [49] Consider the network in Figure 9.16 with a fixed demand of 1000 vehicles from 1 to 9 and 1000 vehicles from 4 to 9. The link performance function on every link a arriving at a “merge node” (that is, nodes 5, 6, 8, and 9) is $4 + (x_a/150)^2 + (x_{a'}/300)^2$ where a' is the other link arriving at the merge. Verify that the link interactions are symmetric, and perform five iterations of the Frank-Wolfe method. Report the link flows and travel times.
12. [49] Consider the network in Figure 9.16 with a fixed demand of 1000 vehicles from 1 to 9 and 1000 vehicles from 4 to 9. The link performance function on every link a arriving at a “merge node” (that is, nodes 5, 6, 8, and 9) is $3 + (x_a/200)^2 + (x_{a'}/400)^2$ if the link is light, and $5 + (x_a/100)^2 + (x_{a'}/200)^2$ if the link is thick, where a' is the other link arriving at the merge. Show that not all link interactions are symmetric, and then perform five iterations of the diagonalization method and report the link flows, travel times, and relative gap.
13. [49] See the network in Figure 9.17, where the numbers indicate the label for each link. The link performance functions are:

$$\begin{aligned} t_1(\mathbf{x}) &= 1 + 4x_1 + 2x_2 \\ t_2(\mathbf{x}) &= 2 + x_1 + 2x_2 \\ t_3(\mathbf{x}) &= 3 + 2x_3 + x_4 \\ t_4(\mathbf{x}) &= 2 + 3x_4 \end{aligned}$$

and the demand from node A to node C is 10 vehicles. For both methods below, start with an initial solution loading all flow on links 1 and 4.

- (a) Use three iterations of the diagonalization method to try to find an equilibrium solution, and report the link flows and average excess cost. (As before, this means finding three \mathbf{x}^* vectors *after* your initial solution.)
- (b) Use three iterations of simplicial decomposition to try to find an equilibrium solution, and report the link flows and average excess cost. Three iterations means that \mathcal{X} should have three vectors in it when the algorithm terminates (unless it terminates early because the \mathbf{x}^* you find is already in \mathcal{X}). For each subproblem, make the number of improvement steps one less than the size of \mathcal{X} (so when \mathcal{X} has 1 vector, perform 0 steps; when it has 2 vectors, perform 1 step, and so on). For each of these steps, try the sequence of μ values $1/2, 1/4, 1/8, \dots$, choosing the first that reduces the restricted average excess cost.

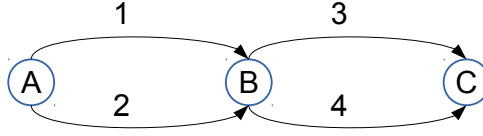


Figure 9.17: Network for Exercise 13.

14. [64] Prove Proposition 9.2.
15. [42] Show that the three path generation methods described in Section 9.3.2 indeed yield totally acyclic path sets, and that they always include the shortest path from the origin to the destination.
16. [32] Show that any set of totally acyclic paths satisfies the segment substitution property. Then show that the set of all paths (including all cyclic paths) also satisfies this property.
17. [53] Complete the example following equation (9.40) by forming the infinite sum for the flow on links (2,3) and (3,2), and showing that they are equal to $1/2$.
18. [62] Reformulate the Markov property, and the results in Section 9.3.3 in terms of conditional probabilities for the link a vehicle used to *arrive* at a given node, rather than the link a vehicle will choose to depart a given node.
19. [34] In the logit formula (9.39), show that the same path choice probabilities are obtained if each link travel time t_{ij} is replaced with $L_i + t_{ij} - L_j$, where \mathbf{L} is a vector of node-specific constants. (In practice, these are usually the shortest path distances from the origin.)
20. [51] What would have to change in the “full cyclic path set” stochastic network loading procedure, if there were multiple links with the same tail and head nodes?
21. [35] Consider the network in Figure 9.18 with 1000 vehicles traveling from node A to node I, where the link labels are the travel times. Use the first criterion in Section 9.3.2 to define the path set, and identify the link flows corresponding to these travel times. Assume $\theta = 1$.
22. [35] Repeat Exercise 21 for the third definition of a path set. Assume $\theta = 1$.
23. [34] Show that the objective function (9.74) is strictly convex.
24. [49] Consider the network in Figure 9.16 with a fixed demand of 1000 vehicles from 1 to 9 and 1000 vehicles from 4 to 9. The light links have

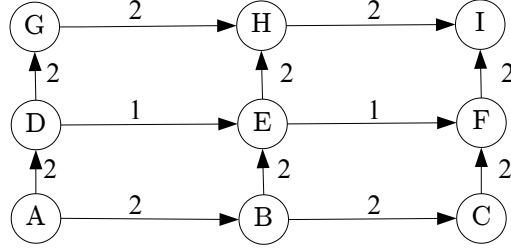
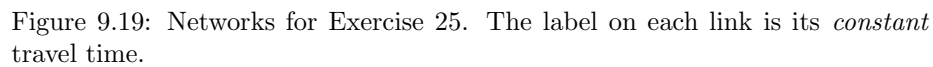


Figure 9.18: Network for Exercises 21 and 22.

delay function $3 + (x_{ij}/200)^2$, and the dark links have delay function $5 + (x_{ij}/100)^2$. Assume that drivers choose paths according to the stochastic user equilibrium principle, with $\theta = 1/5$. Perform three iterations of the method of successive averages, and report the link flows. Assume all paths are allowed.

25. [24] (Limitations of logit assignment). This problem showcases three “problem instances” for the stochastic network loading models described in this chapter (the logit model, and one proposed definition of allowable paths). Throughout, assume that the first definition in Section 9.3.2 is used to define the allowable path set. These instances motivated the development of probit and other, more sophisticated, stochastic equilibrium models.
 - (a) Consider the network in Figure 9.19(a), where the numbers by each link represents the travel time and $z \in (0, 1)$. Let p_{\uparrow} represent the proportion of vehicles choosing the top path. In a typical probit model, we have $p_{PROBIT}^{\uparrow} = 1 - \Phi\left(\sqrt{\frac{z}{2\pi - z}}\right)$, where $\Phi(\cdot)$ is the standard cumulative normal distribution function. Calculate p_{LOGIT}^{\uparrow} for the *logit* model as a function of z , and plot p_{PROBIT}^{\uparrow} and p_{LOGIT}^{\uparrow} as z varies from 0 to 1. Which do you think is more realistic, and why?
 - (b) Consider networks of the type shown in Figure 9.19(b), where there is a top path consisting of a single link, and a number of bottom paths. The network is defined by an integer m ; there are $m - 1$ intermediate nodes in the bottom paths, and each consecutive pair of intermediate nodes is connected by two parallel links with travel time $8/m$. In a common probit model, $p_{PROBIT}^{\uparrow} = \Phi(-0.435\sqrt{m})$. What is p_{LOGIT}^{\uparrow} as a function of m ? Again create plots for small values of m (say 1 to 8), indicate which you think is more realistic, and explain why.
 - (c) In the networks in Figure 9.19(c), identify the proportion of travelers choosing each link if $\theta = 1$. The left and right panels show a network before and after construction of a new link. Again identify the



proportion of travelers choosing each link if $\theta = 1$. Do your findings seem reasonable?

Part III

**Dynamic Traffic
Assignment**

Chapter 10

Network Loading

This chapter discusses *network loading*, the process of modeling the state of traffic on a network, given the route and departure time of every vehicle. In static traffic assignment, this is a straightforward process based on evaluating link performance functions. In dynamic traffic assignment, however, network loading becomes much more complicated due to the additional detail in the traffic flow models — but it is exactly this complexity which makes dynamic traffic assignment more realistic than static traffic assignment. Rather than link performance functions, dynamic network loading models generally rely on some concepts of traffic flow theory. There are a great many theories, and an equally great number of dynamic network loading models, so this chapter will focus on those most commonly used.

To give the general flavor of network loading, we start with two simple *link models*, the point queue and spatial queue (Section 10.1), which describe traffic flow on a single link. We next present three simple *node models* describing how traffic streams behave at junctions (Section 10.2). With these building blocks we can perform network loading on simple networks, showing how link and node models interact to represent traffic flow in a modular way (Section 10.3).

However, the point queue and spatial queue models have significant limitations in representing traffic flow. The most common network loading models for dynamic traffic assignment are based on the hydrodynamic model of traffic flow, reviewed in Section 10.4. This theory is based in fluid mechanics, and assumes that the traffic stream can be modeled as the motion of a fluid, but it can be derived from certain car-following models as well, which have a more behavioral basis. The cell transmission model and link transmission model are link models based on this theory, and both of these are discussed in Section 10.5. Section 10.6 concludes the chapter with a discussion of more sophisticated node models that can represent general intersections.

This chapter aims to present several alternative network loading schemes as part of the general dynamic traffic assignment framework in Figure 1.8, so that any of them can be combined with the other steps in a flexible way.

10.1 Link Model Concepts

Link models are efficient ways to represent congestion and traffic flow. Efficiency is key, since link models must be implemented on each link in a large network over the entire analysis period; and furthermore, due to the iterative nature of solving for equilibrium, the network loading must be repeated many times with different path flow values as input. This section presents the basic concepts of link models, and demonstrates them in the point queue and spatial queue models. More sophisticated link models will be presented later in this chapter, in Section 10.5.

So, in this section, we are solely concerned with a single link. Call the length of this link L , so $x = 0$ corresponds to the upstream end of the link, and $x = L$ corresponds to the downstream end. All of the link models we discuss in this book operate in *discrete time*. That is, we divide the analysis period into small time intervals of length Δt , and assume uniform conditions within each time interval. These time intervals are very small relative to the analysis period; at a minimum, no vehicle should be able to traverse more than one link in a single time interval, so Δt can be no greater than the shortest free-flow travel time on a link. In practice, Δt values on the order of 5–10 seconds are a reasonable choice. For convenience in this section, we assume that the unit of time is chosen such that $\Delta t = 1$, so we will only track the state of the network for times $t \in \{0, 1, 2, \dots, \bar{T}\}$, where \bar{T} is the length of the analysis period (in units of Δt).

In any discrete time model, it is important to clarify what exactly we mean when we index a variable with a time interval, such as $N(t)$ or $y(t)$. Does this refer to the value of x at the start of the t -th time interval, the end of the interval, the average of a continuous value through the interval, or something else? This chapter will consistently use the following convention: when referring to a quantity *measured at a single instant in time*, such as the number of vehicles on a link, or the instantaneous speed of a vehicle, we will take such measurements at the *start* of a time interval. When referring to a quantity *measured over time*, such as the number of vehicles passing a fixed point on the link, we take such measurements *over* the time interval. (Figure 10.1) This distinction can also be expressed by referring, say, to the number of vehicles on a link *at* time t (meaning the start of the t -th interval), or the number of vehicles that exit a link *during* time t .

We load the network in increasing order of time. That is, we start with the network state at $t = 0$ (usually assuming an empty condition). Then, with these values known, we compute the network state at $t = 1$, then at $t = 2$, and so forth. That is, for any point in time, any values at an earlier point in time can be treated as “known” values.

The definition of the “network state” at a timestep is intentionally left a bit vague at this point. It contains all information necessary for modeling traffic flow (e.g., the locations of vehicles, traffic signal indications). The collection of network states at all time steps must also have enough information to check consistency with the equilibrium principle (cf. Section 1.3) after the network loading is complete, and for adjusting vehicle flows if equilibrium is not satisfied.

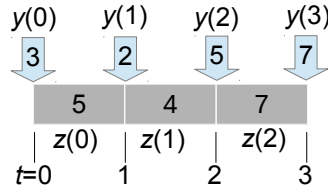


Figure 10.1: Convention for indexing discrete time intervals for two hypothetical variables y (measured at an instant) and z (measured over time).

The exact components of the network state vary from one model to the next (for instance, signals may be modeled in detail, approximately, or not at all), and as you read about the link and node models presented in this chapter, you should think about what information you would need to store in order to perform the computations for each link and node model. At a minimum, it is common to record the cumulative number of vehicles which have entered and left each link at each timestep, since the start of the modeling period. These values are denoted by $N^\uparrow(t)$ and $N^\downarrow(t)$, respectively; the arrows are meant as a mnemonic for “upstream” and “downstream”, since they can be thought of as counters at the ends of the links. So, for instance, $N^\uparrow(3)$ is the number of vehicles which have crossed the upstream end of the link by the start of the 3rd time interval (the cumulative entries), and $N^\downarrow(5)$ is the number of vehicles which have crossed the downstream end of the link by the start of the 5th time interval (cumulative exits). We will assume that the network is empty when $t = 0$ (no vehicles anywhere), and as a result $N^\uparrow(t) - N^\downarrow(t)$ gives the number of vehicles currently on the link at any time t .

It is possible to use different time interval lengths for different links and nodes, and this can potentially reduce the necessary computation time. There are also *continuous time* dynamic network loading models, where flows and other traffic variables are assumed to be functions defined for any real time value, not just a finite number of points. These are not discussed here to keep the focus on the basic network loading ideas, and to avoid technical details associated with infinitesimal calculations. Finally, some formulas may call for the value of a discrete variable at a non-integer point in time, such as $N^\uparrow(3.4)$, in which case a linear interpolation can be used between the neighboring values $N^\uparrow(3)$ and $N^\uparrow(4)$. If possible, the time step should be chosen to minimize or eliminate these interpolation steps, which are time-consuming and which can introduce numerical errors.

10.1.1 Sending and receiving flow

The main outputs of a link model are the *sending flow* and *receiving flow*, calculated for each discrete time interval. The sending flow at time t , denoted $S(t)$, is the number of vehicles which would leave the link during the t -th time

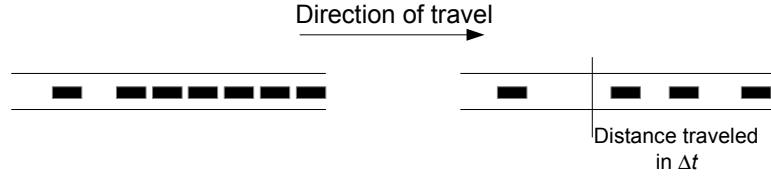


Figure 10.2: Calculating the sending flow when there is a queue on the link (left) and when the link is at free-flow (right).

interval (that is, between times t and $t + 1$) if there was no obstruction from downstream links or nodes (you can imagine that the link is connected to a wide, empty link downstream). You can also think of this as the flow that is ready to leave the link during this time interval. The sending flow is calculated at the downstream end of a link.

To visualize sending flow, Figure 10.2 shows two examples of links. In the left panel, there is a queue at the downstream end of the link. If there is no restriction from downstream, the queue would discharge at the full capacity of the link, and the sending flow would be equal to the capacity of the link multiplied by Δt . In the right panel, the link is uncongested and vehicles are traveling at free-flow speed. The sending flow will be less than the capacity, because relatively few vehicles are close enough to the downstream end of the link to exit within the next time interval. The vertical line in the figure indicates the distance a vehicle would travel at free-flow speed during one time step, so in this case the sending flow would be 3 vehicles. Note that the *actual* number of vehicles which can leave the link during the next time step may be affected by downstream conditions: perhaps the next link is congested, or perhaps there is a traffic signal at the downstream end. These considerations are irrelevant for calculating sending flow, and will be treated with node models, introduced in Section 10.2. In any case, the sending flow is an upper bound on the actual number of vehicles which can depart the link.

The receiving flow during time t , denoted $R(t)$, is the number of vehicles which would enter the link during the t -th time interval if the upstream link could supply a very large (even infinite) number of vehicles: you can imagine that the upstream link is wide, and completely full at jam density. You can also think of this as the maximum amount of flow which can enter the link during this interval, taking into account the available free space. The receiving flow is calculated at the upstream end of the link.

To visualize receiving flow, Figure 10.3 shows two examples of links. In the left panel, the upstream end of the link is empty. This means that vehicles can potentially enter the link at its full capacity, and the receiving flow would equal the link's capacity multiplied by Δt . The actual number of vehicles which will enter the link may be less than this, if there is little demand from upstream — like the sending flow, the receiving flow is simply an upper bound indicating how

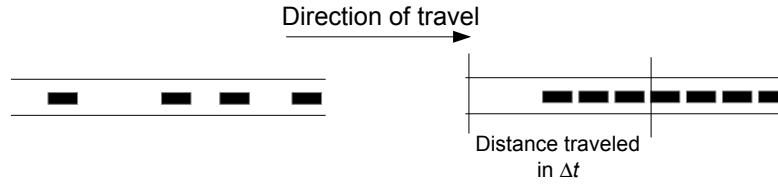


Figure 10.3: Calculating the receiving flow when the link is at free-flow (left) and when there is a queue on the link (right).

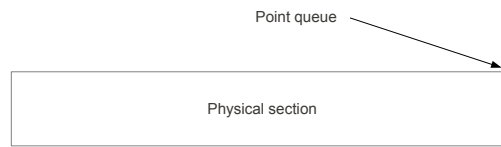


Figure 10.4: Division of a link in point queue models.

many vehicles could potentially enter the link if demand were high enough. In the right panel, there is a stopped queue which nearly fills the entire link. Here the vertical line indicates how far into the link a vehicle would travel at free-flow speed during one time step. Assuming that the stopped vehicles remain stopped throughout the t -th time interval, the receiving flow is the number of vehicles which can physically fit into the link, in this case 2.

Each link model has a slightly different way of calculating the sending and receiving flows, which correspond to different assumptions on traffic behavior with the link, or to different calculation methods. The next two subsections present simple link models. Notice how the different traffic flow assumptions in these models lead to different formulas for sending and receiving flow.

10.1.2 Point queue

Point queue models divide each link into two sections:

- A **physical section** which spans the length of the link and is assumed uncongestible: vehicles will always travel over this section at free-flow speed.
- A **point queue** at the downstream end of the link which occupies no physical space, but conceptually holds vehicles back to represent any congestion delay on the link.

These are shown in Figure 10.4.

Point queues can be thought of in several ways. One can imagine a wide link which necks down at its downstream end — the physical section reflects the

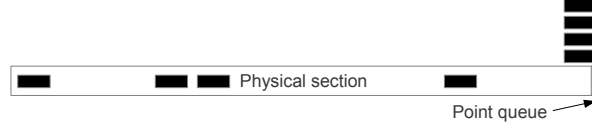


Figure 10.5: Envisioning a point queue as vehicles stacking vertically.

wide portion of the link, and the point queue represents the vehicles which are delayed as the capacity is reduced at the downstream bottleneck. One can also imagine a traffic signal at the downstream end, and magical technology (flying cars?) which allows vehicles to “stack” vertically at the signal — there can be no congestion upstream of this “stack,” since vehicles can always fly to the top. (Figure 10.5) One may even imagine that there is no physical meaning to either of these, and that the physical section and point queue merely represent the delays incurred from traveling the link at free-flow, and the additional travel time due to congestion.

The point queue discharges vehicles at a maximum rate of q_{max}^\downarrow (measured in vehicles per unit time), called the *capacity*. The capacity imposes an upper limit on the sending flow, so we always have

$$S(t) \leq q_{max}^\downarrow \Delta t. \quad (10.1)$$

However, if the queue is empty, or if only a few vehicles are in the queue, the discharge rate may be less than this. Once the queue empties, the only vehicles which can exit the link are ones reaching the downstream end from the uncongested physical section. Since we assume that all vehicles in this section travel at the free-flow speed (which we will denote u_f), this means that only the vehicles that are closer than $u_f \Delta t$ to the downstream end can possibly leave.

We can use the cumulative counts N^\uparrow and N^\downarrow to count the number of vehicles which are close enough to the downstream end to exit in the next time step. Since the entire physical section is traversed at the free-flow speed u_f , a vehicle whose distance from the downstream end of the link is exactly $u_f \Delta t$ distance units must have passed the upstream end of the link exactly $(L - u_f \Delta t)/u_f$ time units ago. We call this a “threshold” vehicle, since any vehicle entering the link after this one has not yet traveled far enough, while any vehicle entering the link before this one is close enough to the downstream end to exit. The number of vehicles between the threshold vehicle and the downstream end of the link can thus be given by

$$N^\uparrow \left(t - \frac{L - u_f \Delta t}{u_f} \right) - N^\downarrow(t) = N^\uparrow \left(t + \Delta t - \frac{L}{u_f} \right) - N^\downarrow(t). \quad (10.2)$$

The sending flow is the smaller of the number of vehicles which are close enough to the downstream end to exit, given by equation (10.2), and the capacity of the queue. Thus

$$S(t) = \min \{ N^\uparrow(t + \Delta t - L/u_f) - N^\downarrow(t), q_{max}^\downarrow \Delta t \}. \quad (10.3)$$

The receiving flow for the point queue model is easy to calculate. Since the physical section is uncongestible, the link capacity is the only limitation on the rate at which vehicles can enter. The capacity of the upstream end of the link may be different than the capacity of the downstream end of the link (perhaps due to a lane drop, or a stop sign at the end of the link), so we denote the capacity of the upstream end by q_{max}^\uparrow . (If we just use q_{max} without an \uparrow or \downarrow superscript, the same capacity applies over the whole link.) The receiving flow is given by

$$R(t) = q_{max}^\uparrow \Delta t. \quad (10.4)$$

In real traffic networks, queues occupy physical space and cannot be confined to a single point. The spatial queue model in the next subsection shows one way to reflect this.

Table 10.1 shows how the point queue model operates, depicting the state of a link over ten time steps. The N^\uparrow and N^\downarrow columns express the number of vehicles which have entered and left the link at each time step, as well as the sending and receiving flows during each timestep. The difference between N^\uparrow and N^\downarrow represents the number of vehicles on the link at any point in time. In this example, we assume that the free-flow speed is $u_f = L/(3\Delta t)$ (so a vehicle takes 3 time steps to traverse the link under free-flow conditions), the upstream capacity is $q_{max}^\uparrow = 10/\Delta t$, and the downstream capacity is $q_{max}^\downarrow = 5/\Delta t$. Initially, the sending flow is zero, because no vehicles have reached the downstream end of the link. The sending flow then increases as flow exits, but eventually reaches the downstream capacity. At this point, a queue forms and vehicles exit at the downstream capacity rate. Eventually, the queue clears, the link is empty, and the sending flow returns to zero. The receiving flow never changes from the upstream capacity, even when a queue is present. In this example, notice that $N^\downarrow(t+1) = N^\downarrow(t) + S(t)$. This happens because we are temporarily ignoring what might be happening from downstream. Depending on downstream congestion, $N^\downarrow(t+1)$ could be less than $N^\downarrow(t) + S(t)$; but it could never be greater, because the sending flow is always a limit on the number of vehicles that can exit. Also notice that for all time steps, $N^\uparrow(t+1) \leq N^\uparrow(t) + R(t)$, because the receiving flow is a limit on the number of vehicles that can enter the link.

In practice, the upstream and downstream capacities are usually assumed the same, in which case we just use the notation q_{max} to refer to capacities at both ends. The network features which would make the capacities different upstream and downstream (such as stop signs or signals) are usually better represented with node models, discussed in the next section.

10.1.3 Spatial queue

The spatial queue model is similar to the point queue model, except that a maximum queue length is now enforced. When the queue reaches this maximum length, no further vehicles are allowed to enter the link. That is, the queue now occupies physical space, and because the link is finite in length, the link can

Table 10.1: Point queue example, with $u_f = L/(3\Delta t)$, $q_{max}^\uparrow = 10/\Delta t$, and $q_{max}^\downarrow = 5/\Delta t$.

t	N^\uparrow	N^\downarrow	R	S
0	0	0	10	0
1	1	0	10	0
2	5	0	10	0
3	10	0	10	1
4	17	1	10	4
5	27	5	10	5
6	30	10	10	5
7	30	15	10	5
8	30	20	10	5
9	30	25	10	5
10	30	30	10	0

become completely blocked. This will result in *queue spillback*, as vehicles on upstream links will be unable to enter the blocked link. Like the point queue model, we assume that the queue is always at the downstream end of the link: there is at most one uncongested physical section at the upstream end of the link, and at most one stopped queue at the downstream end, in that order. This is still a simplification of real traffic (where links can have multiple congested and uncongested sections), but by allowing the length of the physical section to shrink as the queue grows, one can model the queue spillback phenomenon which is common in congested networks.

The sending flow for the spatial queue model is calculated in exactly the same way as for the point queue model: the smaller of the number of vehicles close enough to the downstream end of the link to exit in the next time step, and the capacity of the link:

$$S(t) = \min\{N^\uparrow(t + \Delta t - L/u_f) - N^\downarrow(t), q_{max}^\downarrow \Delta t\}. \quad (10.5)$$

The receiving flow includes an additional term to reflect the finite space on the link for the queue, alongside the link capacity. Since vehicles in the queue are stopped, the space they occupy is given by the jam density k_j , expressed in vehicles per unit length. The maximum number of vehicles the link can hold is $k_j L$, while the number of vehicles currently on the link is $N^\uparrow(t) - N^\downarrow(t)$. The receiving flow cannot exceed the difference between these:

$$R(t) = \min\{k_j L - (N^\uparrow(t) - N^\downarrow(t)), q_{max}^\uparrow \Delta t\}. \quad (10.6)$$

By assuming that the queue is always at the downstream end of the link, the spatial queue model essentially assumes that all vehicles in a queue move together. In reality, there is some delay between when the head of the queue starts moving, and when the vehicle at the tail of the queue starts moving —

Table 10.2: Spatial queue example, with $u_f = L/(3\Delta t)$, $q_{max}^\uparrow = 10/\Delta t$, $q_{max}^\downarrow = 5/\Delta t$, and $k_j L = 20$.

t	N^\uparrow	N^\downarrow	R	S
0	0	0	10	0
1	1	0	10	0
2	5	0	10	0
3	10	0	10	1
4	17	1	4	4
5	21	5	4	5
6	25	10	5	5
7	30	15	5	5
8	30	20	10	5
9	30	25	10	5
10	30	30	10	0

when a traffic light turns green, vehicles start moving one at a time, with a slight delay between when a vehicle starts moving and when the vehicle behind it starts moving. These delays cannot be captured in a spatial queue model. To represent this behavior, we will need a better understanding of traffic flow theory. Section 10.4 will present this information, and we will ultimately build more realistic link models. The point queue and spatial queue models nevertheless illustrate the basic principles of link models, and what the sending and receiving flow represent.

Table 10.2 shows how the spatial queue model operates. It is similar to the point queue example (Table 10.1), but we now introduce a jam density, assuming that the maximum number of vehicles which can fit on the link is $k_j L = 20$. In this example, the receiving flow drops once the queue reaches a certain size. This reflects the finite space available on the link. As a result, it takes longer for the 30 vehicles to enter the link. The queue is still able to completely discharge by the end of the ten time steps. As before, the difference between $N^\downarrow(t)$ and $N^\downarrow(t+1)$ is never more than $S(t)$ (in this example, exactly equal because we are ignoring downstream conditions), and the difference between $N^\uparrow(t)$ and $N^\uparrow(t+1)$ is never more than $R(t)$.

10.2 Node Model Concepts

Node models complement the link models discussed in the previous section, by representing how flows between different links interact with each other. By focusing on sending and receiving flows from each link, one can separate the different processes used to model traffic flow within a single link, and the interactions between links at junctions. Calculations of sending and receiving flows can be done independently, in parallel, for each link. It is even possible to use different link models on different links (say, one link with a spatial queue model

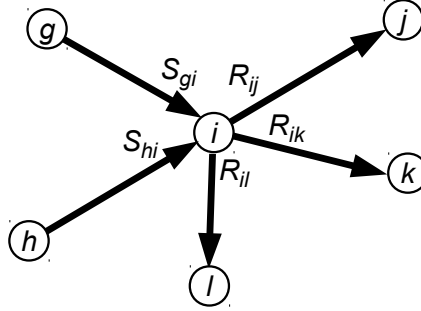


Figure 10.6: Sending and receiving flows interact at nodes to produce turning movement flows.

and another with a point queue).

All that a node model needs is the sending flow of each of the links which enter this node, and the receiving flow for each link leaving this node. (Figure 10.6) That is, when processing node i at time t , we assume that $S_{hi}(t)$ and $R_{ij}(t)$ have respectively been calculated for each incoming link $(h, i) \in \Gamma^{-1}(i)$ and each outgoing link $(i, j) \in \Gamma(i)$. The task is to determine how many vehicles move from each incoming link to each outgoing link during the t -th time interval; denote this value by $y_{hij}(t)$. This is called the *turning movement flow* from (h, i) to (i, j) , or more compactly, the turning movement $[h, i, j]$. Let $\Xi(i)$ denote the set of allowable turning movements at node i ; this provides a natural way to model turn prohibitions, U-turn prohibitions, and so forth.

As an example of this notation, consider node i in Figure 10.6. If all of the turning movements are allowed, then $\Xi(i)$ has six elements: $[g, i, j]$, $[g, i, k]$, $[g, i, l]$, $[h, i, j]$, $[h, i, k]$, and $[h, i, l]$. This set might not include all six elements — for instance, if the left turn from approach (g, i) to (i, j) is prohibited, then $\Xi(i)$ would exclude $[g, i, j]$.

Many different node models can be used. This section discusses three simple node models — links in series, diverges, and merges — to illustrate the general concepts. Later in this chapter, we discuss node models that can be used for signalized intersections, all-way stops, and other types of intersections. There are many variations of all of these, but they share some common principles, which are discussed here.

Any node model must satisfy a number of constraints and principles, for all time intervals t :

1. Vehicles will not voluntarily hold themselves back. That is, it should be impossible to increase any turning movement flow $y_{hij}(t)$ without violating one of the constraints listed below, or one of the other constraints imposed by a specific node model.
2. Each turning movement flow must be nonnegative, that is, $y_{hij}(t) \geq 0$ for

each $[h, i, j] \in \Xi(i)$.

3. Any turning movement which is not allowable has zero flow, that is, $y_{hij}(t) = 0$ whenever $[h, i, j] \notin \Xi(i)$.
4. For each incoming link, the sum of the turning movement flows out of this link cannot exceed the sending flow, since by definition those are the only vehicles which could possibly leave that link in the next time step:

$$\sum_{(i,j) \in \Gamma(i)} y_{hij}(t) \leq S_{hi}(t) \quad \forall (h, i) \in \Gamma^{-1}(i) \quad (10.7)$$

The sum on the left may be less than $S_{hi}(t)$, because it is possible that some vehicles cannot leave (h, i) due to obstructions from a downstream link or from the node itself (such as a red signal).

5. For each outgoing link, the sum of the turning movement flows into this link cannot exceed the receiving flow, since that is the maximum number of vehicles that link can accommodate:

$$\sum_{(h,i) \in \Gamma^{-1}(i)} y_{hij}(t) \leq R_{ij}(t) \quad \forall (i, j) \in \Gamma(i) \quad (10.8)$$

The sum on the left may be less than $R_{ij}(t)$, because there may not be enough vehicles from upstream links to fill all of the available space in the link.

6. Route choices must be respected. That is, the values $y_{hij}(t)$ must be compatible with the directions travelers wish to go based on their chosen paths; we cannot reassign them on the fly in order to increase a y_{hij} value.
7. The *first-in, first-out (FIFO) principle* must be respected. This is closely related to the previous property; we cannot allow vehicles to “jump ahead” in queue to increase a y_{hij} value, unless there is a separate turn lane or other roadway geometry which can separate vehicles on different routes. (Figure 10.7).
8. The *invariance principle* must be respected. If the outflow from a link is less than its sending flow, then increasing the sending flow further could not increase the outflow beyond its current value. Likewise, if the inflow to a link is less than its receiving flow, then increasing the receiving flow could not increase the inflow to the link beyond its current value. In other words, if the sending (or receiving) flow is not “binding,” then its specific value cannot matter for the actual flows.

The invariance principle warrants additional explanation. In the first case, assume that flow into link (i, j) is restricted by its receiving flow R_{ij} . This means that more vehicles wish to enter the link than can be accommodated, so a queue must form on the upstream link. But when there is a queue on the

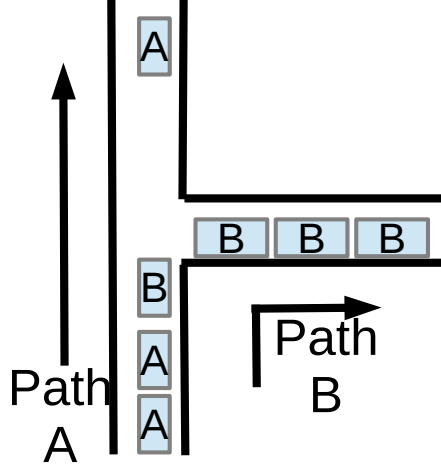


Figure 10.7: The first-in, first-out principle implies that vehicles waiting to turn will obstruct vehicles further upstream.

upstream link, its sending flow will increase to the capacity $q_{max}^{hi}\Delta t$. If the flow y_{hij} increases in response, then the queue might be cleared in the next time interval, the sending flow would drop, a new queue would form, and so on, with the flows oscillating between time intervals. This is unrealistic, and is an artifact introduced by choosing a particular discretization, not a traffic phenomenon one would expect in the field. In the second case, if flow onto link (i, j) is restricted by the sending flow S_{hi} , then there is more space available on link (i, j) than vehicles wish to enter, a short time later the receiving flow will increase to the capacity $q_{max}^{ij}\Delta t$. If y_{hij} would increase because of this increase in the receiving flow, again an unrealistic oscillation would be seen. The exercises ask you to compare some node models which violate the invariance principle to those which do not.

10.2.1 Links in series

The simplest node to model is one with exactly one incoming link, and exactly one outgoing link. (Figure 10.8). This may seem like a trivial node, since there is no real “intersection” here. However, they are often introduced to reflect changes within a link. For instance, if a freeway reduces from three lanes in a direction to two lanes, this reduction in capacity and jam density can be modeled by introducing a node at the point where the lane drops. In this way, each link can have a homogeneous capacity and jam density, and we can simplify the notation — in the link models above, we distinguished between q_{max}^{\uparrow} and q_{max}^{\downarrow} at the two ends of the link. Now we can just use q_{max} for the entire link,

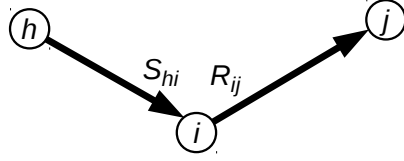


Figure 10.8: Two links in series.

including both ends.

For concreteness, let the incoming link be (h, i) and the outgoing link be (i, j) . In this case, there is only one turning movement, so $\Xi(i) = \{[h, i, j]\}$, and the only turning movement flow we need concern ourselves with is $y_{hij}(t)$. In the case of two links in series, the formula is simple:

$$y_{hij}(t) = \min\{S_{hi}(t), R_{ij}(t)\}, \quad (10.9)$$

that is, the number of vehicles moving from link (h, i) to (i, j) during the t -th time interval is the lesser of the sending flow from the upstream link in that time interval, and the receiving flow of the downstream link. For instance, if the upstream link sending flow is 10 vehicles, while the downstream link receiving flow is 5 vehicles, a total of 5 vehicles will successfully move from the upstream link to the downstream one, because that is all there is space for. If the upstream link sending flow is 3 vehicles and the downstream link receiving flow is 10 vehicles, 3 vehicles will move from the upstream link to the downstream one, because that is all the vehicles that are available.

We can check that this node model satisfies all of the desiderata from Section 10.2. Going through each of these conditions in turn:

1. The flow $y_{hij}(t)$ is chosen to be the minimum of the upstream sending flow, and the downstream receiving flow; any value larger than this would violate either the sending flow constraint or the receiving flow constraint.
2. The sending and receiving flows should both be nonnegative, regardless of the link model, so $y_{hij}(t)$ is as well.
3. There is only one turning movement, so this constraint is trivially satisfied.
4. The formula for $y_{hij}(t)$ ensures it cannot be greater than the upstream sending flow. (This condition simplifies since there is only one incoming and outgoing link, so the summation and “for all” quantifier can be disregarded.)
5. The formula for $y_{hij}(t)$ ensures it cannot be greater than the downstream receiving flow. (This condition simplifies in the same way.)

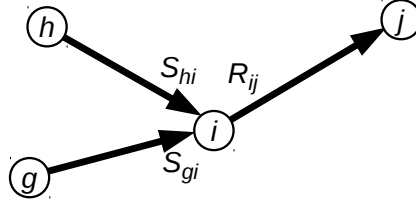


Figure 10.9: Prototype merge node.

6. Route choice is irrelevant when two links meet in series, since all incoming vehicles must exit by the same link.
7. FIFO is also irrelevant, since all vehicles entering the node behave in the same way. (This would not be the case if there was more than one exiting link, and vehicles were on different paths.) So we don't have to worry about FIFO when calculating the turning movement flow.
8. To see that the formula satisfies the invariance principle, we have to check two conditions. If the outflow from (h, i) is less than the sending flow, this means that $y_{hij}(t) = R_{ij}(t)$, and the second term in the minimum of equation (10.9) is binding. Increasing the sending flow (the first term in the minimum) further would not affect its value. Similarly, if the inflow to (i, j) is less than its receiving flow, this means that $y_{hij}(t) = S_{hi}(t)$, and the first term in the minimum is binding. Increasing the receiving flow (the second term) would not affect its value either.

10.2.2 Merges

A merge node has only one outgoing link (i, j) , but more than one incoming link, here labeled (g, i) and (h, i) , as in Figure 10.9. This section only concerns itself with the case of only two upstream links, and generalizing to the case of additional upstream links is left as an exercise. Here $\Xi(i) = \{[g, i, j], [h, i, j]\}$, and we want to calculate the rate of flow from the upstream links to the downstream one, that is, the flow rates $y_{gij}(t)$ and $y_{hij}(t)$. As you might expect, the main quantities of interest are the upstream sending flows $S_{gi}(t)$ and $S_{hi}(t)$, and the downstream receiving flow $R_{ij}(t)$. We assume that these values have already been computed by applying a link model. For brevity, we will omit the time index in the rest of this section — it is implicit that all calculations are done with the sending and receiving flows at the current time step.

There are three possibilities, one corresponding to free flow conditions at the merge, one corresponding to congestion with queues growing on both upstream links, and one corresponding to congestion on only one upstream link. For the merge to be freely flowing, both upstream links must be able to transmit all of the flow which seeks to leave them, and the downstream link must be able to

accommodate all of this flow. Mathematically, we need $S_{gi} + S_{hi} \leq R_{ij}$, and if this is true then we simply set $y_{gij} = S_{gi}$, and $y_{hij} = S_{hi}$.

In the second case, there is congestion (so $S_{gi} + S_{hi} > R_{ij}$), and furthermore, flow is arriving fast enough on both upstream links for a queue to form at each of them. Empirically, in such cases the flow rate from the upstream links is approximately proportional to the capacity on these links, that is,

$$\frac{y_{gij}}{y_{hij}} = \frac{q_{max}^{gi}}{q_{max}^{hi}} \quad (10.10)$$

A little thought should convince you that this relationship is plausible. Furthermore, in the congested case, all of the available downstream capacity will be used, so

$$y_{gij} + y_{hij} = R_{ij} \quad (10.11)$$

Substituting (10.10) into (10.11) and solving, we obtain

$$y_{gij} = \frac{q_{max}^{gi}}{q_{max}^{gi} + q_{max}^{hi}} R_{ij} \quad (10.12)$$

with a symmetric expression for y_{hij} .

The third case is perhaps a bit unusual. The merge is congested ($S_{gi} + S_{hi} > R_{ij}$), but a queue is only forming on one of the upstream links. This may happen if the flow on one of the upstream links is much less than the flow on the other. In this case, the proportionality rule allows all of the sending flow from one link to enter the downstream link, with room to spare. This “spare capacity” can then be consumed by the other approach. If link (g, i) is the link which cannot send enough flow to meet the proportionality condition, so that

$$S_{gi} < \frac{q_{max}^{gi}}{q_{max}^{gi} + q_{max}^{hi}} R_{ij}, \quad (10.13)$$

then the two flow rates are $y_{gij} = S_{gi}$ and $y_{hij} = R_{ij} - S_{gi}$: one link sends all of the flow it can, and the other link consumes the remaining capacity. The formulas are reversed if it is link (h, i) that cannot send enough flow to meet its proportionality condition.

Exercise 7 asks you to show that the second and third cases can be handled by the single equation

$$y_{gij} = \text{med} \left\{ S_{gi}, R_{ij} - S_{hi}, \frac{q_{max}^{gi}}{q_{max}^{gi} + q_{max}^{hi}} R_{ij} \right\}, \quad (10.14)$$

where med refers to the median of a set of numbers. This formula applies whenever $S_{gi} + S_{hi} > R_{ij}$. An analogous formula holds for the other approach by swapping the g and h indices. Exercise 8 asks you to verify the desiderata of Section 10.2 are satisfied by this equation.

For certain merges, it may not be appropriate to assign flow proportional to the capacity of the incoming links. Rules of the road, signage, or signalization

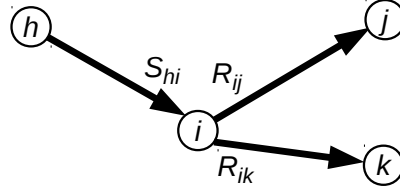


Figure 10.10: Prototype diverge node.

might allocate the capacity of the downstream link differently. In general, the share of the downstream receiving flow that is allocated to approaches (g, i) and (h, i) can be written as ψ_{gi} and ψ_{hi} , respectively, with $\psi_{gi} + \psi_{hi} = 1$ and both of them nonnegative. A more general form of the merge equation can then be written as

$$y_{gij} = \text{med} \{S_{gi}, R_{ij} - S_{hi}, \psi_{gi} R_{ij}\}, \quad (10.15)$$

with a similar formula for (h, i) .

10.2.3 Diverges

A diverge node is one with only one incoming link (h, i) , but more than one outgoing link, as in Figure 10.10. This section concerns itself with the case of only two downstream links. The exercises ask you to generalize to the case of three downstream links, using the same concepts. Let these two links be called (i, j) and (i, k) , so $\Xi(i) = \{[h, i, j], [h, i, k]\}$. Our interest is calculating the rate of flow from the upstream link to the downstream ones, that is, the flow rates y_{hij} and y_{hik} . We assume that the sending flow S_{hi} and the receiving flows R_{ij} and R_{ik} have already been calculated. Unlike links in series or merges, we also need to represent some model of route choice, since some drivers may choose link (i, j) , and others link (i, k) . Let p_{ij} and p_{ik} be the proportions of drivers choosing these two links during the t -th time interval, respectively. Naturally, p_{ij} and p_{ik} are nonnegative, and $p_{ij} + p_{ik} = 1$. Like the sending and receiving flows, these values can change with time, but to avoid cluttering formulas we will leave the time indices off of p values unless it is unclear which time step we are referring to.

There are two possibilities, one corresponding to free flow conditions at the diverge, and the other corresponding to congestion. What does “free flow” mean? For the diverge to be freely flowing, *both* of the downstream links must be able to accommodate the flow which seeks to enter them. The rates at which vehicles want to enter the two links are $p_{ij}S_{hi}$ and $p_{ik}S_{hi}$, so if both downstream links can accommodate this, we need $p_{ij}S_{hi} \leq R_{ij}$ and $p_{ik}S_{hi} \leq R_{ik}$. In this case we simply have $y_{hij} = p_{ij}S_{hi}$ and $y_{hik} = p_{ik}S_{hi}$: all of the flow which wants to leave the diverge can.

The case of congestion is slightly more interesting, and requires making assumptions about how drivers will behave. One common assumption is that

flow waiting to enter one link at a diverge will obstruct every other vehicle on the link (regardless of which link it is destined for). This most obviously represents the case where the upstream link has only a single lane, so any vehicle which has to wait will block any vehicle behind it; but this model is commonly used even in other cases.¹ When there is congestion, only some fraction ϕ of the upstream sending flow can move. The assumption that any vehicle waiting blocks every vehicle upstream implies that this same fraction applies to both of the downstream links, so $y_{hij} = \phi p_{ij} S_{hi}$ and $y_{hik} = \phi p_{ik} S_{hi}$.

So, how to calculate ϕ ? The inflow rate to a link cannot exceed its receiving flow, so $y_{hij} = \phi p_{ij} S_{hi} \leq R_{ij}$ and $y_{hik} = \phi p_{ik} S_{hi} \leq R_{ik}$, or equivalently $\phi \leq R_{ij}/p_{ij} S_{hi}$ and $\phi \leq R_{ik}/p_{ik} S_{hi}$. Every vehicle which can move will, so

$$\phi = \min \left\{ \frac{R_{ij}}{p_{ij} S_{hi}}, \frac{R_{ik}}{p_{ik} S_{hi}} \right\} \quad (10.16)$$

Furthermore, we can introduce the uncongested case into this equation as well, and state

$$\phi = \min \left\{ \frac{R_{ij}}{p_{ij} S_{hi}}, \frac{R_{ik}}{p_{ik} S_{hi}}, 1 \right\} \quad (10.17)$$

regardless of whether there is congestion at the diverge or not. Why? If the diverge is at free flow, then $\phi = 1$, but $R_{ij}/p_{ij} S_{hi} \geq 1$ and $R_{ik}/p_{ik} S_{hi} \geq 1$. Introducing 1 into the minimum therefore gives the correct answer for free flow. Furthermore, if the diverge is not at free flow, then either $R_{ij}/p_{ij} S_{hi} < 1$ or $R_{ik}/p_{ik} S_{hi} < 1$, so adding 1 does not affect the minimum value. Therefore, this formula is still correct even in the congested case. Exercise 13 asks you to verify the desiderata of Section 10.2 are satisfied by this equation.

10.3 Combining Node and Link Models

This section describes how node and link models are combined, to complete the network loading process. We must also describe what happens at origins and destination (zone) nodes. For the algorithm in this section, we assume that the only links connected to zones are special links called *centroid connectors*, which do not represent a specific roadway so much as a collection of small local streets used by travelers entering or leaving a specific neighborhood. It is common to give centroid connectors between origins and ordinary nodes a very large (even infinite) jam density, and centroid connectors between ordinary nodes and destinations a very large (even infinite) capacity. In both cases the free flow time should be small. These considerations reflect the ideas that centroid connectors should not experience significant congestion (or else they should be modeled as proper links in the network), and simply convey flow from origin nodes and to destination nodes with as little interference as possible. It is common to forbid

¹For instance, this can represent drivers attempting to “queue jump” by cutting into the turn lane at the last moment. Or, if the turn lane is long, it may be appropriate to treat the diverge at the point where the turn lane begins, as opposed to at the physical diverge.

travelers from using centroid connectors except to start and end their trips, that is, to exclude the use of centroid connectors as “shortcuts.” This can be done either by transforming the underlying network, having origins and destinations be distinct nodes adjacent to “one-way” centroid connectors, or by excluding such paths when finding paths for travelers, as discussed in Chapter 11.

If centroid connectors are set up in this way, then flow entering the network can simply be added to the upstream ends of their centroid connectors, and flow leaving the network at destinations can simply vanish, without any constraints in either case. The network loading algorithm can then be stated as follows:

1. Initialize all counts and the time index: $N_{ij}^\uparrow(0) \leftarrow 0$ and $N_{ij}^\downarrow(0) \leftarrow 0$ for all links (i, j) , $t \leftarrow 0$.
2. Use a link model to calculate sending and receiving flows S_{ij} and R_{ij} for all links.
3. Use a node model to calculate transition flows y_{ijk} for all nodes j except for zones.
4. Update cumulative counts: for each non-zone node i , perform

$$N_{hi}^\downarrow(t+1) \leftarrow N_{hi}^\downarrow(t) + \sum_{(i,j) \in \Gamma(i)} y_{hij} \quad (10.18)$$

for each upstream link (h, i) , and

$$N_{ij}^\uparrow(t+1) \leftarrow N_{ij}^\uparrow(t) + \sum_{(h,i) \in \Gamma^{-1}(i)} y_{hij} \quad (10.19)$$

for each downstream link (i, j) .

5. Load trips: for all origins r , let $D_r(t) = \sum_{s \in Z} d^{rs}(t)$ be the total demand starting at this node, and for each centroid connector (r, i) , let p_{ri} be the fraction of demand beginning their trips on that connector. Set $N_{ri}^\uparrow(t+1) \leftarrow N_{ri}^\uparrow(t) + p_{ri}D_r(t)$ for each connector (r, i) .
6. Terminate trips at each destination s : for each centroid connector (i, s) , set $N_{is}^\downarrow(t+1) \leftarrow N_{is}^\downarrow(t) + S_{is}$.
7. Increment the time: $t \leftarrow t+1$. If t equals the time horizon \bar{T} , then stop. Otherwise, return to step 2.

Figure 10.11 and Table 10.3 illustrate this process on a network with three links in series, using the spatial queue link model. The link parameters are shown in Figure 10.11; notice that link (i, j) has a larger capacity at the upstream end than at the downstream end. The node model at i is the “links in series” model discussed in Section 10.2.1. The node model at j is a modified version of the “links in series” node model: when $5 \leq t < 10$, $y_{ijs} = 0$ regardless of the sending and receiving flows of (i, j) and (j, s) ; at all other times the “links

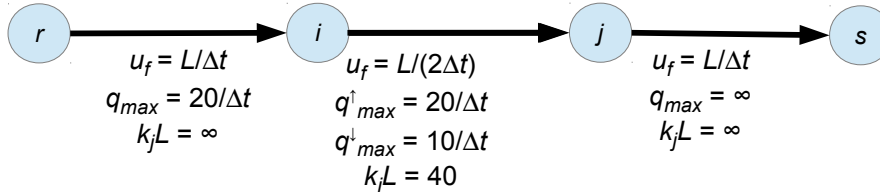


Figure 10.11: Example for network loading algorithm.

in series” model is used. This might reflect a red traffic signal, or a closed drawbridge between these time intervals, since no flow can move through the node. In this example, the flow downstream is first interrupted when moving from the centroid connector (r, i) to the link (i, j) , whose capacity is lower than the rate at which vehicles are being loaded at origin r . This can be seen by examining the difference between the N^\uparrow and N^\downarrow values for link (r, i) during the initial timesteps. The difference between these values gives the total number of vehicles on the link at that instance in time. Since N^\uparrow is increasing at a faster rate than N^\downarrow , vehicles are accumulating on the link, in a queue at the downstream end. There is no queue at node j , because no bottleneck exists there. Compare $N_{ij}^\uparrow(t)$ and $N_{ij}^\downarrow(t)$ when $3 \leq t \leq 5$. Both the upstream and downstream count values increase at the same rate, which means there is no net accumulation of vehicles.

At $t = 5$, the flow through node j drops to zero, which introduces a further bottleneck. The impacts of the bottleneck are first seen at $t = 6$: 40 vehicles are now on link (i, j) , up from 30. As a result, the receiving flow R_{ij} drops to zero, because the link is full. Therefore, the node model at i restricts any additional inflow to link (i, j) , and the queue on (r, i) grows at an even faster rate than before, even though the number of new vehicles loaded onto the network has dropped. At $t = 10$, the bottleneck at node j is released, and vehicles begin to move again. At $t = 25$, the upstream and downstream counts are equal on all links, which means that the network is empty. All vehicles have reached their destination.

10.4 Elementary Traffic Flow Theory

This section provides an introduction to the hydrodynamic theory of traffic flow, the basis of several widely-used link models which are more realistic than the point or spatial queue models. In this theory, traffic is modeled as a compressible fluid. Its primary advantage is its simplicity: it can capture many important congestion phenomena, while remaining tractable for large networks. Of course, vehicles are not actually molecules of a fluid, but are controlled by drivers with heterogeneous behavior, who drive vehicles of heterogeneous size, power, and so on. By treating vehicles as identical particles, we are ignoring such distinctions; as an example of these limitations, we assume that no overtaking occurs within

Table 10.3: Example of network loading algorithm.

t	d_{rs}	Link (r, i)				Node i	Link (i, j)				Node j	Link (j, s)			
		R_{ri}	N_{ri}^\uparrow	N_{ri}^\downarrow	S_{ri}	$y_{r \rightarrow i}$	R_{ij}	N_{ij}^\uparrow	N_{ij}^\downarrow	S_{ij}	$y_{i \rightarrow s}$	R_{js}	N_{js}^\uparrow	N_{js}^\downarrow	S_{js}
0	15	∞	0	0	0	0	20	0	0	0	0	∞	0	0	0
1	15	∞	15	0	15	15	20	0	0	0	0	∞	0	0	0
2	15	∞	30	15	15	15	20	15	0	0	0	∞	0	0	0
3	15	∞	45	30	15	10	10	30	0	10	10	∞	0	0	0
4	15	∞	60	40	20	10	10	40	10	10	10	∞	10	0	10
5	15	∞	75	50	20	10	10	50	20	10	0	∞	20	10	10
6	10	∞	90	60	20	0	0	60	20	10	0	∞	20	20	0
7	10	∞	100	60	20	0	0	60	20	10	0	∞	20	20	0
8	10	∞	110	60	20	0	0	60	20	10	0	∞	20	20	0
9	10	∞	120	60	20	0	0	60	20	10	0	∞	20	20	0
10	10	∞	130	60	20	0	0	60	20	10	10	∞	20	20	0
11	10	∞	140	60	20	10	10	60	30	10	10	∞	30	20	10
12	10	∞	150	70	20	10	10	70	40	10	10	∞	40	30	10
13	0	∞	160	80	20	10	10	80	50	10	10	∞	50	40	10
14	0	∞	160	90	20	10	10	90	60	10	10	∞	60	50	10
15	0	∞	160	100	20	10	10	100	70	10	10	∞	70	60	10
16	0	∞	160	110	20	10	10	110	80	10	10	∞	80	70	10
17	0	∞	160	120	20	10	10	120	90	10	10	∞	90	80	10
18	0	∞	160	130	20	10	10	130	100	10	10	∞	100	90	10
19	0	∞	160	140	20	10	10	140	110	10	10	∞	110	100	10
20	0	∞	160	150	10	10	10	150	120	10	10	∞	120	110	10
21	0	∞	160	160	0	0	10	160	130	10	10	∞	130	120	10
22	0	∞	160	160	0	0	20	160	140	10	10	∞	140	130	10
23	0	∞	160	160	0	0	20	160	150	10	10	∞	150	140	10
24	0	∞	160	160	0	0	20	160	160	0	0	∞	160	150	10
25	0	∞	160	160	0	0	20	160	160	0	0	∞	160	160	0

a link (because the particles are identical, none of them has any reason to move faster than another). But seen as a first-order approximation, the hydrodynamic theory provides a simple and tractable way to model network traffic. It is also possible to derive certain aspects of the hydrodynamic theory from more behavioral models representing car-following.

In contrast to the link models we will ultimately use in dynamic traffic assignment, fluid-based traffic models are formulated in *continuous* space and time, so it makes sense to talk about the state of traffic at any point in time (not just at the integer timesteps $0, 1, \dots, \bar{T}$). Link models which are based on fluid models will convert these continuous quantities to discrete ones.

10.4.1 Traffic state variables

We start by modeling a roadway link as a one-dimensional object, using x to index the distance from the upstream end of the link, and using t to index the current time. At any point x , and at any time t , the state of traffic can be described by three fundamental quantities: the *flow* q , the *density* k , and the *speed* u . Each of these can vary over space and time, so $q(x, t)$, $k(x, t)$, and $u(x, t)$ can be seen as functions defined over all x values on the link, and all t values in the analysis period. However, when there is no ambiguity (e.g., only looking at a single point at a single time), we can simply write q , k , and u without providing the space and time coordinates. Both x and t are treated as continuous variables, as are the vehicles themselves, allowing derivatives of q , k , and u to be meaningfully defined. This assumption is imported from fluid mechanics, where the molecules are so small and numerous that there is essentially no error in approximating the fluid as a continuum. For traffic flow, this assumption is not so trivial, and is one of the drawbacks of hydrodynamic models.

Flow is defined as the rate at which vehicles pass a stationary point, and commonly has units of vehicles per hour. In the field, flow can be measured using point detectors (such as inductive loops) which record the passage of each vehicle at a fixed location. Flow can be thought of as a temporal concentration of vehicles. Density, on the other hand, is defined to be the spatial concentration of vehicles at a given time, and is measured in vehicles per unit length (commonly vehicles per kilometer or vehicles per mile). Density can be obtained from taking a photograph of a link, noting the concentration of vehicles at different locations at a single instant in time. The speed is the instantaneous rate at which the vehicles themselves are traveling, and is measured in units such as kilometers per hour or miles per hour. Speed can be directly measured from radar detectors.

These three quantities are not independent of each other. As a start, there is the basic relationship

$$q = uk \tag{10.20}$$

which must hold at each point and time. (You should check the dimensions of the quantities in this formula to verify their compatibility.) If this equation is not evident to you from the definitions of flow, density, and speed, imagine that

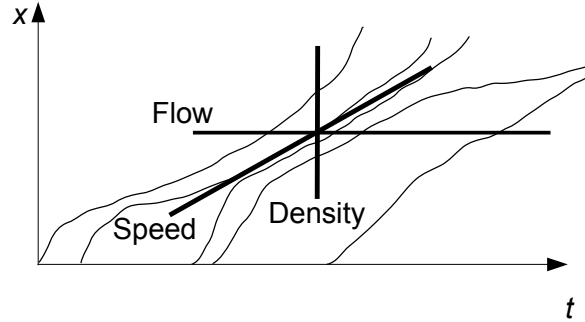


Figure 10.12: Trajectory diagram illustrating speed, flow, and density.

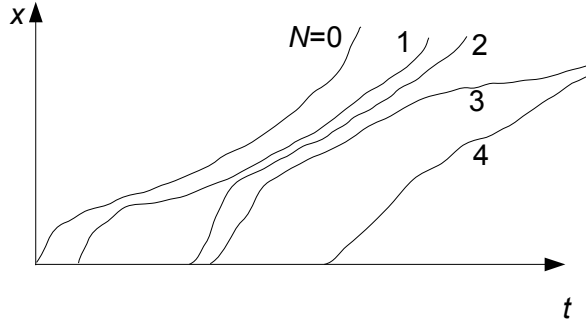
we want to know the number of vehicles ΔN which will pass a fixed point over a small, finite time interval Δt . If the speed of vehicles is u , any upstream vehicle within a distance of $u\Delta t$ from the fixed point will pass during the next Δt time units, and the number of such vehicles is

$$\Delta N = ku\Delta t. \quad (10.21)$$

The flow rate is approximately $\Delta N/\Delta t$, and the formula (10.20) then follows from taking limits as the time increment Δt shrinks to zero.

Figure 10.12 is a *trajectory diagram* showing the locations of vehicles on the link over time — the horizontal axis denotes time, and the vertical axis denotes space, with the upstream end of the link at the bottom and the downstream end at the top. Speed, flow, and density can all be interpreted in terms of these trajectories. The speed of a vehicle at any point in time corresponds to the slope of its trajectory there. Flow is the rate at which vehicles pass a fixed point: on a trajectory diagram, a fixed point in space is represented by a horizontal line. Time intervals when more trajectories cross this horizontal line have higher flow, and when fewer trajectories cross this line, the flow is lower. Density is the spatial concentration of vehicles at a particular instant in time: on a trajectory diagram, a specific instant is represented by a vertical line. Where more trajectories cross this vertical line, the density is higher, and where fewer trajectories cross, the density is lower.

However, this equation by itself is not enough to describe anything of real interest in traffic flow. Another equation, based on vehicle conservation principles, is described in the next subsection. The Lighthill-Whitham-Richards model, described at the end of this section, makes a further assumption about the relationships of three state variables. These relationships are enough to specify the network loading problem as the solution to a well-defined system of partial differential equations.

Figure 10.13: Vehicle trajectories are contours of $N(x, t)$

10.4.2 Cumulative counts and conservation

The hydrodynamic theory can be simplified by introducing a fourth variable N , again defined at each location x and at each time t . This new variable N is referred to as the *cumulative count*, and has a similar interpretation to the N^\uparrow and N^\downarrow counts defined at the start of Section 10.1, but now applied at any point on the link, not just the upstream and downstream ends. Imagine that each vehicle is labeled with a number; for instance, the vehicle at the upstream end of the link at $t = 0$ may be labeled as zero. The next vehicle that enters the link is then labeled as one, the next vehicle as two, and so forth. The numbering does not have to start at zero. What is important that each entering vehicle be given consecutive numbers. Then $N(x, t)$ gives the number of the vehicle at location x at time t (keeping in mind that we are modeling vehicles as a continuous fluid, so $N(x, t)$ need not be an integer). The contours of $N(x, t)$ then give the trajectories of individual vehicles in the traffic stream. (Figure 10.13)

It is worthwhile to point out a potential source of confusion here. In a trajectory diagram like Figure 10.13, space (indexed by x) is conventionally denoted on the vertical axis, and time (indexed by t) on the horizontal axis. This has the unfortunate side-effect of making the “ x -axis” the vertical one. Second, it is conventional to list the spatial component before the time component, so a point like $(1, 2)$ refers to $x = 1$ and $t = 2$. This means that on a trajectory diagram, the *vertical* component is listed first, as opposed to typical Cartesian coordinates where the *horizontal* component is given first. Unfortunately, both of these conventions are so well-established in the transportation engineering literature that it is best to simply highlight them and become comfortable using them.

The quantity $N(x, t)$ is called a cumulative count for the following reason: at the upstream end of the link ($x = 0$), the quantity $N(0, t)$ gives the cumulative number of vehicles which have entered the link up to time t . Furthermore, at any location x , the difference $N(x, t') - N(x, t)$ gives the number of vehicles which passed point x between times t and t' . Taking the limit as t' approaches

t provides the relationship

$$\frac{\partial N}{\partial t} = q, \quad (10.22)$$

which holds for every x and t where N is differentiable. That is, at any fixed spatial location, the rate at which N increases in time is exactly the flow rate q .

The cumulative counts can also be related to density with a similar argument. At any point in time t , the difference $N(x', t) - N(x, t)$ gives the number of vehicles which lie between locations x and x' at time t . However, we have to be careful regarding the sign convention regarding N described above: since vehicles are numbered in the order they enter the link, in any platoon of vehicles N decreases as we move from the following vehicles to the lead vehicle. So, if $x' > x$, then $N(x', t) \leq N(x, t)$. Thus, taking the limit as x' approaches x , we must have

$$\frac{\partial N}{\partial x} = -k, \quad (10.23)$$

which again applies wherever N is differentiable, and where the negative sign in the formula results from our sign convention. In this respect, the cumulative counts N can be seen as the most basic description of traffic flow: if we are given $N(x, t)$ at all points x and times t , we can calculate $q(x, t)$ and $k(x, t)$ by using equations (10.22) and (10.23), and therefore $u(x, t)$ everywhere by using (10.20).

Furthermore, wherever the flow and density are themselves continuously differentiable functions, Clairaut's theorem states that the mixed second partial derivatives of N must be equal, that is,

$$\frac{\partial^2 N}{\partial x \partial t} = \frac{\partial^2 N}{\partial t \partial x}, \quad (10.24)$$

so substituting the relationships (10.22) and (10.23) and rearranging, we have

$$\frac{\partial q}{\partial x} + \frac{\partial k}{\partial t} = 0. \quad (10.25)$$

This is an expression of vehicle conservation, that is, vehicles do not appear or disappear at any point. This equation must hold everywhere that these derivatives exist.²

Equations (10.22) and (10.23) are useful in another way. If (x_1, t_1) and (x_2, t_2) are any two points in space and time, the difference in cumulative count number between these points is given by the line integral

$$N(x_2, t_2) - N(x_1, t_1) = \int_C q \, dt - k \, dx, \quad (10.26)$$

where C is any curve connecting (x_1, t_1) and (x_2, t_2) . Because vehicles are conserved, this line integral does not depend on the specific path taken. This

²Of course, vehicle conservation must hold even when these derivatives do not exist, it is just that the formula (10.25) is meaningless there. We have to enforce flow conservation in a different way at such points.

is helpful, because we can choose a path which is easy to integrate along. In what follows, we often choose the straight line connecting these two points as the integration path.

Two issues concerning cumulative counts are often confusing, and are worth further explanation. First, the cumulative counts can only be meaningfully compared within the same link. Each link maintains its own counts, and the number associated with a vehicle may change when traveling between links. Instead of thinking of the cumulative count as being a label permanently associated with a vehicle, it is better to think about it as a label given to the vehicle *by a specific link*, and each link maintains its labels independently of all of the other links. Each link simply gives successive numbers to each new vehicle entering the link, without having to coordinate its numbering scheme with other links. This is needed to ensure that link models can function independently, and because in complex networks it is usually impossible to assign “permanent” numbers to vehicles such that any two vehicles entering a link consecutively have consecutive numbers. Second, within any given link, only the *difference* in cumulative counts is meaningful, the absolute numbers do not have specific meaning. For instance, it may be relevant that 10 vehicles entered the link between times 5 and 6, but the specific numbers of these vehicles are not important. A common convention is to have the first vehicle entering the link be assigned the number 0, the next vehicle the number 1, and so forth, but this is not required. In cases where there are already vehicles on the link at the start of the modeling period, the first vehicle entering the link may be assigned a higher number (because the vehicles already on the link must have lower numbers, and a modeler’s aesthetics may prefer nonnegative vehicle counts). But there would be nothing wrong with a negative cumulative count either. In this regard, different choices of the “zero point” are analogous to the different zero points in the Fahrenheit and Celsius temperature scales: either one will give you correct answers as long as you are consistent, and a negative number is not necessarily cause for alarm.

For example, consider a link where $N(x, t) = 1000t - 100x$, with t measured in hours and x measured in miles. The flow at any point and time is $\partial N / \partial t$, which is a constant of 1000 vehicles per hour. Likewise, the density is $-\partial N / \partial x = 100$ vehicles per mile. Using the basic relationship (10.20), the speed must be 10 miles per hour uniformly on the link.

As a more involved example, consider a link which is 1 mile long, with all times measured in hours and distances in miles. If we are given that

$$N(x, t) = 3600t - 120x + \frac{60x^2}{60t + 1}, \quad (10.27)$$

we can calculate the densities and flows everywhere:

$$k(x, t) = -\frac{\partial N}{\partial x} = 120 \left(1 - \frac{x}{60t + 1} \right) \quad (10.28)$$

$$q(x, t) = \frac{\partial N}{\partial t} = 3600 \left(1 - \left(\frac{x}{60t + 1} \right)^2 \right). \quad (10.29)$$

Exercise 24 asks you to verify that the conservation relationship (10.25) is satisfied by explicit computation.

Ordinarily, the $N(x, t)$ map is not given — indeed, the goal of network loading is to calculate it. For if we know $N(x, t)$ everywhere, we can calculate flow, density, and speed everywhere, using equations (10.22), (10.23), and (10.20). So let's assume that we *only* know the density and flow maps (10.28) and (10.29), and try to recover information about the cumulative counts. For the given N map, the vehicle at $x = 1/2$ at $t = 0$ has the number 0. (As discussed above, we do not necessarily have to set the zero point at $N(0, 0)$.) To calculate the number of the vehicle at $x = 1$ and $t = 1$ (the downstream end of the link, one minute later), we can use equation (10.26).

As this equation involves a line integral, we must choose a path between $(x, t) = (1/2, 0)$ and $(1, 1)$. Because of the conservation relationship (10.25), we can choose any path we wish. For the purposes of an example, we will calculate this integral along three different paths, and verify that they give the same answer. Figure 10.14 shows the three paths of integration.

Path A : This path consists of the line segment from $(1/2, 0)$ to $(1, 0)$, followed by the segment from $(1, 0)$ to $(1, 1)$. Because these line segments are parallel to the axes, this reduces the line integral to two integrals, one over x alone, and the other over t alone. We thus have

$$\begin{aligned} N(1, 1) &= \int_A q \, dt - k \, dx = - \int_{1/2}^1 k(x, 0) \, dx + \int_0^1 q(1, t) \, dt \\ &= - \int_{1/2}^1 120(1 - x) \, dx + \int_0^1 60 \left(1 - \left(\frac{1}{t+1} \right)^2 \right) \, dt \\ &= -15 + 30 = 15, \end{aligned}$$

and the vehicle at the downstream end of the link at $t = 1$ has the number 15.

Path B : This path consists of the line segment from $(1/2, 0)$ to $(1/2, 1)$, followed by the segment from $(1/2, 1)$ to $(1, 1)$. As before, we have

$$\begin{aligned} N(1, 1) &= \int_B q \, dt - k \, dx = \int_0^1 q(1/2, t) \, dt - \int_{1/2}^1 k(x, 1) \, dx \\ &= \int_0^1 60 \left(1 - \left(\frac{1/2}{t+1} \right)^2 \right) \, dt - \int_{1/2}^1 120 \left(1 - \frac{x}{2} \right) \, dx \\ &= 52.5 - 37.5 = 15. \end{aligned}$$

Path C : This path is the line segment directly connecting $(1/2, 0)$ to $(1, 1)$. Although this line is not parallel to either axis, the integral actually ends up being the easiest to evaluate, because $x/(t+1)$ is constant along this line, and equal to $1/2$. Therefore $k(x, t) = 60$ at all points along this line,

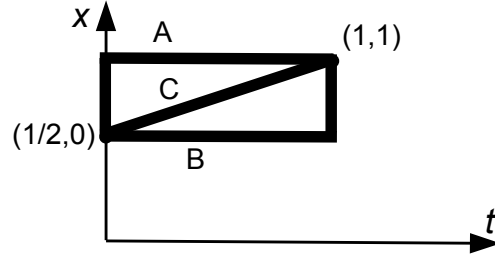


Figure 10.14: Three possible paths for the line integral between $(1/2, 0)$ and $(1, 1)$.

and $q(x, t) = 45$. Since $dx = (1/2)dt$ on this line segment, we have

$$\begin{aligned} N(1, 1) &= \int_C q \, dt - k \, dx = \int_0^1 (45dt - 60(1/2)dt) = \int_0^1 15 \, dt \\ &= 15. \end{aligned}$$

All three integrals gave the same answer (as they must), which we can verify by checking $N(1, 1)$ with equation (10.27). So, we can choose whichever integration path is easiest. In this example, the integrals in Path B involved the most work. The integral in Path C required a bit more setup, but the actual integral ended up being very easy, since q and k were constants along the integration path. Such a path is called a *characteristic*, and will be described in more detail later in this chapter.

10.4.3 The Lighthill-Whitham-Richards model

At this point, we have two relationships between the flow, density, and speed variables: the basic relationship (10.20) and the conservation relationship (10.25). These first two relationships can be derived directly from the definitions of these variables, and can describe a wide range of fluid phenomena — at this point nothing yet has been specific to vehicle flow. The Lighthill-Whitham-Richards (LWR) model provides a third relationship, completing the hydrodynamic theory.³

Specifically, the LWR model postulates that the flow at any point is a function of the density at that point, that is,

$$q(x, t) = Q(k(x, t)) \quad (10.30)$$

for some function Q . Equivalently, by the relationship (10.20), we can assume that the speed at any point depends only on the density at that point. It must

³Lighthill and Whitham published this model in 1955, as the sequel to a paper on flow in rivers. Richards independently proposed an equivalent model in 1956. All three are now given credit for this model.

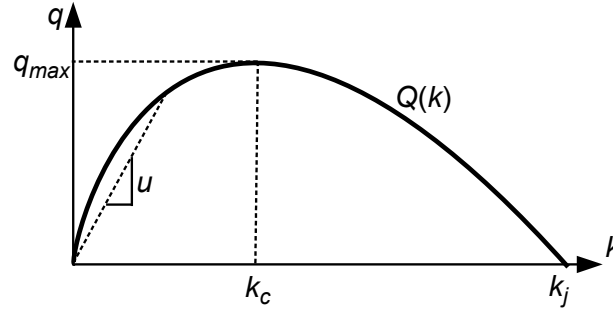


Figure 10.15: Fundamental diagram for the LWR model.

be emphasized that this relationship, unlike (10.20) and (10.25), is an assumed behavior and does not follow from basic principles. In dynamic network loading, we typically assume that this function Q is uniform over space and time on a link, an assumption we adopt in this section for simplicity. It is possible to generalize the results in this section when Q varies over space and time.

The function Q is commonly called the *fundamental diagram*; an example of such a diagram is shown in Figure 10.15. A few features of fundamental diagrams are worth noting: they are concave functions⁴, and typically assumed to be continuous and piecewise differentiable. They also have two zeros: one at $k = 0$ (zero density means zero flow, because no vehicles are present), and another at the *jam density* k_j , corresponding to a maximum density where there is no flow because all vehicles are stopped. At intermediate values of density, the flow is positive, although for a given flow value q , there can be two possible density values corresponding to this flow, one corresponding to uncongested conditions and the other to congested conditions.

If Q is concave, and has two zeros at 0 and k_j , then there is some intermediate point where $Q(k)$ is maximal. This maximal value of Q is called the *capacity* of the link, denoted q_{max} , and the *critical density* k_c is defined to be a value such that $Q(k_c) = q_{max}$. The k values for which $Q(k) < q_{max}$ and $k < k_c$ are referred to as *subcritical*, and reflect uncongested traffic flow; the k values for which $Q(k) < q_{max}$ and $k > k_c$ are *supercritical*, and reflected congested flow.⁵

Using $q = uk$, the speed at any point can be seen as the slope of the secant line connecting the origin to the point on the fundamental diagram correspond-

⁴A function f is concave if $-f$ is convex.

⁵The definitions of subcritical and supercritical in the transportation field are exactly opposite to how these terms are used in fluid mechanics. This is a bit annoying, but this difference in convention reflects differences in the “default state” of flow. Traffic engineers view the default state of traffic flow as being uncongested, when vehicles can move freely and (as we show later) shockwaves only travel downstream. Most fluids have to be moving rather quickly for the same state to occur, and the default state of rivers and many other fluid systems is a slower rate of travel, where waves can move both upstream and downstream — what traffic engineers would describe as a “congested” state.

ing to the density at that point. That is, in the LWR model, the density k at a point completely determines the traffic state. The flow q at that point is obtained from the fundamental diagram Q , and the speed u can then be obtained from equation (10.31).

To summarize, the three equations relating flow, density, and speed are:

$$q(x, t) - k(x, t)u(x, t) = 0 \quad (10.31)$$

$$q(x, t) - Q(k(x, t)) = 0 \quad (10.32)$$

$$\frac{\partial q}{\partial x} + \frac{\partial k}{\partial t} = 0 \quad (10.33)$$

and these equations must hold everywhere (with the exception that (10.33) may not be defined if q or k is not differentiable at a point). Together with initial conditions (such as the values of k along the link at $t = 0$) and boundary conditions (such as the “inflow rates” q at the upstream end $x = 0$ throughout the analysis period, or restrictions on q at the downstream end from a traffic signal), this system of equations can in principle be solved to yield $k(x, t)$ everywhere. Exercise 24 asks you to verify that the $N(x, t)$ map used in the example in the previous section is consistent with the fundamental diagram $Q(k) = \frac{1}{4}k(240 - k)$.

The points where k is not differentiable are known as *shockwaves*, and often correspond to abrupt changes in the density. Figure 10.16 shows an example of several shockwaves associated with the changing of a traffic light. Notice that in region A, the density is subcritical (uncongested); in region B, traffic is at jam density; and in region C, traffic is at critical density and flow is at capacity. The speed of a shockwave can still be determined from conservation principles, even though the conservation equation (10.33) does not apply because the density and flow derivatives do not exist at a shock.

Assume that k_A and k_B are the densities immediately upstream and immediately downstream of the shockwave (Figure 10.17). The corresponding flow rates $q_A = Q(k_A)$ and $q_B = Q(k_B)$ can be calculated from the fundamental diagram, and finally the speeds are obtained as $u_A = q_A/k_A$ and $u_B = q_B/k_B$.

Furthermore, let u_{AB} denote the speed of the shockwave. Then the speed of vehicles in region A *relative to the shockwave* is $u_A - u_{AB}$, and the rate at which vehicles cross the shockwave from region A is $(u_A - u_{AB})k_A$; this is nothing more than equation (10.20) as viewed from the perspective of an observer moving with the shockwave. Likewise, the relative speed of the vehicles in region B is $u_B - u_{AB}$, and the rate at which vehicles cross the shockwave and enter region B is $(u_B - u_{AB})k_B$. Obviously these two quantities must be equal, since vehicles do not appear or disappear at the shock. Equating these flow rates from the left and right sides of the shockwave, we can solve for the shockwave speed:

$$u_{AB} = \frac{q_A - q_B}{k_A - k_B}. \quad (10.34)$$

Notice that this calculated speed is the same regardless of whether A is the upstream region and B the downstream region, or vice versa. This equation also

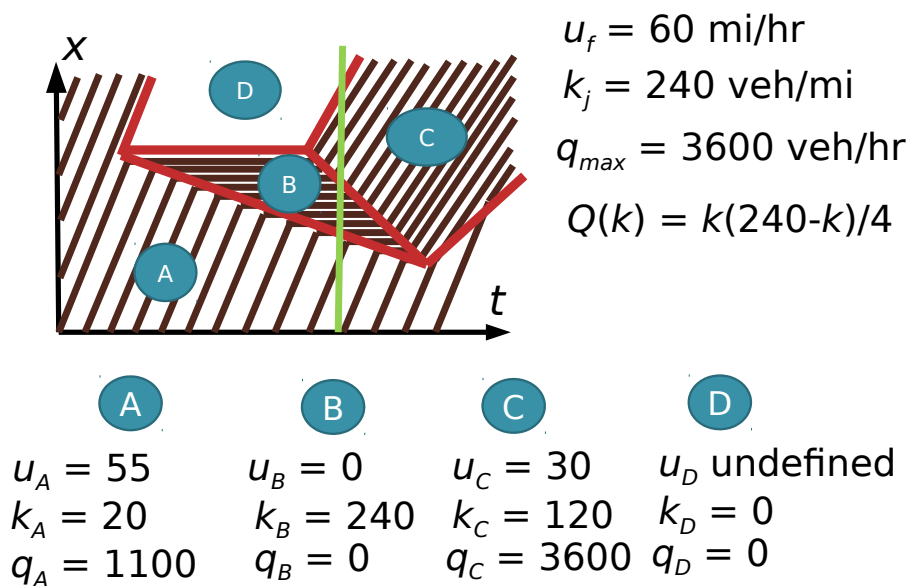


Figure 10.16: Shockwaves associated with a traffic signal; vehicle trajectories indicated in brown and shockwaves in red.

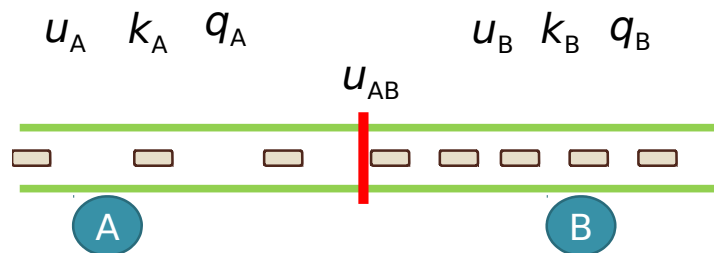


Figure 10.17: Flow conservation produces an equation for shockwave speed.

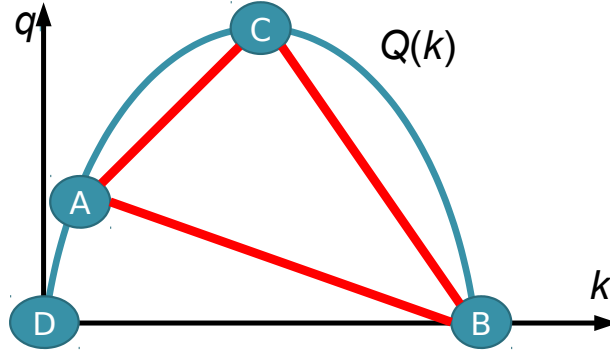


Figure 10.18: Identifying shockwave speed on the fundamental diagram.

has a nice geometric interpretation: the speed of the shockwave is the slope of the line connecting regions A and B on the fundamental diagram (Figure 10.18).

For instance, in Figure 10.16, in region A the flow and density are 1100 vehicles per hour and 20 vehicles per mile, and in region B the flow and density are 0 vehicles per hour and 240 vehicles per mile. Therefore, using (10.34), the shockwave between regions A and B has a speed of $(1100 - 0)/(20 - 240) = -5$ miles per hour. The negative sign indicates that the shockwave is moving *upstream*. Since A represents uncongested traffic, and region B represents the stopped queue at the traffic signal, the interpretation is that the queue is growing at 5 miles per hour. Tracing the derivation of equation (10.34), the rate at which vehicles enter the shockwave is $(55 + 5) \times 20$ from the perspective of region A, or 1200 vehicles per hour. You should check that the same figure is obtained from the perspective of region B. Vehicles are entering the queue faster than the upstream flow rate (1200 vs. 1100 vph) because the queue is growing upstream, moving to meet vehicles as they arrive.

The system of equations (10.31)–(10.33) can then be solved, introducing shockwaves as necessary to accommodate the initial and boundary conditions, using (10.34) to determine their speed. This is the LWR model. However, this theory does not immediately suggest a technique for actually solving the system of partial differential equations, which is the topic of the next subsection.

Notice also that the fundamental diagram determines the maximum speed at which a shockwave can move. Because the fundamental diagram is concave, its slope at any point can never be greater than the free-flow speed $u_f = Q'(0)$, nor less than the slope at jam density $Q'(k_j)$. The absolute values of these slopes give the fastest speeds shockwaves can move in the downstream and upstream directions, respectively, because shockwave speeds are the slopes of lines connecting points on the fundamental diagram. This leads to an important notion, the *domain of dependence*. Consider a point (x, t) in space and time. Through this point, draw lines with slopes $Q'(0)$ and $Q'(k_j)$. The area between

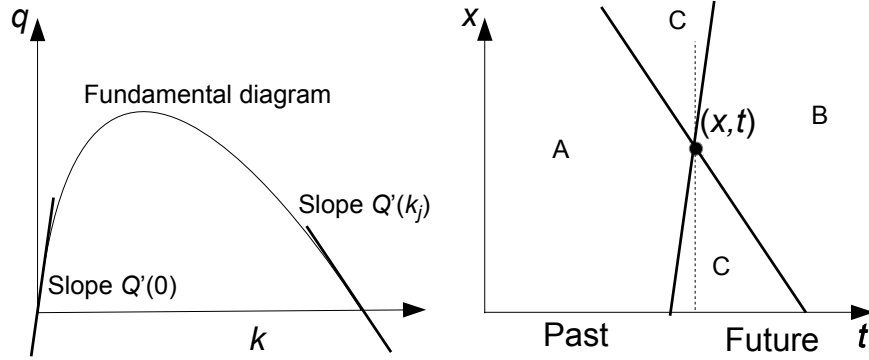


Figure 10.19: The fundamental diagram (left) and the domain of dependence (right).

these lines to the left of (x, t) represents all the points which can potentially influence the traffic state at (x, t) (labeled as region A in Figure 10.19), and the area between the lines to the right of (x, t) represents all of the points the traffic state at (x, t) can potentially influence (labeled as region B). In the LWR model, points outside of these regions (the regions labeled C) are independent of what happens at (x, t) (say, a signal turning red or green): in the past, they are either too recent or too distant to affect what is happening at (x, t) . In the future, they are too soon or too distant to be affected by an event at (x, t) . *This is a crucial fact for dynamic network loading.* If you want to know what is happening at (x, t) , it is sufficient to know what has happened in the past, in the domain of dependence. We do not need to know what is happening simultaneously at other points in the network, and as a result we can perform network loading in an decentralized fashion, performing calculations in parallel since they do not depend on each other.⁶

10.4.4 Characteristics and the Newell-Daganzo method

By substituting equation (10.32) into (10.33), and setting aside (10.31) for the moment, we can obtain a partial differential equation in k alone:

$$\frac{dQ}{dk} \frac{\partial k}{\partial x} + \frac{\partial k}{\partial t} = 0 \quad (10.35)$$

It is in this form that the LWR model is most frequently solved, obtaining $k(x, t)$ values everywhere. This type of partial differential equation can be approached using the *method of characteristics*, which is briefly described below.

A *characteristic* is a curve in (x, t) space, along which the density $k(x, t)$ varies in a predictable way. If the fundamental diagram Q does not vary in space

⁶If you have studied relativistic physics, there are many similarities with the notion of light cones and causality.

and time, then straight lines form characteristics of (10.35). To demonstrate, assume that we know the density $k(x_1, t_1)$ at some point, and consider a straight line C through (x_1, t_1) with some slope v . Then the directional derivative of k along this line is given by

$$k' = \nabla k \cdot [v \ 1] = \frac{\partial k}{\partial x} v + \frac{\partial k}{\partial t}. \quad (10.36)$$

Substituting the conservation equation (10.35) and rearranging, we have

$$k' = \left(v - \frac{dQ}{dk} \right) \frac{\partial k}{\partial x}. \quad (10.37)$$

This gives us a formula for the change in density along the line. But what is really useful is noticing that if $v = \frac{dQ}{dk}$, then (10.37) vanishes and *the density is constant along the line*. Put another way, the density is constant along any line whose slope is equal to the slope of the *tangent* of the fundamental diagram at that density value.⁷

To see this, return to the example we have been using where $Q(k) = \frac{1}{4}k(240 - k)$ and $N(x, t) = 60t - 120x + \frac{60x^2}{t+1}$. We have calculated the density and flow at all points and times in equations (10.28) and (10.29). In our calculation, we observed that the line integral along path C was the simplest to evaluate, because $k(x, t)$ was a constant 60 vehicles per mile along the line $x/(t+1) = 1/2$. This path is in fact a characteristic: the derivative of the fundamental diagram at this point is $Q'(60) = 60 - k/2 = 30$ miles per hour. Rearranging the equation of the line, we have $x = t/2 + 1/2$; in other words, this is a line whose location moves half a mile in one minute: 30 miles per hour, the same as the derivative of the fundamental diagram.

These slopes will be different at points where the density is different, which means that characteristic lines can potentially intersect. This indicates the presence of a shockwave separating regions of different density. One interpretation of these characteristic lines is as “directions of influence,” since the density at a point will determine the density at all later points along this line. In uncongested regions, where k is subcritical, dQ/dk is positive, meaning that the characteristics have positive slope: uncongested states will propagate downstream. In congested regions with supercritical density, dQ/dk is negative, and the characteristics have negative slope: congested states will propagate upstream.

Furthermore, by combining knowledge of characteristics with equation (10.26), we can determine the cumulative count and the density at any point, given sufficient initial and boundary data. This method was first suggested by G. F. Newell, and developed further by Carlos Daganzo. Suppose we wish to calculate the density at a point (x, t) . If we knew the density at this point, then we would know the slope of the characteristic through this point. This characteristic could be traced back until it intersected a point (x_0, t_0) where the cumulative count

⁷Note that the speed of the characteristic is different from the speed of the vehicles themselves, or the speed of shockwaves (which are given by slopes of secant lines on the fundamental diagram).

N was known, either because it corresponds to an initial or boundary point, or a point where N has already been calculated. Then, applying (10.26), we would have

$$N(x, t) = N(x_0, t_0) + \int_C q \, dt - k \, dx \quad (10.38)$$

where C is the straight line between (x_0, t_0) and (x, t) . Since this line is a characteristic, we have $dx/dt = dQ/dk$, and so

$$N(x, t) = N(x_0, t_0) + \int_C \left(q - k \frac{dQ}{dk} \right) dt. \quad (10.39)$$

Furthermore, as a characteristic, k (and therefore q) are constant. So the integral is easy to evaluate, and we have

$$N(x, t) = N(x_0, t_0) + \left(q - k \frac{dQ}{dk} \right) (t - t_0). \quad (10.40)$$

The only trouble is that we do not actually know the density at (x, t) . Each possible value of density corresponds to a slightly different cumulative count, based on equation (10.40). The insight of the Newell-Daganzo method is that *the correct value of the cumulative count is the lowest possible value*. That is, imagine that the density at (x, t) is k , and let (x_k, t_k) be the known point corresponding to the characteristic slope of density k . Then

$$N(x, t) = \inf_{k \in [0, k_j]} \left\{ N(x_k, t_k) + \left(q - k \frac{dQ}{dk} \right) (t - t_k) \right\}. \quad (10.41)$$

Rigorously validating this insight requires knowledge of the calculus of variations, which is beyond the scope of this text. An intuitive justification is that (10.40) represents an upper bound on the cumulative count imposed by a boundary or initial point — we know the number of the vehicle passing that boundary or initial point, and (10.40) expresses one possible value of the cumulative count at the point in question, if the density took a particular value and no shockwave intervened. However, due to the possible presence of shockwaves, there may be another, more restrictive constraint imposed on the cumulative count. The equation (10.41) thus finds the “most restrictive” boundary or initial condition, which gives the correct cumulative count.

Better yet, this method becomes exceptionally easy if we assume that the fundamental diagram takes a simple form, such as a triangular shape (Figure 10.20), where the equation is given by

$$Q(k) = \min \{ u_f k, w(k_j - k) \}. \quad (10.42)$$

In this case, there are only two possible characteristic speeds: one $(+u_f)$ corresponding to uncongested conditions, and the other $(-w)$ corresponding to congested conditions. The uncongested speed u_f is the free-flow speed, and w is known as the *backward wave speed*. In the case where the characteristic speed

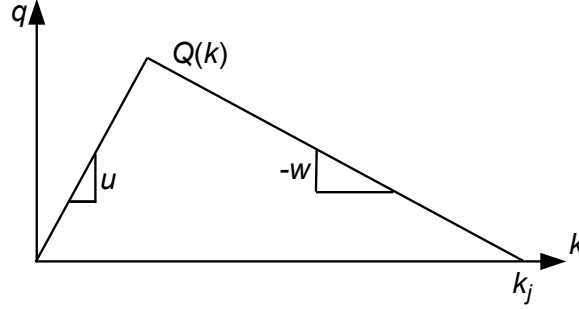


Figure 10.20: A triangular fundamental diagram, defined by three parameters u_f , w , and k_j .

is u_f , the vehicle speed u equals the characteristic speed u_f , and both are equal to q/k . Therefore, the line integral along the characteristic in (10.40) is

$$(q - ku_f)(t - t_0) = (q - uk)(t - t_0) = 0 \quad (10.43)$$

because $q = uk$. In other words, *the vehicle number is constant along characteristics at free-flow speed.*

For the congested characteristic with slope $-w$, we have

$$(q - k(-w))(t - t_0) = w \left(\frac{q}{w} + k \right) (t - t_0) = k_j w (t - t_0) \quad (10.44)$$

since $k + \frac{q}{w} = k_j$, as can be seen in Figure 10.20. This expression can also be written as $k_j(x_0 - x)$. Since these are the only two characteristics which can prevail at any point, equation (10.41) gives

$$N(x, t) = \min \{N(x_U, t_U), N(x_C, t_C) + k_j(x_C - x)\} \quad (10.45)$$

where (x_U, t_U) is the known point intersected by the uncongested characteristic, and (x_C, t_C) is the known point intersected by the congested characteristic.

Trapezoidal fundamental diagrams, such as that in Figure 10.21, are also commonly used, given by the equation

$$Q(k) = \min \{u_f k, q_{max}, w(k_j - k)\} . \quad (10.46)$$

In this case, there is a third possible characteristic speed, corresponding to the flat region where flow is at capacity. Since the derivative of the fundamental diagram at this point is zero, this characteristic has zero speed, represented by a horizontal line on a space-time diagram. Tracing this characteristic back to a known point (x, t_R) , the change in vehicle number between this known point and the unknown point (x, t) is just

$$q_{max}(t - t_R) \quad (10.47)$$

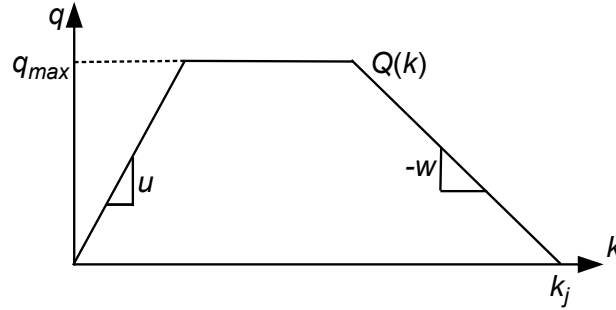


Figure 10.21: A trapezoidal fundamental diagram, defined by four parameters u_f , q_{max} , w , and k_j .

since $dx = 0$ in the line integral (10.26). (Note that the location of this third known point is the same as the location of the point we are solving for, since the characteristic is stationary.) This adds a third term to the minimum in (10.45), giving

$$N(x, t) = \min \{N(x_U, t_U), N(x, t_R) + q_{max}(t - t_R), N(x_C, t_C) + k_j(x_C - x)\} \quad (10.48)$$

for trapezoidal fundamental diagrams.

The trapezoidal fundamental diagram requires four parameters to calibrate: the free-flow speed u_f , the capacity q_{max} , the jam density k_j , and the backward wave speed $-w$. The first three of these are fairly straightforward to estimate from traffic engineering principles. The backward wave speed is a bit trickier; empirically it is often a third to a half of the free-flow speed.

As a demonstration of this method, consider a link which is 1 mile long. The fundamental diagram is shown in the left panel of Figure 10.22, and has the equation

$$Q(k) = \min \left\{ k, 120 - \frac{1}{2}k \right\} \quad (10.49)$$

when flow is measured in vehicles per minute, and density in vehicles per mile. Initially, vehicles on the link flow at an uncongested 48 veh/min — this state has existed for a long time in the past, and vehicles continue entering the link at this rate. However, at the downstream end of the link there is an obstruction which prevents any vehicles from passing, such as a red light or an incident blocking all lanes. This will cause a queue of stopped vehicles to form. Assume that we want to know how many vehicles lie between the obstruction and three given points: half a mile upstream of the obstruction, 30 seconds after it begins; an eighth of a mile upstream of the obstruction, 30 seconds after it begins; and an eighth of a mile upstream of the obstruction, one minute after it begins. We do not know, *a priori*, whether these points lie within the queue or in the portion of the link which is still uncongested.

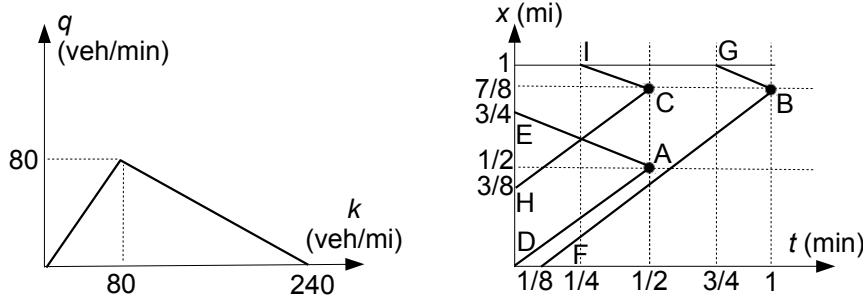


Figure 10.22: Example of the Newell-Daganzo method.

The first step is to establish a coordinate system. As always, we set $x = 0$ at the upstream end of the link. For this problem, it will be convenient to set $t = 0$ at the time when the obstruction begins, and to count vehicles starting from the first vehicle stopped at the obstruction, that is, $N(1, 0) = 0$. This way, $N(x, t)$ will immediately give the number of vehicles between point x and the obstruction at time t . Next, we use the given data from the problem to construct initial and boundary conditions where we already know $N(x, t)$. Since no vehicles can pass the obstruction, we know that $N(1, t) = 0$ for all t . Since the link is initially uncongested at a flow rate of 48 veh/min, the fundamental diagram (10.49) gives the initial density to be 48 veh/mi. Therefore, the initial condition is $N(x, 0) = 48 - 48x$, and the vehicle number at the origin of the coordinate system is 48. Since vehicles continue to enter the link at a rate of 48 veh/min, we have $N(0, t) = 48 + 48t$ along the upstream boundary of the link. The three points where we must calculate $N(x, t)$ are labeled as A, B, and C in the right panel of Figure 10.22.

We start with point A, half a mile upstream of the obstruction and 30 seconds after it begins. There are two possible characteristics at this point, one with slope $+1$ (corresponding to uncongested conditions) and one with slope $-\frac{1}{2}$ (corresponding to congested conditions). We can trace back these characteristics until they reach a point where $N(x, t)$ is known, in this case an initial or boundary condition. These points of intersection are labeled D and E in Figure 10.22. From the initial condition, we know that $N(D) = 48$ and $N(E) = 12$. Along the uncongested characteristic, there is no change in the cumulative count, while along the congested characteristic the cumulative count increases at a rate of $k_j = 240$ veh/mi for each mile traveled. Therefore, equation (10.43) tells us that $N(A) = 48 + 0 = 48$ if point A is uncongested, while equation (10.44) tells us that $N(A) = 12 + \frac{1}{4}240 = 72$ if point A is congested. The correct value is the smaller of the two: $N(A) = 48$, and this point is uncongested (the queue has not yet reached this point).

We next move to point B, an eighth of a mile upstream of the obstruction and 60 seconds after it begins. Tracing back the two possible characteristics from point B leads us to the points labeled F and G, and from the boundary

conditions we know $N(F) = 54$ and $N(G) = 0$. So $N(B)$ is the lesser of $N(F) + 0 = 54$ and $N(G) + \frac{1}{8}240 = 30$. This means that there are 30 vehicles between point B and the obstruction, and since the congested characteristic produced the lower value, point B lies within the queue.

At point C, the two characteristics lead to the points labeled H and I, where $N(H) = 30$ from an initial condition and $N(I) = 0$ from a boundary condition. So $N(C) = \min\{30 + 0, 0 + \frac{1}{8}240\} = 30$. In this case, both characteristics led to the *same* value of $N(C)$. This indicates that the shockwave passes exactly through point C.

An alternative approach for this problem would be to explicitly calculate the location of the shockwave, determine the densities in each region, and apply (10.26) directly. For this problem, this approach would be simpler than the Newell-Daganzo method. However, if there were multiple shockwaves introduced into the problem (say, from multiple red/green cycles), it would become very tedious to track the locations of all of the shockwaves and determine which region the points lie in. The Newell-Daganzo method can be applied just as easily in such a case, once the boundary conditions are determined.

10.5 LWR-based Link Models

The hydrodynamic traffic flow model developed by Lighthill, Whitham, and Richards forms the basis for several popular link models. The LWR model is simple enough to be usable in large-scale dynamic network loading, while capturing enough key properties of traffic flow for its results to be meaningful. Through shockwaves, we can capture how congestion grows and shrinks over time. These shockwaves allow us to model queue spillback (when a congestion shockwave reaches the upstream end of a link) and to account for delays in queue startup (unlike the spatial queue model). This section describes two link models based on the LWR model — the cell transmission model, which is essentially an explicit solution scheme for the LWR system of partial differential equations, and the link transmission model, which uses the Newell-Daganzo method to directly calculate sending and receiving flows. Lastly, we show how the point and spatial queue models can be seen as special cases of LWR-based link models, with a suitable choice of the fundamental diagram.

10.5.1 Cell transmission model

In the cell transmission model, in addition to discretizing time into intervals of length Δt , we also discretize space, dividing the link into cells of length Δx . These two discretizations are not chosen independently. Rather, they are related by

$$\Delta x = u_f \Delta t, \quad (10.50)$$

that is, the length of each cell is the distance a vehicle would travel in time Δt at free flow. The reasons for this choice are discussed at the end of this section.

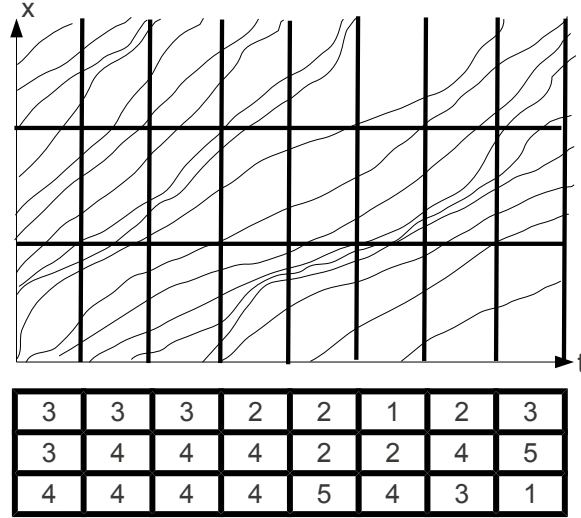


Figure 10.23: Discretizing space and time into cells, $n(x, t)$ values below.

With this discretization in mind, we will use the notation $n(x, t)$ to describe the number of vehicles in cell x at time t , where x and t are both integers — we must convert the continuous LWR variables k and q into discrete variables for dynamic network loading, which we will call n and y . If the cell size is small, we can make the approximation that

$$n(x, t) \approx k(x, t)\Delta x, \quad (10.51)$$

essentially assuming that the density within the cell is constant. Further define $y(x, t)$ to be the number of vehicles which *enter* cell x during the t -th time interval. Making a similar assumption, we can make the approximation

$$y(x, t) \approx q(x, t)\Delta t. \quad (10.52)$$

In a space-time diagram showing vehicle trajectories, such as Figure 10.23, $n(x, t)$ and $y(x, t)$ respectively correspond to the number of trajectories crossing the vertical line at t between locations x and $x + \Delta x$, and the number of trajectories crossing the horizontal line at x between times t and $t + \Delta t$.

The cell transmission model provides methods for solving for $n(x, t)$ for all integer values of x and t ; these can then be converted to density values $k(x, t)$ through (10.51). These density values can then be used to calculate flows $q(x, t)$ through the fundamental diagram, and $u(x, t)$ values through equation (10.20), finally providing an approximate solution to the system of partial differential equations (10.31)–(10.33). Recall that the goal of a link model is to determine

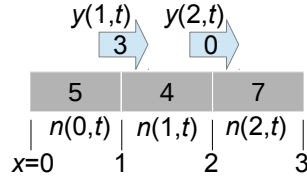


Figure 10.24: Where discrete values are calculated in the cell transmission model.

sending and receiving flow at each time step. For this reason, we will be content with determining how the $n(x, t + \Delta t)$ values can be calculated, given the $n(x, t)$ values (which are already known), and the $y(x, t)$ values, which must be calculated.

Since $q(x, t) = Q(k(x, t))$, substitution into equations (10.51) and (10.52) give

$$y(x, t) \approx Q\left(\frac{n(x, t)}{\Delta x}\right) \Delta t. \quad (10.53)$$

Substituting the particular form of the fundamental diagram gives an equation for $y(x, t)$ in terms of $n(x, t)$. Using the trapezoidal diagram of Figure 10.21, we have

$$y(x, t) \approx \min \left\{ u_f n(x, t) \frac{\Delta t}{\Delta x}, q_{max} \Delta t, w \Delta t \left(k_j - \frac{n(x, t)}{\Delta x} \right) \right\}. \quad (10.54)$$

Using the fact that $\Delta x / \Delta t = u_f$, the first term in the minimum is simply $n(x, t)$. The third term can be simplified by defining $\bar{n}(x) = k_j \Delta x$ to be the maximum number of vehicles which can fit into a cell and $\delta = w / u_f$ to be the ratio between the backward wave speed and free-flow speed. Then, factoring out $1 / \Delta x$ from the term in parentheses and again using $\Delta x / \Delta t = u_f$, the third term simplifies to $\delta(\bar{n}(x) - n(x, t))$.

There is one more point which is subtle, yet incredibly important. Being a “flow” variable, $y(x, t)$ is calculated at a single point (over a time interval), while $n(x, t)$ is calculated at a single time (over a longer spatial interval). As shown in Figure 10.24, the x in $y(x, t)$ refers to a single location, while the x in $n(x, t)$ refers to an entire cell. So, when we are calculating the flow across the (single) point x , which is the boundary between two cells, do we look at the adjacent cell upstream $n(x - 1, t)$, or the adjacent cell downstream $n(x, t)$?

The correct answer depends on the fundamental diagram, and the meaning of characteristics. In uncongested conditions, corresponding to the increasing part of the fundamental diagram and the first term in the minimum, the traffic state moves from upstream to downstream (because the characteristic has positive speed). In congested conditions, corresponding to the decreasing part of the fundamental diagram and the third term in the minimum, the traffic state moves from downstream to upstream (because the characteristic has negative speed.)

So, if traffic is uncongested at the (single) point x , we need to refer to the upstream cell, while if traffic is congested at the (single) point x , we must refer to the downstream cell. So, the final expression for the cell transmission model flows is

$$y(x, t) = \min\{n(x - \Delta x, t), q_{max}\Delta t, \delta(\bar{n}(x) - n(x, t))\}. \quad (10.55)$$

This expression also has a nice intuitive interpretation. The number of vehicles moving from cell $x - 1$ to cell x is limited either by the number of vehicles in the upstream cell (the first term), the capacity of the roadway (the second term), or by the available space in the downstream cell (the third term).

Expression (10.55) is the discrete equivalent of the differential equation (10.32). We now derive the discrete form of the partial differential equation giving the conservation law (10.33). The derivative $\frac{\partial q}{\partial x}$ can be approximated as

$$\frac{\partial q}{\partial x}(x, t) \approx \frac{1}{\Delta t \Delta x} (y(x + \Delta x, t) - y(x, t)) \quad (10.56)$$

and the derivative $\frac{\partial k}{\partial t}$ can be approximated as

$$\frac{\partial k}{\partial t}(x, t) \approx \frac{1}{\Delta t \Delta x} (n(x, t + \Delta t) - n(x, t)). \quad (10.57)$$

Substituting into (10.33), we have

$$\frac{1}{\Delta t \Delta x} (y(x + \Delta x, t) - y(x, t) + n(x, t + \Delta t) - n(x, t)) = 0 \quad (10.58)$$

or, in a more convenient form,

$$n(x, t + \Delta t) = n(x, t) + y(x, t) - y(x + \Delta x, t) \quad (10.59)$$

This also has a simple intuitive interpretation: the number of vehicles in cell x at time $t + 1$ is simply the number of vehicles in cell x at the previous time t , plus the number of vehicles which flowed into the cell during the t -th time interval, minus the number of vehicles which left.

Together, the equations (10.55) and (10.59) define the cell transmission model for trapezoidal fundamental diagrams. There are only two pieces of “missing” information, at the boundaries of the link. Refer again to Figure 10.24. How should $y(0, t)$ and $y(L, t)$ be calculated? For $y(0, t)$, the first term in formula (10.55) involves $n(-\Delta x, t)$, while for $y(L, t)$, the third term in the formula involves $n(L + \Delta x, t)$, and both of these cells are “out of range.” The answer is that these boundary flows are used to calculate the sending and receiving flows for the link, and a node model will then give the actual values link inflows $y(0, t)$ and link outflows $y(L, t)$.

Remember that the sending flow is the maximum number of vehicles which could leave the link if there was no obstruction from downstream. In terms of (10.55), this means that the third term in the minimum (which corresponds to downstream congestion) is ignored. Then, the first two terms in the minimum

(which only refer to cells on the link) are the possible restrictions on the flow leaving the link, so

$$S(t) = \min\{n(C, t), q_{max}\Delta t\}. \quad (10.60)$$

Likewise, the receiving flow is the maximum number of vehicles which could enter the link if there were a large number of vehicles wanting to enter from upstream. In terms of (10.55), this means that the first term in the minimum (which corresponds to the number of vehicles wanting to enter) is ignored. The second two terms in the minimum refer to cells on the link, and

$$R(t) = \min\{q_{max}\Delta t, \delta(\bar{n}(0) - n(0, t))\}. \quad (10.61)$$

It remains to explain the choice of the discretization (10.50). An intuitive explanation for linking the cell length and time discretizations in this way is that this choice limits vehicles to moving at most one cell between time steps. In fact, in uncongested conditions, *all* vehicles in a cell will move to the next cell downstream in the next time interval, simplifying calculations — for instance, the simplification in (10.54) only works because of the choice made in (10.50). The underlying mathematical reason has to do with the speed of characteristics, which for the trapezoidal fundamental diagram lie between $-w$ and $+u_f$. Empirically, $w < u_f$, so the fastest moving characteristic (in either direction) is one with speed u_f . More important than vehicles moving at most one cell between time steps is that *characteristics cannot move more than one cell between time steps*, in either the upstream or downstream directions. This condition is needed for stability of finite-difference approximations to partial differential equations, and corresponds to the *Courant-Friedrich-Lewy (CFL) condition* for explicit solution methods such as the cell transmission model.

Table 10.4 provides a demonstration of how the cell transmission model works. In this example, a link is divided into three cells, implying that it takes three time steps for a vehicle to traverse the link at free-flow. The fundamental diagram is such that at most 10 vehicles can move between cells in one time step, at most 30 vehicles can fit into one cell at jam density, and the ratio between the backward and forward characteristic speeds is $\delta = 2/3$. There is a red light at the downstream end of the link which turns green at $t = 10$, and remains green thereafter. In this table, $d(t)$ represents the number of vehicles that wish to enter the link during the t -th timestep (perhaps the sending flow from an upstream link), while $R(t)$ is the receiving flow for the link, calculated from (10.61). The middle columns of the table show the main cell transmission model calculations: the number of vehicles in each cell at the start of each timestep, $N(x, t)$, and the number of vehicles moving into each cell during the t -th time step. These values are calculated from (10.59) and (10.55), respectively, along with the initial condition that the link is empty, that is, $N(x, 0) = 0$ for all x . The rightmost columns of the table show the link's sending flow, calculated from (10.60), and the actual flow which leaves the link, denoted $y(3, t)$. This latter value is constrained to be zero as long as the light at the downstream end of the link is red.

Notice that the table has non-integer values: we do not need to round cell occupancies and flows to whole values, since the LWR model assumes vehicles are a continuously-divisible fluid. Preserving non-integer values also ensures that the cell transmission model remains accurate no matter how small the timestep Δt is (in fact, its accuracy should increase as this happens). Insisting that flows and occupancies be rounded to whole numbers can introduce significant error if the timestep is small, unless one is careful with implementation.

Table 10.5 shows only the cell occupancies at each timestep, color-coded according to the density in the cells. In this example, the link is initially at free-flow, until the first vehicles encounter the red light and must stop. A queue forms, and a shockwave begins moving backward. When this shockwave reaches the cell at the upstream end of the link, the receiving flow of the link decreases, and the inflow to the link is limited. When the light turns green, a second shockwave begins moving backward as the queue clears. Once this shockwave overtakes the first, vehicles can begin entering the link again. For a few time steps, the inflow is greater than $d(t)$, representing demand which was blocked when the receiving flow was restricted by the queue and which was itself queued on an upstream link (the “queue spillback” phenomenon). Unlike the point queue and spatial queue link models, the cell transmission model tells us what is happening in the interior of a link, not just at the endpoints. This is both a blessing and a curse: sometimes this additional information is helpful, while other times we may not be concerned with such details. The link transmission model, described next, can simplify computations if we do not need information on the internal state of a link.

10.5.2 Link transmission model

The link transmission model allows us to calculate sending and receiving flows for links with any trapezoidal fundamental diagram.⁸ In contrast to the cell transmission model, it only involves calculations at the ends of the links — details of what happen in the middle of the link are ignored. As a consequence, the link transmission model does not require us to keep track of information within the link. However, it does require us to keep track of information on the past state of the link, whereas the cell transmission model calculations only involve quantities at the current time step, and the previous time step. The link transmission model can also overcome the “shock spreading” phenomenon, where backward-moving shockwaves in the cell transmission model can diffuse across multiple cells, even though they are crisp in the LWR model. (See Exercise 29.)

Assume that the trapezoidal fundamental diagram is parameterized as in Figure 10.21. There are three characteristic speeds, $+u_f$ at free-flow, 0 at capacity flow, and $-w$ at congested flow. As with the point queue and spatial queue models, we apply the Newell-Daganzo method to calculate sending and receiving flows. We will only need to refer to cumulative counts N at the

⁸In fact, it can be generalized to any piecewise linear fundamental diagram without too much difficulty; see Exercise 34.

Table 10.5: Cell occupancies from the example in Table 10.4, green is zero density and red is jam density.

t	$N(0,t)$	$N(1,t)$	$N(2,t)$
0	0	0	0
1	10	0	0
2	10	10	0
3	10	10	10
4	10	10	20
5	10	13.3	26.7
6	9	21.1	28.9
7	11.1	26.3	29.6
8	15.6	28.5	29.9
9	20.6	29.4	30
10	25.2	29.8	30
11	28.3	29.9	20
12	29.4	23.3	16.7
13	25.3	18.9	15.6
14	21	16.7	15.2
15	13.4	15.7	15.1
16	3.9	15.3	15
17	0	9.2	15
18	0	0	14.2
19	0	0	4.2
20	0	0	0

upstream and downstream ends of the link, that is, at $N(0, \cdot)$ and $N(L, \cdot)$, respectively. In keeping with the notation introduced in Section 10.1, we will refer to these as N^\uparrow and N^\downarrow . If there is no obstruction from a downstream link or node, then the end of the link will be uncongested, and the characteristic at this end will either have slope $+u_f$ or slope zero. The Newell-Daganzo method thus gives

$$N^\downarrow(t + \Delta t) = \min\{N^\uparrow((t + \Delta t) - L/u_f), N^\downarrow(t) + q_{max}\Delta t\} \quad (10.62)$$

and the equation for the sending flow is obtained as the difference between $N^\downarrow(t + 1)$ and $N^\downarrow(t)$:

$$S(t) = \min\{N^\uparrow((t + \Delta t) - L/u_f) - N^\downarrow(t), q_{max}\Delta t\}. \quad (10.63)$$

For the receiving flow, we have to take into account the two relevant characteristic speeds of 0 and $-w$, since the receiving flow is calculated assuming an inflow large enough that the upstream end of the link is congested (or at least at capacity). The stationary characteristic corresponds to the known point $N(0, t)$, while the backward-moving characteristic corresponds to the known point $N(L, t - L/w)$. Thus, applying the last two terms of equation (10.46) would give

$$N^\uparrow(t + \Delta t) = \min\{N^\uparrow(t) + q_{max}\Delta t, N^\downarrow((t + \Delta t) - L/w) + k_j L\} \quad (10.64)$$

and

$$R(t) = \min\{q_{max}\Delta t, (N^\downarrow((t + \Delta t) - L/w) + k_j L) - N^\uparrow(t)\}. \quad (10.65)$$

It is possible to show that equations (10.63) and (10.65) ensure that the number of vehicles on the link is always nonnegative, and less than $k_j L$.

The link transmission model is demonstrated on an example similar to the one used for the cell transmission model; the only difference is that the ratio of backward-to-forward characteristics has been adjusted from 2/3 to 3/4. In particular, $L/u_f = 3\Delta t$, and $L/w = 4\Delta t$, so forward-moving characteristics require three time steps to cross the link, and backward-moving characteristics require four time steps. The total number of vehicles which can fit on the link is $k_j L = 90$. Otherwise, the example is the same: the demand profile is identical, and a red light prevents outflow from the link until $t = 10$. The results of the calculations are shown in Table 10.6. The rightmost column shows the number of vehicles on the link, which is the difference between the upstream and downstream cumulative counts at any point in time. Notice that inflow to the link is completely blocked during the 12th and 13th time intervals, even though the number of vehicles on the link is less than the jam density of 90. This happens because the queue has started to clear at the downstream end, but the clearing shockwave has not yet reached the upstream end of the link. The vehicles at the upstream end are still stopped, and no more vehicles can enter. In contrast, the spatial queue model would allow vehicles to start entering the link as soon as they began to leave.

Table 10.6: Link transmission model example, with $L/u_f = 3\Delta t$, $L/w = 4\Delta t$, and $k_j L = 90$.

t	$d(t)$	$R(t)$	Inflow	$N^\uparrow(t)$	$N^\downarrow(t)$	$S(t)$	Outflow	Vehicles on link
0	10	10	10	0	0	0	0	0
1	10	10	10	10	0	0	0	10
2	10	10	10	20	0	0	0	20
3	10	10	10	30	0	10	0	30
4	10	10	10	40	0	10	0	40
5	9	10	9	50	0	10	0	50
6	8	10	8	59	0	10	0	59
7	7	10	7	67	0	10	0	67
8	6	10	6	74	0	10	0	74
9	5	10	5	80	0	10	0	80
10	4	5	4	85	0	10	10	85
11	3	1	1	89	10	10	10	79
12	2	0	0	90	20	10	10	70
13	1	0	0	90	30	10	10	60
14	0	10	5	90	40	10	10	50
15	0	10	0	95	50	10	10	45
16	0	10	0	95	60	10	10	35
17	0	10	0	95	70	10	10	25
18	0	10	0	95	80	10	10	15
19	0	10	0	95	90	5	5	5
20	0	10	0	95	95	0	0	0

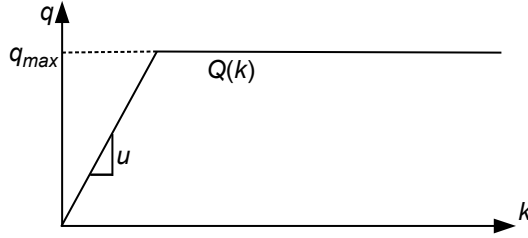


Figure 10.25: Fundamental diagram in the point queue model.

10.5.3 Point and spatial queues and the LWR model (*)

(This optional section shows how the previously-introduced point and spatial queue models can be seen as special cases of the LWR model.)

The first link models introduced in this chapter were the point and spatial queue models, in Section 10.1. These were presented as simple link models to illustrate concepts like the sending and receiving flow, rather than realistic depictions of traffic flow. Nevertheless, it is possible to view the point and spatial queue models as special cases of the LWR model by making an appropriate choice of the fundamental diagram, as shown in this section. Applying the Newell-Daganzo method with this fundamental diagram gives us a second way to derive the expressions for sending and receiving flow for these models.

The point queue model is equivalent to assuming that the flow-density relationship is as shown in Figure 10.25. This diagram is unlike others we've seen, because there is no jam density. This represents the idea that the point queue occupies no physical space: no matter how many vehicles are in queue, there is nothing to prevent additional vehicles from entering the link and joining the queue. It is also the simplest diagram which we have seen so far, and is defined by only two parameters: the free-flow speed u_f and the capacity q_{max} . (Even the triangular fundamental diagram in Figure 10.20 required a third parameter, either $-w$ or k_j .)

The Newell-Daganzo method leads to a simple expression for the sending and receiving flows in a point queue model. To calculate the sending flow $S(t)$, we need to examine the downstream end of the link, so $x = L$. Since we are solving in increasing order of time, we already know $N^\downarrow(0), N^\downarrow(\Delta t), \dots, N^\downarrow(t)$ (the number of vehicles which have left the link at each time interval). Likewise, we know how many vehicles have entered the link at earlier points in time, so we know $N^\uparrow(0), N^\uparrow(\Delta t), \dots, N^\uparrow(t)$. For the sending flow, we are assuming that there are no obstructions from downstream. If this were the case, then we can calculate $N^\downarrow(t + \Delta t)$ using the Newell-Daganzo method, and

$$S(t) = N^\downarrow(t + \Delta t) - N^\downarrow(t). \quad (10.66)$$

In the point queue model, there are two possible wave speeds, u_f (corre-

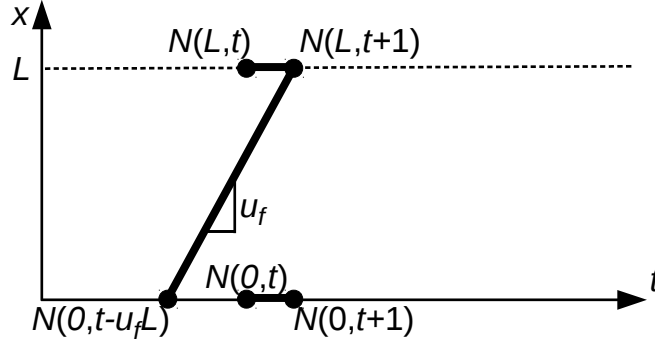


Figure 10.26: Point queue model characteristics for sending and receiving flow.

sponding to free-flow conditions) and zero (corresponding to flow moving at capacity). Figure 10.26 shows how these characteristics can be traced back to known data, either at the upstream end $x = 0$, or at the downstream end $x = L$. Therefore, $t_R = t$, $x_U = 0$, and $t_U = (t + \Delta t) - L/u_f$. We can think of this as a special case of a “trapezoidal” diagram where the jam density k_j is infinite. Applying equation (10.48), we see that if $k_j = \infty$ then the minimum must occur in one of the first two terms, so

$$N^\downarrow(t + \Delta t) = \min\{N^\uparrow((t + \Delta t) - L/u_f), N^\downarrow(t) + q_{max}\Delta t\} \quad (10.67)$$

and

$$S(t) = \min\{N^\uparrow((t + \Delta t) - L/u_f) - N^\downarrow(t), q_{max}\Delta t\}. \quad (10.68)$$

The two terms in the minimum in (10.68) correspond to the case when the queue is empty, and when there are vehicles in queue. In the first term, since there is no queue, we just need to know how many vehicles will finish traversing the physical section of the link between t and $t + \Delta t$; this is exactly the difference between the total number of vehicles which have entered by time $t + \Delta t - L/u_f$ and the total number that have left by time t . When there is a queue, the vehicles exit the link at the full capacity rate.

In these expressions, it is possible that $(t + \Delta t) - L/u_f$ is not an integer, that is, it does not line up with one of the discretization points exactly. In this case the most accurate choice is to interpolate between the known time points on either side (remember that we chose Δt so that $L/u_f \geq 1$). If you are willing to sacrifice some accuracy for efficiency, you can choose to round to the nearest integer, or to adjust the length of the link so that L/u_f is an integer.

For the receiving flow, we look at the upstream end of the link. We can treat the same points as known — $N^\uparrow(0)$, $N^\uparrow(\Delta t)$, \dots , $N^\uparrow(t)$ and $N^\downarrow(0)$, $N^\downarrow(\Delta t)$, \dots , $N^\downarrow(t)$. *Since the fundamental diagram for the point queue model has no decreasing portions, the known data at the downstream end can never be relevant.* (A line con-

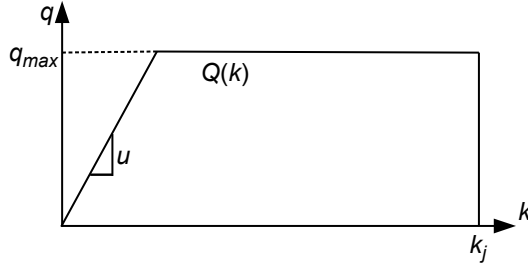


Figure 10.27: Fundamental diagram in the spatial queue model.

necting one of these points to the unknown point $N^\uparrow(t + \Delta t)$ must have negative slope, see Figure 10.26.) Furthermore, for the receiving flow, the characteristic with positive slope $+u_f$, corresponding to upstream conditions, is irrelevant because we are assuming an unlimited number of vehicles are available to move from upstream — and therefore its term in (10.48) will never be the minimum. We are only left with the middle term, corresponding to capacity, so

$$N^\uparrow(t + \Delta t) = N^\uparrow(t) + q_{max}\Delta t \quad (10.69)$$

and

$$R(t) = q_{max}\Delta t. \quad (10.70)$$

In terms of the fundamental diagram, the spatial queue model takes the form in Figure 10.27. This diagram requires three parameters to calibrate: the free-flow speed u_f , the capacity q_{max} , and the jam density k_j . Notice, however, that the fundamental diagram is discontinuous, and immediately drops from q_{max} to zero once jam density is reached. This implies that backward-moving shockwaves can have infinite speed in the spatial queue model — a physical interpretation is that when vehicles at the front of the queue begin moving, vehicles at the rear of the queue immediately start moving as well. In reality, there is a delay before vehicles at the rear of the queue begin moving, and this can be treated as an artifact arising from simplifying assumptions made in the spatial queue model.⁹

There are thus three possible characteristic speeds: $+u_f$ at free-flow, 0 at capacity flow, and $-\infty$ when the queue reaches jam density. The Newell-Daganzo method is applied in much the same way as was done for the point queue model. In particular, the sending flow expression is exactly the same, because the two characteristics with nonnegative velocity are the same. We thus have

$$N(L, t + 1) = \min\{N^\uparrow((t + \Delta t) - L/u_f), N^\downarrow(t) + q_{max}\Delta t\} \quad (10.71)$$

and

$$S(t) = \min\{N^\uparrow((t + \Delta t) - L/u_f) - N^\downarrow(t), q_{max}\Delta t\}. \quad (10.72)$$

⁹Alternatively, connected and autonomous vehicles may be able to exhibit such behavior if an entire platoon of vehicles coordinates its acceleration.

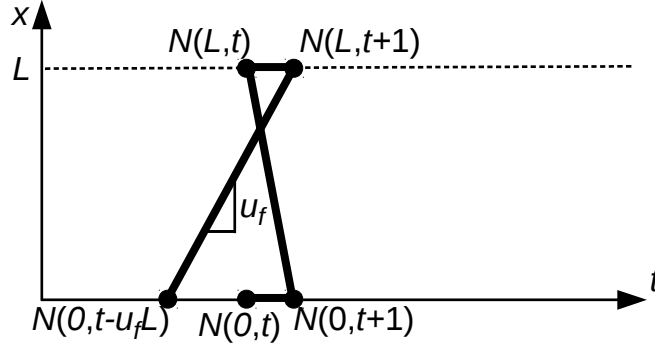


Figure 10.28: Spatial queue model characteristics for sending and receiving flow; the upstream-moving wave is an approximation of the vertical component of the fundamental diagram.

For the receiving flow, we have to take into account the new shockwave speed. Dealing with an infinite speed can be tricky, since, taken literally, would mean that the upstream cumulative count $N^\uparrow(t + \Delta t)$ could depend on the downstream cumulative count at the same time $N^\downarrow(t + \Delta t)$. Since we are solving the model in forward order of time, however, we do not know the value $N^\downarrow(t + \Delta t)$ when calculating $N^\uparrow(t + \Delta t)$. In an acyclic network, we could simply do the calculations such that $N^\downarrow(t + \Delta t)$ is calculated first before $N^\uparrow(t + \Delta t)$, using the concept of a topological order. In networks with cycles — virtually all realistic traffic networks — this will not work. Instead, what is best is to approximate the “infinite” backward wave speed with one which is as large as possible, basing the calculation on the most recent known point $N^\downarrow(t)$ (Figure 10.28). Effectively, this replaces the infinite backward wave speed with one of speed $L/\Delta t$. Equation (10.48) thus gives

$$N^\uparrow(t + \Delta t) = \min\{N^\uparrow(t) + q_{max}\Delta t, N^\downarrow(t) + k_j L\} \quad (10.73)$$

and

$$R(t) = \min\{q_{max}\Delta t, (N^\downarrow(t) + k_j L) - N^\uparrow(t)\}. \quad (10.74)$$

Equation (10.74) will ensure that the number of vehicles on the link will never exceed $k_j L$, assuming that this is true at time zero, as you are asked to show in Exercise 4.

10.5.4 Discussion

This chapter has presented four different link models: point queues, spatial queues, the cell transmission model, and the link transmission model. Although not initially presented this way, all four can be seen as special cases of the Lighthill-Whitham-Richards model. The point and spatial queue models can

be derived from particularly simple forms of the fundamental diagram (as well as from physical first principles, as in Section 10.1), while the cell transmission model and link transmission model are more general methods which can handle more sophisticated fundamental diagrams (typically triangular or trapezoidal in practice). The cell transmission model directly solves the LWR system of partial differential equations by discretizing in space and time, and applying a finite-difference approximation. The link transmission model is based on the Newell-Daganzo method. The primary distinction between these methods is that the Newell-Daganzo method only requires tracking the cumulative counts N at the upstream and downstream ends of each link in time, while the cell transmission model also requires tracking the number of vehicles n at intermediate cells within the link. However, the cell transmission model does not require storing any values from previous time steps, and can function entirely using the number of vehicles in each cell at the current time. The Newell-Daganzo method requires that some past cumulative counts be stored, for the amount of time needed for a wave to travel from one end of the link to the other. Which is more desirable depends on implementation details, and on the specific application context — at times it may be useful to know the distribution of vehicles within a link (as the cell transmission model gives), while for other applications this may be an irrelevant detail. One final advantage of the Newell-Daganzo method is that the values it gives are exact. In the cell transmission model, backward-moving shockwaves will tend to “spread out” as a numerical issue involved in the discretization; this will not happen when applying the Newell-Daganzo method. The exercises explore this issue in more detail. On the other hand, the cell transmission model is easier to explain to decision-makers, and its equations have intuitive explanations in terms of vehicles moving within a link and the amount of available space. The Newell-Daganzo method is a “deeper” method requiring knowledge of partial differential equations, and seems more difficult to convey to nontechnical audiences.

The point queue and spatial queue models have their places as well, despite their strict assumptions. The major flaw in the point queue model, from the standpoint of realism, is its inability to model queue spillbacks which occur when links are full. On the other hand, by ignoring this phenomenon, the point queue model is much more tractable, and is amenable even to closed-form expressions of delay and sensitivity to flows. It is also more robust to errors in input data, because queue spillback can introduce discontinuities in the network loading. There are cases where this simplicity and robustness may outweigh the (significant) loss in realism induced by ignoring spillbacks. The spatial queue model can represent spillbacks, but will tend to underestimate its effect due to its assumption of infinitely-fast backward moving shockwaves. Nevertheless, it can also lead to simpler analyses than the link transmission model.

10.6 Fancier Node Models

Section 10.2 introduced node modeling concepts, and how the sending and receiving flows from adjacent links are mapped to transition flows showing how many vehicles move from each incoming link to each outgoing link in a single time step. In analogy with Section 10.5, we now expand our discussion of node models beyond the simple intersection types presented thus far. A general intersection can have any number of incoming and outgoing links. General intersections can represent signal-controlled intersections, stop-controlled intersections, roundabouts, and so forth. Node models for general intersections are less standardized, and less well-understood, than node models for merges and diverges, and representations of these can vary widely in dynamic traffic assignment implementations. This section presents four alternative models for general intersections: the first for simple traffic signals, the second for intersections where drivers strictly take turns or have equal priority for all turning movements (such as an all-way stop), and the third where approaches have different priority levels and crossing conflicts must be considered, such as a two-way stop or signal with permitted phasing. A fourth node model is presented which allows one to use arbitrary models for turning movement capacity and intersection delay. As you read through this section, it would be helpful to think about how these models can be adapted or extended to represent more sophisticated signal types, roundabouts, and other types of junctions. Many of the ideas in these models are based on the ideas described for merges and diverges.

10.6.1 Basic signals

For the purposes of this section, a “basic signal” is one which (1) only has protected phases (no permitted turns which must yield to oncoming traffic) and (2) all turn movements corresponding to the same approach move simultaneously (for instance, where there are no turn lanes). See Figure 10.29 for an illustration. In this case, the node can be modeled as a diverge intersection, where the “upstream” link varies over time, depending on which approach has the green indication. Flows from other approaches (which have red indications) are set to zero. Following the notation for diverges, we use p_{hij} to reflect the proportion of the sending flow from approach (h, i) which wishes to leave via link (i, j) : these values must be nonnegative, and $\sum_{(i,j):(h,i,j) \in \Xi(i)} p_{hij} = 1$ for all approaches (h, i) . The algorithm is as follows:

1. Let $(h^*, i) \in \Gamma^{-1}(i)$ be the approach which has the green indication at the current time.
2. Calculate the fraction of flow which can move:

$$\phi = \min_{(i,j):(h^*,i,j) \in \Xi(i)} \left\{ \frac{R_{ij}}{p_{h^*ij} S_{h^*i}}, 1 \right\}. \quad (10.75)$$

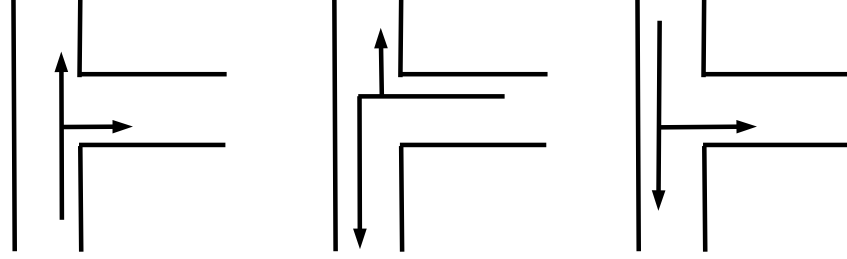


Figure 10.29: Phasing plan for a basic signal.

3. Calculate the transition flows for each turning movement:

$$y_{hij} = \begin{cases} \phi S_{hi} p_{hij} & \text{if } h = h^* \text{ and } (h, i, j) \in \Xi(i) \\ 0 & \text{otherwise} \end{cases} \quad (10.76)$$

In this implementation, one must be a little bit careful if the green times in the signal are not multiples of the time step Δt . It is possible to round the green times so that they are multiples of Δt , but this approach can introduce considerable error over the analysis period: for instance, assume that Δt is equal to six seconds, and a two-phase intersection has green times of 10 seconds and 14 seconds, respectively. Rounding to multiples of the time step would give both phases twelve seconds each, which seems reasonable enough; but over a three-hour analysis period, the phases would receive 75 and 105 minutes of green time in reality, as compared to 90 minutes each in simulation. In highly congested situations, this can introduce considerable error. This issue can be avoided if, instead of rounding, one gives the green indication to the approach which would have green in reality at that time. In the example above, the intersection has a cycle length of 24 seconds. So, when $t = 60$ seconds, we are 12 seconds into the third cycle; and at this point the green indication should be given to the second phase. In this way, there is no systematic bias introduced into the total green time each approach receives.

10.6.2 Equal priority movements

An all-way stop intersection is characterized by turn-taking: that is, vehicles have the opportunity to depart the intersection in the order in which they arrive. No turning movement has priority over any other, but the turning movements from different approaches interact with each other and may compete for space on the same outgoing link. This is different from the basic signal model, where the phasing scheme ensures that at most one approach is attempting to use an outgoing link at any given point of time.

Intersections with equal priority movements have characteristics of both diverges and merges. Like a diverge, if a vehicle is unable to turn into a down-

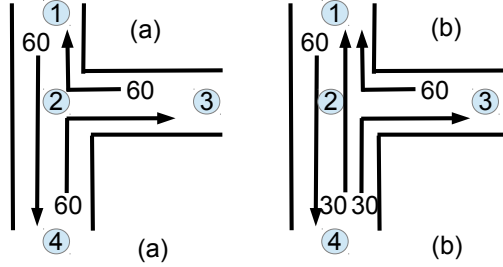


Figure 10.30: Oriented capacities for a three-leg intersection, where all approaches have $q_{max} = 60$. (a) All approaches map to a unique outgoing link. (b) The flow on approach (4,2) is split between two outgoing links.

stream link because of an obstruction, we assume that the vehicle obstructs all other vehicles from the same approach, respecting the FIFO principle. This means that the outflows for all of the turning movements corresponding to any approach must follow the same proportions as the number of drivers *wishing* to use all of these movements. Similar to a merge, we assume that if there are high sending flows from all the approaches, the fraction of the receiving flow allocated to each approach is divided up proportionally. However, instead of allocating the receiving flow R_{ij} to approach (h, i) based on the full capacity q_{max}^{hi} , we instead divide up the receiving flow based on the *oriented capacity*

$$q_{max}^{hij} = q_{max}^{hi} p_{hij}, \quad (10.77)$$

where p_{hij} is the proportion of the flow from approach (h, i) which wishes to exit on link (i, j) . (If turn movement $[h, i, j]$ is not in the allowable set $\Xi(i)$, then $p_{hij} = 0$.)

Multiplying the capacity by this proportion reflects the fact that an upstream approach can only make use of an available space on a downstream link if there is a vehicle wishing to turn. In Figure 10.30(a), each incoming link uses a unique exiting link, and thus can claim its full capacity. In Figure 10.30(b), half the vehicles on link (4,2) want to turn right and half wish to go straight, whereas all the vehicles on link (3,2) wish to turn right. Link (3,2) therefore has twice as many opportunities to fill available space on link (2,1), and thus its rightful share is twice that of link (4,2). For any two approaches $[h, i, j]$ and $[h', i, j]$ using the same outgoing link, we thus require that

$$\frac{y_{hij}}{y_{h'ij}} = \frac{q_{max}^{hij}}{q_{max}^{h'ij}}. \quad (10.78)$$

assuming that both approaches are fully competing for the link (i, j) . If an approach has a small sending flow, it may use less of its assigned receiving flow than equation (10.78) allocates, and this unused receiving flow may be used by other approaches.

We also define the *oriented sending flow*

$$S_{hij} = S_{hi}p_{hij} \quad (10.79)$$

to reflect the demand for the turning movement $[h, i, j]$. Following the same principles as the merge model, if the oriented sending flow from an approach is less than its proportionate share of a downstream link's receiving flow, its unused share will be divided among the other approaches with unserved demand still remaining, in proportion to their oriented capacities. If the oriented sending flow for a turn movement is greater than the oriented receiving flow for that movement, then by the FIFO principle applied to diverges, it will restrict flow to all other downstream links by the same proportion, and for any two turning movements (h, i, j) and (h, i, j') from the same approach we must have

$$\frac{y_{hij}}{y_{hij'}} = \frac{S_{hij}}{S_{hij'}} = \frac{p_{hij}}{p_{hij'}}. \quad (10.80)$$

We can rearrange this equation to show that the ratio y_{hij}/S_{hij} (the ratio of actual flow and desired flow for any turning movement) is uniform for all the turning movements approaching from link (h, i) — this ratio plays the same role as ϕ in a diverge.

The presence of multiple incoming and outgoing links causes another complication, in that the flows between approaches are all linked together. If an approach is restricted by the receiving flow of a downstream link, flow from that approach is restricted to *all* other downstream links. This means that the approach may not fully consume its “rightful share” of another downstream link, thereby freeing up additional capacity for a different approach. Therefore, we cannot treat the approaches or downstream links separately or even sequentially in a fixed order, because we do not know *a priori* how these will be linked together.

However, there is an algorithm which generates a consistent solution despite these mutual dependencies. In this algorithm, each approach link can be *demand-constrained*, or *supply-constrained by a downstream link*. If an approach is demand-constrained, its oriented sending flow to all downstream links is less than its rightful share, and therefore all of the sending flow can move. If an approach is supply-constrained by link (i, j) , then the approach is unable to move all of its sending flow, and the fraction which can move is dictated by link (i, j) . (That is, receiving flow on (i, j) is the most restrictive constraint for the approach). The algorithm must determine which links are demand-constrained, and which are supply-constrained by a downstream link.

To find such a solution, we define two sets of auxiliary variables, \tilde{S}_{hij} to reflect the amount of unallocated sending flow for movement $[h, i, j]$, and \tilde{R}_{ij} to reflect the amount of unallocated receiving flow for outgoing link (i, j) . These are initialized to the oriented sending flows and link receiving flows, and reduced iteratively as flows are assigned and the available sending and receiving flows are used up. The algorithm also uses the notion of *active* turning movements; these are turning movements whose flows can still be increased. A turning movement

$[h, i, j]$ becomes inactive either when \tilde{S}_{hij} drops to zero (all vehicles that wish to turn have been assigned), or when $\tilde{R}_{ij'}$ drops to zero for *any* outgoing link (i, j') that approach (h, i) is using (that is, for which $p_{hij'} > 0$). Allocating all of the receiving flow for one outgoing link can thus impact flow on turning movements which use other outgoing links, because of the principle that vehicles wishing to turn will block others, as expressed in equation (10.80). The set of active turning movements will be denoted by A ; a turning movement remains active until we have determined whether it is supply-constrained or demand-constrained.

At each stage of the algorithm, we will increase the flows for all active turning movements. We must increase these flows in a way which is consistent both with the turning fractions (10.80), and with the division of receiving flow for outgoing links given by equation (10.78), and we will use α_{hij} to reflect the rate of increase for turning movement (h, i, j) . The absolute values of these α values do not matter, only their proportions, so you can scale them in whatever way is most convenient to you. Often it is easiest to pick one turning movement (h, i, j) and fix its α value either to one, or to its oriented sending flow. The turning proportions from (h, i) then fix the α values for all other turning movements from the same approach. You can then use equation (10.78) to determine α values for turning movements competing for the same outgoing link, then use the turning fractions for the upstream link on those turning movements, and so on until a consistent set of α values has been determined.

The algorithm then increases the active turning movement flows in these proportions until some movement becomes inactive, because its sending flow or the receiving flow on its outgoing link becomes exhausted. The process is then repeated with the smaller set of turning movements which remain active, and continues until all possible flows have been assigned. This algorithm is a bit more involved than the node models seen thus far, and you may find it helpful to follow the example below as you read through the algorithm steps.

1. Initialize by calculating oriented capacities and sending flows using equations (10.77) and (10.79); by setting $y_{hij} \leftarrow 0$ for all $[h, i, j] \in \Xi(i)$; by setting $\tilde{S}_{hij} \leftarrow S_{hij}$ for all $[h, i, j] \in \Xi(i)$ and $\tilde{R}_{ij} \leftarrow R_{ij}$ for all $(i, j) \in \Gamma(i)$; and by declaring active all turning movements with positive sending flow: $A \leftarrow \{[h, i, j] \in \Xi(i) : S_{hij} > 0\}$.
2. Identify a set of α_{hij} values for all active turning movements which is consistent with the turning fractions ($\alpha_{hij}/\alpha_{hij'} = p_{hij}/p_{hij'}$ for all turning movements from the same incoming link) and oriented capacities ($\alpha_{hij}/\alpha_{h'ij} = q_{max}^{hij}/q_{max}^{h'ij}$ for all turning movements to the same outgoing link).
3. For each outgoing link $(i, j) \in \Gamma(i)$, identify the rate α_{ij} at which its receiving flow will be reduced, by adding α_{hij} for all active turn movements whose outgoing link is (i, j) : $\alpha_{ij} \leftarrow \sum_{[h, i, j] \in A} \alpha_{hij}$.
4. Determine the point at which some turning movement will become inac-

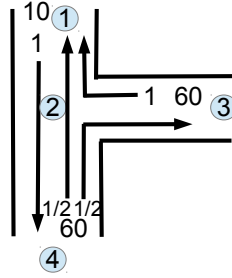


Figure 10.31: Example of equal priorities algorithm, showing sending flows and turning proportions. All links have capacity of 60.

tive, by calculating

$$\theta = \min \left\{ \min_{[h,i,j] \in A} \left\{ \frac{\tilde{S}_{hij}}{\alpha_{hij}} \right\}, \min_{(i,j) \in \Gamma(i): \alpha_{ij} > 0} \left\{ \frac{\tilde{R}_{ij}}{\alpha_{ij}} \right\} \right\}. \quad (10.81)$$

5. Increase flows for active turning movements, and update unallocated sending and receiving flows: for all $[h, i, j] \in A$ update $y_{hij} \leftarrow y_{hij} + \theta \alpha_{hij}$, $\tilde{S}_{hij} \leftarrow \tilde{S}_{hij} - \theta \alpha_{hij}$, and $\tilde{R}_{ij} \leftarrow \tilde{R}_{ij} - \theta \alpha_{hij}$.
6. Update the set of active turning movements, by removing from A any turning movement $[h, i, j]$ for which $\tilde{S}_{hij} = 0$ or for which $\tilde{R}_{ij} = 0$ for any $(i, j') \in \Gamma(i)$ which is being used ($p_{hij'} > 0$).
7. If there are any turning movements which are still active ($A \neq \emptyset$), return to step 3. Otherwise, stop.

As a demonstration, consider the intersection in Figure 10.31, where all links (incoming and outgoing) have the same capacity of 60 vehicles per time step, and the sending flows S_{hij} and turning proportions p_{hij} are shown. None of the downstream links is congested, so their receiving flows are equal to the capacity. (In the figure, two sets of numbers are shown for each approach; the “upstream” number is the sending flow and the “downstream” number(s) are the proportions.) The oriented capacities can be seen in Figure 10.30(b). For this example, the six turning movements will be indexed in the following order:

$$\Xi(2) = \{[1, 2, 3], [1, 2, 4], [3, 2, 1], [3, 2, 4], [4, 2, 1], [4, 2, 3]\} \quad (10.82)$$

All vectors referring to turning movements will use this ordering for their components.

Step 1 of the algorithm initializes the oriented sending flows using equation (10.79),

$$\mathbf{S} = [S_{hij}] = \begin{bmatrix} 0 & 10 & 60 & 0 & 30 & 30 \end{bmatrix} \quad (10.83)$$

and the oriented capacities using equation (10.77)

$$\mathbf{q}_{\max} = [q_{max}^{hij}] = \begin{bmatrix} 0 & 60 & 60 & 0 & 30 & 30 \end{bmatrix}. \quad (10.84)$$

The step also initializes the turning movement flows and auxiliary variables:

$$\mathbf{y} \leftarrow [0 \ 0 \ 0 \ 0 \ 0 \ 0], \quad (10.85)$$

$$\tilde{\mathbf{S}} \leftarrow [S_{hij}] = [0 \ 10 \ 60 \ 0 \ 30 \ 30], \quad (10.86)$$

and

$$\tilde{\mathbf{R}} = [\tilde{R}_{21} \ \tilde{R}_{23} \ \tilde{R}_{24}] = [60 \ 60 \ 60]. \quad (10.87)$$

The set of active turn movements is $A = \{[1, 2, 4], [3, 2, 1], [4, 2, 1], [4, 2, 3]\}$.

Step 2 of the algorithm involves calculation of a set of consistent α_{hij} values. One way of doing this is to start by setting $\alpha_{423} \leftarrow S_{423} = 30$. The turning fractions from (4, 2) then require that $\alpha_{421} = \alpha_{423} = 30$. The allocation rule for outgoing link (2, 1) then forces $S_{321} = 60$: the oriented capacity for $[3, 2, 1]$ is twice that of $[4, 2, 1]$, and the α values must follow the same proportion. Turning movement $[1, 2, 4]$ is independent of all of the other turning movements considered thus far, so we can choose its α value arbitrarily; say, $\alpha_{124} \leftarrow S_{124} = 10$. (You should experiment around with different ways of calculating these α values, and convince yourself that the final flows are the same as long as the proportions of α values for interdependent turning movements are the same.) We thus have

$$\boldsymbol{\alpha} = [\alpha_{hij}] = [0 \ 10 \ 60 \ 0 \ 30 \ 30]. \quad (10.88)$$

The α values for inactive turning movements have been set to zero for clarity; their actual value is irrelevant because they will not be used in any of the steps that follow.

With these flow increments, flow on outgoing links (2, 1), (2, 3), and (2, 4) will be $\alpha_{21} = 90$, $\alpha_{23} = 30$, and $\alpha_{24} = 10$, as dictated by Step 3. In Step 4, we determine how much we can increase the flow at the rates given by $\boldsymbol{\alpha}$ until some movement becomes inactive. We have

$$\theta = \min \left\{ \frac{\tilde{S}_{124}}{10}, \frac{\tilde{S}_{321}}{60}, \frac{\tilde{S}_{421}}{30}, \frac{\tilde{S}_{423}}{30}, \frac{\tilde{R}_{21}}{90}, \frac{\tilde{R}_{23}}{30}, \frac{\tilde{R}_{24}}{10} \right\} \quad (10.89)$$

or $\theta = 2/3$.

We can now adjust the flows, as in Step 5. We increase the flow on each active turning movement by $\frac{2}{3}\alpha_{hij}$, giving

$$\mathbf{y} \leftarrow [0 \ 6\frac{2}{3} \ 40 \ 0 \ 20 \ 20]. \quad (10.90)$$

We subtract these flow increments from the auxiliary sending and receiving flows, giving

$$\tilde{\mathbf{S}} \leftarrow [0 \ 3\frac{1}{3} \ 20 \ 0 \ 10 \ 10] \quad (10.91)$$

and

$$\tilde{\mathbf{R}} \leftarrow [0 \ 40 \ 53\frac{1}{3}]. \quad (10.92)$$

Step 6 updates the set of active turning movements. With the new $\tilde{\mathbf{S}}$ and $\tilde{\mathbf{R}}$ values, we see that $[3, 2, 1]$ and $[4, 2, 1]$ have become inactive, since there is no

remaining receiving flow on link $(2, 1)$. Furthermore, this inactivates movement $[4, 2, 3]$: even though there are still travelers that wish to turn in this direction, and space on the downstream link (\tilde{S}_{423} and \tilde{R}_{23} are still positive), they are blocked by travelers waiting to use movement $[4, 2, 1]$. So, there is only one active movement remaining, $A \leftarrow \{[1, 2, 4]\}$, and we must return to step 3.

In Step 3, we must recalculate the α_{ij} values because some of the turning movements are inactive. With the new set A , we have $\alpha_{21} = \alpha_{23} = 0$ and $\alpha_{24} = 10$. The new step size is

$$\theta = \min \left\{ \frac{\tilde{S}_{124}}{3\frac{1}{3}}, \frac{\tilde{R}_{24}}{53\frac{1}{3}} \right\} = \frac{1}{3}. \quad (10.93)$$

We then increase the flows, increasing y_{124} by $10 \times \frac{1}{3}$ to 10, decreasing \tilde{S}_{124} to zero, and decreasing \tilde{R}_{24} to 50. This change inactivates movement $[1, 2, 4]$. Since there are no more active turning movements, the algorithm terminates, and the final vector of flows is

$$\mathbf{y} = [0 \quad 10 \quad 40 \quad 0 \quad 20 \quad 20]. \quad (10.94)$$

10.6.3 Intersections with priority

Intersections which allow crossing conflicts are more complex to model than the intersection types described above. These include intersections with stop control only on some of the approaches, but not all, or signalized intersections with permitted movements that must yield to another traffic stream. One way to model this type of intersection is to introduce a set C of conflict points. Like outgoing links, we associate a receiving flow R_c indicating the maximum number of vehicles which can pass through this conflict point during a single time step Δt . For each conflict point $c \in C$, let $\Xi(c)$ denote the turning movements which make use of this conflict point. Then, for each turning movement $[h, i, j] \in \Xi(c)$, we define a (strictly positive) priority parameter β_{hij}^c . These priority parameters are interpreted through their ratios: the share of the receiving flow allocated to turning movement $[h, i, j]$ relative to that allocated to turning movement $[h', i, j']$ is in the same proportion as the ratio

$$\beta_{hij}^c p_{hij} / \beta_{h'ij'}^c p_{h'ij'}. \quad (10.95)$$

These play an analogous role to the oriented capacities defined in the previous subsection, in suggesting how the conflict point receiving flow should be divided among the competing approaches. They are used more generally, however, to reflect priority rules. For instance, a through movement often has priority over a turning movement which crosses it. Even if the capacity and sending flows of the through lane and turn lane are the same, the through lane should have access to a greater share of the crossing point's receiving flow. If, at full saturation, ten through vehicles move for every turning vehicle, then the β value for the through movement should be ten times the β value for the turning movement.

The node model is simplified if we assume that the β values are all strictly positive. You might find this unrealistic: in the example above, if the turning movement must yield to the through movement, then at full saturation perhaps *no* turning vehicles could move. In practice, however, priority rules are not strictly obeyed as traffic flows near saturation. Polite through drivers may stop to let turning drivers move, or aggressive turning drivers might force their way into the through stream. (Think of what would happen at a congested freeway if vehicles merging from an onramp took the “yield” sign literally!) The requirement of strictly positive β values thus has some practical merits, as well as mathematical ones. The exercises explore ways to generalize this node model, including strict priority, and cases where different turning movements may consume different amounts of the receiving flow (for instance, if they are moving at different speeds).

We can now adapt the algorithm for equal priority for the case of intersections with different priorities. The algorithm is augmented by adding receiving flows R_c and auxiliary receiving flows \tilde{R}_c for each conflict point, and we extend some of the computations to include the set of conflict points. First, we require that the ratio of α values for two turning movements using the same crossing point follow the ratio of the β values, assuming saturated conditions:

$$\frac{\alpha_{hij}}{\alpha_{h'ij'}} = \frac{\beta_{hij}^c p_{hij}}{\beta_{h'ij'}^c p_{h'ij'}} \quad \forall c \in C; [h, i, j], [h', i, j'] \in \Xi(c). \quad (10.96)$$

Note that the β values are multiplied by the relevant turning fraction for the movements. As with the oriented capacity, this respects the fact that the more flow is attempting to turn in a particular direction, the more opportunities or gaps will be available for it to claim. Second, we must calculate the inflow rates to conflict points given the α values from active turning movements:

$$\alpha_c = \sum_{[h, i, j] \in \Xi(c) \cap A} \alpha_{hij}. \quad (10.97)$$

Third, the calculation of the step size θ must now include obstructions from conflict points:

$$\theta = \min \left\{ \min_{[h, i, j] \in A} \left\{ \frac{\tilde{S}_{hij}}{\alpha_{hij}} \right\}, \min_{(i, j) \in \Gamma(i): \alpha_{ij} > 0} \left\{ \frac{\tilde{R}_{ij}}{\alpha_{ij}} \right\}, \min_{c \in C: \alpha_c > 0} \left\{ \frac{\tilde{R}_c}{\alpha_c} \right\} \right\}. \quad (10.98)$$

With these modifications, the algorithm proceeds in the same way as before. Specifically,

1. Initialize by calculating oriented capacities and sending flows using equations (10.77) and (10.79); by setting $y_{hij} \leftarrow 0$ for all $[h, i, j] \in \Xi(i)$; by setting $\tilde{S}_{hij} \leftarrow S_{hij}$ for all $[h, i, j] \in \Xi(i)$, $\tilde{R}_{ij} \leftarrow R_{ij}$ for all $(i, j) \in \Gamma(i)$, and $\tilde{R}_c \leftarrow R_c$ for each $c \in C$; and by declaring active all turning movements with positive sending flow: $A \leftarrow \{[h, i, j] \in Xi(i) : S_{hij} > 0\}$.

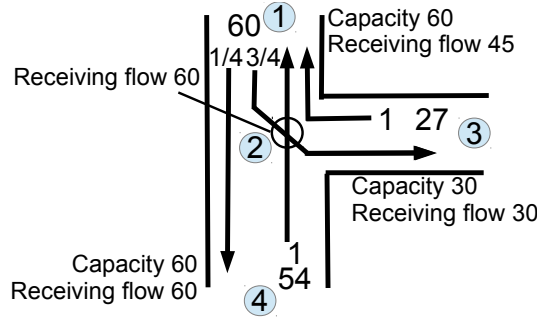


Figure 10.32: Example of intersection priority algorithm, showing sending flows and turning proportions. Link capacity and receiving flows shown.

2. Identify a set of α_{hij} values for all active turning movements which is consistent with the turning fractions ($\alpha_{hij}/\alpha_{hij'} = p_{hij}/p_{hij'}$ for all turning movements from the same incoming link), oriented capacities ($\alpha_{hij}/\alpha_{h'ij} = q_{max}^{hij}/q_{max}^{h'ij}$ for all turning movements to the same outgoing link), and conflict points based on equation (10.96).
3. For each outgoing link $(i, j) \in \Gamma(i)$ and conflict point $c \in C$, identify the rate α_{ij} at which its receiving flow will be reduced, by adding α_{hij} for all active turn movements whose outgoing link is (i, j) : $\alpha_{ij} \leftarrow \sum_{[h,i,j] \in A} \alpha_{hij}$ for links, and equation (10.97) for conflict points.
4. Determine the point at which some turning movement will become inactive, by calculating θ using equation (10.98).
5. Increase flows for active turning movements, and update unallocated sending and receiving flows: for all $[h, i, j] \in A$ update $y_{hij} \leftarrow y_{hij} + \theta \alpha_{hij}$, $\tilde{S}_{hij} \leftarrow \tilde{S}_{hij} - \theta \alpha_{hij}$; for all $(i, j) \in \Gamma(i)$ update $\tilde{R}_{ij} \leftarrow \tilde{R}_{ij} - \theta \alpha_{ij}$, and for all $c \in C$ update $\tilde{R}_c \leftarrow \tilde{R}_c - \theta \alpha_c$.
6. Update the set of active turning movements, by removing from A any turning movement $[h, i, j]$ for which $\tilde{S}_{hij} = 0$, for which $\tilde{R}_{ij'} = 0$ for any $(i, j') \in \Gamma(i)$ which is being used ($p_{hij'} > 0$), or for which $\tilde{R}_c = 0$ for any movement $[h, i, j'] \in \Xi(c)$ and conflict point c .
7. If there are any turning movements which are still active ($A \neq \emptyset$), return to step 3. Otherwise, stop.

As a demonstration, consider the intersection in Figure 10.32, where all links (incoming and outgoing) have the same capacity of 60 vehicles per time step, and the sending flows S_{hi} and turning proportions p_{hij} are shown. (In the figure, two sets of numbers are shown for each approach; the “upstream” number is the sending flow and the “downstream” number(s) are the proportions.) Links (2, 4)

and (2, 3) have receiving flows of 60 vehicles, while link (2, 1) has a receiving flow of only 45 vehicles. There is one conflict point, indexed c , which is marked with a circle in Figure 10.32). Conflict point c has a receiving flow of 60, and the turning movement $[1, 2, 3]$ must yield to the through movement $[4, 2, 1]$, as reflected by the ratio $\beta_{421}/\beta_{123} = 5$. For this example, the six turning movements will be indexed in the following order:

$$\Xi(2) = \{[1, 2, 3], [1, 2, 4], [3, 2, 1], [3, 2, 4], [4, 2, 1], [4, 2, 3]\} \quad (10.99)$$

All vectors referring to turning movements will use this ordering for their components.

Step 1 of the algorithm initializes the oriented sending flows using equation (10.79),

$$\mathbf{S} = [S_{hij}] = \begin{bmatrix} 45 & 15 & 27 & 0 & 54 & 0 \end{bmatrix} \quad (10.100)$$

and the oriented capacities using equation (10.77)

$$\mathbf{q}_{\max} = [q_{max}^{hij}] = \begin{bmatrix} 30 & 30 & 30 & 0 & 60 & 0 \end{bmatrix}. \quad (10.101)$$

The step also initializes the turning movement flows and auxiliary variables:

$$\mathbf{y} \leftarrow \begin{bmatrix} 0 & 0 & 0 & 0 & 0 & 0 \end{bmatrix}, \quad (10.102)$$

$$\tilde{\mathbf{S}} \leftarrow [\tilde{S}_{hij}] = \begin{bmatrix} 45 & 15 & 27 & 0 & 54 & 0 \end{bmatrix}, \quad (10.103)$$

and

$$\tilde{\mathbf{R}} = [\tilde{R}_{21} \quad \tilde{R}_{23} \quad \tilde{R}_{24} \quad \tilde{R}_c] = \begin{bmatrix} 45 & 30 & 60 & 60 \end{bmatrix}. \quad (10.104)$$

The set of active turn movements is $A = \{[1, 2, 3], [1, 2, 4], [3, 2, 1], [4, 2, 1]\}$.

Step 2 of the algorithm involves calculation of a set of consistent α_{hij} values. One way of doing this is to start by setting $\alpha_{123} \leftarrow S_{123} = 15$. The turning fractions from (1, 2) then require that $\alpha_{124} = \alpha_{123} = 45$. The allocation rule (10.96) for conflict point c forces $S_{421} = 300$: the β value for $[4, 2, 1]$ is five times that of the β value for $[1, 2, 3]$, and the turning fraction for $[4, 2, 1]$ is a third higher. Since $S_{421} = 300$, we must have $S_{321} = 150$ because the oriented capacity of $[3, 2, 1]$ is half that of $[4, 2, 1]$. This gives the flow increments

$$\boldsymbol{\alpha} = [\alpha_{hij}] = \begin{bmatrix} 45 & 15 & 150 & 0 & 300 & 0 \end{bmatrix}. \quad (10.105)$$

With these flow increments, flow on outgoing links (2, 1), (2, 3), (2, 4), and the conflict point c will be $\alpha_{21} = 450$, $\alpha_{23} = 45$, $\alpha_{24} = 15$, and $\alpha_c = 345$, as dictated by Step 3. In Step 4, we determine how much we can increase the flow at the rates given by $\boldsymbol{\alpha}$ until some movement becomes inactive. We have

$$\theta = \min \left\{ \frac{\tilde{S}_{123}}{15}, \frac{\tilde{S}_{124}}{45}, \frac{\tilde{S}_{321}}{150}, \frac{\tilde{S}_{421}}{300}, \frac{\tilde{R}_{21}}{450}, \frac{\tilde{R}_{23}}{45}, \frac{\tilde{R}_{24}}{15}, \frac{\tilde{R}_c}{345} \right\} \quad (10.106)$$

or $\theta = \frac{1}{10}$.

We can now adjust the flows, as in Step 5. We increase the flow on each active turning movement by $\frac{1}{10}\alpha_{hij}$, giving

$$\mathbf{y} \leftarrow [4.5 \quad 1.5 \quad 15 \quad 0 \quad 30 \quad 0] . \quad (10.107)$$

We subtract these flow increments from the auxiliary sending and receiving flows, giving

$$\tilde{\mathbf{S}} \leftarrow [40.5 \quad 13.5 \quad 12 \quad 0 \quad 24 \quad 0] \quad (10.108)$$

and

$$\tilde{\mathbf{R}} \leftarrow [0 \quad 25.5 \quad 58.5 \quad 25.5] . \quad (10.109)$$

Step 6 updates the set of active turning movements. With the new $\tilde{\mathbf{S}}$ and $\tilde{\mathbf{R}}$ values, we see that $[3, 2, 1]$ and $[4, 2, 1]$ have become inactive, since there is no remaining receiving flow on link $(2, 1)$. So, there are two active movements remaining, $A \leftarrow \{[1, 2, 3], [1, 2, 4]\}$, and we must return to step 3.

In Step 3, we must recalculate the α_{ij} values because some of the turning movements are inactive. With the new set A , we have $\alpha_{21} = 0$, $\alpha_{23} = 45$, $\alpha_{24} = 15$, and $\alpha_c = 45$. The new step size is

$$\theta = \min \left\{ \frac{\tilde{S}_{123}}{45}, \frac{\tilde{S}_{124}}{15}, \frac{\tilde{R}_{23}}{45}, \frac{\tilde{R}_{24}}{15}, \frac{\tilde{R}_c}{45} \right\} = \frac{17}{30} . \quad (10.110)$$

We then increase the flows, increasing y_{123} by $45 \times \frac{17}{30}$ to 30 and y_{124} by $15 \times \frac{17}{30}$ to 10, decreasing \tilde{S}_{123} to 15, \tilde{S}_{124} to 5, \tilde{R}_{23} to 0, \tilde{R}_{24} to 50, and \tilde{R}_c to 0. This change inactivates movement $[1, 2, 3]$; and movement $[1, 2, 4]$ is then inactivated because these movement's flows are blocked by vehicles waiting to take $[1, 2, 3]$. Since there are no more active turning movements, the algorithm terminates, and the final vector of flows is

$$\mathbf{y} = [30 \quad 10 \quad 15 \quad 0 \quad 30 \quad 0] . \quad (10.111)$$

10.7 Exercises

1. [14] Table 10.7 shows cumulative inflows and outflows to a link with a capacity of 10 vehicles per time step, and a free-flow time of 2 time steps. Use the point queue model to calculate the sending and receiving flow for each time step.
2. [14] Repeat Exercise 1 with the spatial queue model. Assume that the jam density is such that at most 20 vehicles can fit on the link simultaneously.
3. [33] In the point queue model, if the inflow and outflow rates q_{in} and q_{out} are constants with $q_{in} \geq q_{out}$, show that the travel time experienced by the n -th vehicle is $L/u_f + \frac{n}{q_{in}} \left(\frac{q_{in}}{q_{out}} - 1 \right)$. The same result holds for the spatial queue model, if there is no spillback.

Table 10.7: Upstream and downstream counts for Exercises 1 and 2.

t	$N^\uparrow(t)$	$N^\downarrow(t)$
0	0	0
1	5	0
2	10	0
3	15	2
4	20	4
5	20	6
6	20	10
7	20	15
8	20	20
9	20	20

4. [24] In the spatial queue model, show that if the total number of vehicles on a link is at most $k_j L$ at time t , then the total number of vehicles on the link at time $t + 1$ is also at most $k_j L$.
5. [13] In a merge, approach 1 has a sending flow of 50 and a capacity of 100. Approach 2 has a sending flow of 100 and a capacity of 100. Report the number of vehicles moving from each approach if (a) the outgoing link has a receiving flow of 150; (b) the outgoing link has a receiving flow of 120; and (c) the outgoing link has a receiving flow of 100.
6. [32] Extend the merge model of Section 10.2.2 to a merge node with three incoming links.
7. [43] Show that the formula (10.14) indeed captures both case II and case III of the merge model presented in Section 10.2.2.
8. [32] Show that the merge model of Section 10.2.2 satisfies all the desiderata of Section 10.2.
9. [41] Consider an alternative merge model for the congested case, which allocates the receiving flow proportional to the *sending flows* of the incoming links, rather than proportional to the *capacities* of the incoming links as was done in Section 10.2.2. Show that this model does *not* satisfy the invariance principle.
10. [13] In a diverge, the incoming link has a sending flow of 120, 25% of the vehicles want to turn onto outgoing link 1, and the remainder want to turn onto outgoing link 2. Report the number of vehicles moving to outgoing links 1 and 2 if their respective receiving flows are (a) 80 and 100; (b) 80 and 60; (c) 10 and 40.
11. [10] Extend the diverge model of Section 10.2.3 to a diverge node with three outgoing links.

12. [21] Develop a model for a diverge node with two outgoing links, in which flows waiting to enter one link do *not* block flows entering the other link. When might this model be more appropriate?
13. [25] Show that the diverge model of Section 10.2.3 satisfies all the desiderata of Section 10.2.
14. [10] In the network loading procedure in Section 10.3, we specified that centroid connectors starting at origins should have high jam density, and those ending at destinations should have high capacity. Would anything go wrong if centroid connectors starting at origins also had high capacity? What if centroid connectors ending at destinations had high jam density?
15. [12] Draw trajectory diagrams which reflect the following situations: (a) steady-state traffic flow, no vehicles speeding up or slowing down; (b) vehicles approaching a stop sign, then continuing; (c) a slow semi truck merges onto the roadway at some point in time, then exits at a later point in time. Draw at least five vehicle trajectories for each scenario.
16. [33] The relationship $q = uk$ can be used to transform the fundamental diagram (which relates density and flow) into a relationship between density and speed, and vice versa. For each of these three speed-density relationships, derive the corresponding fundamental diagram.
 - (a) The Greenshields (linear) model: $u = u_f(1 - k/k_j)$.
 - (b) The Greenberg (logarithmic) model: $u = C \log(k/k_j)$ where C is a constant.
 - (c) The Underwood (exponential) model: $u = u_f \exp(-k/k_c)$.
 - (d) The Pipes model: $u = u_f(1 - (k/k_j)^n)$ where n is a constant. (The Greenshields model is a special case when $n = 1$.)
 - (e) What features of the Greenberg and Underwood models make them less suitable for dynamic network loading? (Hint: Draw plots of these speed-density diagrams.)
17. [51] Which of the following statements are true with the LWR model and a concave fundamental diagram?
 - (a) Speed uniquely defines the values of flow and density.
 - (b) It is possible for higher density to be associated with higher speed.
 - (c) No shockwave can move downstream faster than the free-flow speed.
 - (d) Flow uniquely defines the values of speed and density.
 - (e) With a triangular fundamental diagram, traffic speed is constant for subcritical densities.
 - (f) Density uniquely defines the values of speed and flow.

18. [44] The relationship $q = uk$ can transform the fundamental diagram (which relates density and flow) into a relationship between speed and flow. The speed-flow relationship is traditionally plotted with the speed on the vertical axis and flow on the horizontal axis.
- (a) Show that for any concave fundamental diagram and any flow value q less than the capacity, there are exactly two possible speeds u_1 and u_2 producing the flow q , one corresponding to subcritical (uncongested) conditions and the other corresponding to supercritical (congested) conditions.
 - (b) Derive and plot the speed-flow relationship corresponding to the Greenshields model of Exercise 16. Express this relationship with two functions $u_1(q)$ and $u_2(q)$ corresponding to uncongested and congested conditions, respectively; these functions should have a domain of $[0, q_{max}]$ and intersect at capacity.
 - (c) Derive and plot the fundamental diagram corresponding to the Highway Capacity Manual speed-flow relation for basic freeway segments. The capacity of such a segment is given by $q_{max} = 1800 + 5u_f$. In this equation, speeds are measured in km/hr and flows in veh/hr/lane:

$$u_1(q) = \begin{cases} u_f & \text{if } q \leq 3100 - 15u_0 \\ u_f - \frac{23u_0 - 1800}{28} \left(\frac{q + 15u_f - 3100}{20u_f - 1300} \right)^{2.6} & \text{if } 3100 - 15u_0 \leq q \leq q_{max} \end{cases} \quad (10.112)$$

$$u_2(q) = 28q \quad (10.113)$$

19. [22] For each of these fundamental diagrams, derive the speed-density function (that is, the travel speed for any given density value), and provide a sketch.
- (a) $Q(k) = Ck(k_j - k)$, where C is a constant and k_j is the jam density.
 - (b) $Q(k) = \min\{u_f k, w(k_j - k)\}$
 - (c) $Q(k) = \min\{u_f k, q_{max}, w(k_j - k)\}$
20. [44] Consider a long, uninterrupted freeway with a capacity of 4400 vehicles per hour, a jam density of 200 vehicles per mile, and a free-flow speed of 75 miles per hour. Initially, freeway conditions are uniform and steady with a subcritical flow of 2000 vehicles per hour. An accident reduces the roadway capacity to 1000 veh/hr for thirty minutes. Draw a shockwave diagram to show the effects of this accident, reporting the space-mean speed, volume, and density in each region of your diagram, and the speed and direction of each shockwave. Assume that the fundamental diagram takes the shape of the Greenshields model (Exercise 16), and that a stopped queue discharges at capacity.

21. [46] Consider a roadway with a linear-speed density relationship (cf. Exercise 16) whose capacity is 2000 veh/hr and free-flow speed is 40 mi/hr. Initially, the flow is 1000 veh/hr and uncongested. A traffic signal is red for 45 seconds, causing several shockwaves. When the light turns green, the queue discharges at capacity.
- Sketch a time-space diagram, indicating all of the shockwaves which are formed.
 - Calculate the speed and direction of each shockwave from your diagram.
 - What is the minimum green time needed to ensure that no vehicle has to stop more than once before passing the intersection? (Neglect any yellow time, reaction time, etc. Assume that when the signal is green, people move immediately, and that when it is red, people stop immediately.)
22. [56] Consider a single-lane roadway with a triangular fundamental diagram, a free-flow speed of 60 mi/hr, a backward wave speed of 30 mi/hr, and a jam density of 200 veh/mi. Initially, traffic flow is uncongested, and the volume is half of capacity. A slow-moving truck enters the roadway at time $t = 0$, and travels at 20 mi/hr. This vehicle turns off of the roadway one mile later.
- What is the capacity of the roadway?
 - At time $t = 2$ minutes, you are a quarter of a mile behind the truck. Use the Newell-Daganzo method to determine how many vehicles are between you and the truck.
 - In total, how many shockwaves are generated by the slow-moving truck? Sketch them on a trajectory diagram.
23. [11] Show that a shockwave connecting two uncongested (subcritical) traffic states always moves downstream, while a shockwave connecting two congested (supercritical) traffic states always moves upstream. This is related to the observation in the chapter that “uncongested states propagate downstream, and congested states propagate upstream.”
24. [36] This exercise asks you to fill in some details of the example in Section 10.4 where the fundamental diagram was $Q(k) = \frac{1}{240}k(240 - k)$ and the cumulative count map was $N(x, t) = 60t - 120x + \frac{60x^2}{t+1}$. Times are measured in minutes, and distances in miles.
- Calculate the capacity, jam density, and free-flow speed associated with this fundamental diagram.
 - Verify that the conservation relationship (10.25) is satisfied by the flow and density maps $q(x, t)$ and $k(x, t)$.

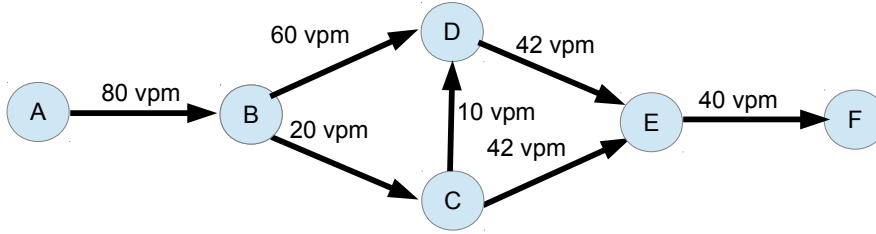


Figure 10.33: Network for Exercise 25

Table 10.8: Current cell occupancies for Exercise 27.

Cell 1	Cell 2	Cell 3	Cell 4
8	10	30	5

- (c) Verify that the density and flow maps are consistent with the given fundamental diagram.
- (d) Calculate the speed $u(x, t)$ at each point and time. Are vehicles accelerating, decelerating, or maintaining a constant speed?
25. [45] Consider the network in Figure 10.33, where each link has a free-flow time of 5 minutes and a capacity shown on the figure, and vehicles split equally at each diverge (that is, $p_{ABC} = p_{ABD} = p_{BCD} = p_{BCE} = 1/2$ at all times). Vehicles enter the network at a rate of 80 veh/min for 20 minutes, and then the inflow rate drops to zero. Perform dynamic network loading, using point queues for the link models. For each link in the network, plot the cumulative counts N^\uparrow and N^\downarrow over time, as well as the sending flow and receiving flow over time. At what time does the last vehicle leave the network?
26. [25] Write the formula for the fundamental diagram $Q(k)$ in the cell transmission model example depicted in Table 10.4.
27. [13] A link is divided into four cells; on this link the capacity is 10 vehicles per time step, each cell can hold at most 40 vehicles, and the ratio of backward wave speed to free-flow speed is 0.5. Currently, the number of vehicles in each cell is as in Table 10.8 (Cell 1 is at the upstream end of the link, Cell 4 at the downstream end.) Calculate the number of vehicles that will move between each pair of cells in the current time interval (that is, the y_{12} , y_{23} , and y_{34} values.), and the number of vehicles in each cell at the start of the next time interval. Assume no vehicles enter or exit the link.
28. [23] Table 10.9 shows cumulative inflows and outflows to a link with a capacity of 10 vehicles per time step, a free-flow time of 2 time steps, and

Table 10.9: Upstream and downstream counts for Exercise 28.

t	$N^\uparrow(t)$	$N^\downarrow(t)$
0	0	0
1	5	0
2	10	0
3	15	2
4	16	4
5	17	6
6	20	10
7	20	15

a backward wave time of 4 time steps. At jam density, there are 20 vehicles on the link. Use the link transmission model to calculate the sending flow $S(7)$ and the receiving flow $R(7)$.

29. [44] (*Exploring shock spreading.*) A link is seven cells long; at most 15 vehicles can fit into each cell, the capacity is 5 vehicles per timestep, and $w/u_f = 1/2$. Each time step, 2 vehicles wish to enter the link, and will do so if the receiving flow can accommodate. There is a traffic signal at the downstream end of the link. During time steps 0–9, and from time step 50 onwards, the light is green and all of the link’s sending flow can leave. For the other time steps, the light is red, and the sending flow of the link is zero.
- (a) Use the cell transmission model to propagate flow for 80 time steps, portraying the resulting cell occupancies in a time-space diagram (time on the horizontal axis, space on the vertical axis). At what time interval does the receiving flow first begin to drop; at what point does it reach its minimum value; and what is that minimum value? Is there any point at which the entire link is at jam density?
- (b) Repeat, but with $w/u_f = 1$.
- (c) Repeat, but instead use the link transmission model (with the same time step) to determine how much flow can enter or leave the link.
30. [68] Consider the network in Figure 10.34. The figure shows each link’s length, capacity, jam density, free-flow speed, and backward wave speed. The inflow rate at node A is 4320 veh/hr for $0 \leq t < 30$ (t measured in seconds), 8640 veh/hr for $25 \leq t < 120$, and 0 veh/hr thereafter. The splitting proportion towards node C is $\frac{2}{3}$ for $0 \leq t < 50$, $\frac{1}{6}$ for $50 \leq t < 90$, and $\frac{1}{2}$ for $t \geq 90$.
- (a) Use a point queue model to propagate the vehicle flow with the time step $\Delta t = 5$ s. Plot the turning movement flows y_{ABC} , y_{ABD} , y_{BDE} , and y_{BCE} from $t = 0$ until the last vehicle has left the network. (q_{12}

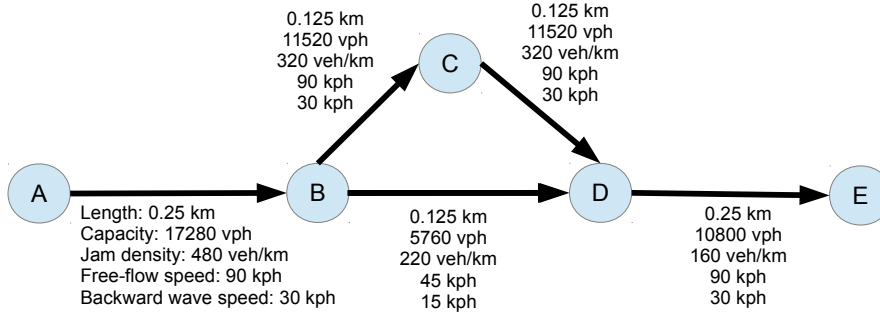


Figure 10.34: Network for Exercise 30

is the rate at which flow leaves the downstream end of link 1 to enter the upstream end of link 2).

- (b) Use the cell transmission model to propagate the vehicle flow with the time step $\Delta t = 5$ s. Plot the same flow rates as in the previous part.
 - (c) Use the link transmission model to propagate the vehicle flow with the time step $\Delta t = 10$ s. Plot the same flow rates as in the previous part.
 - (d) Comment on any differences you see in these plots for the three flow models.
31. [21] Assuming that a cell initially has between 0 and \bar{n} vehicles, show the cell transmission model formula (10.55) ensures that it will have between 0 and \bar{n} vehicles at all future time steps, regardless of upstream or downstream conditions.
 32. [21] On a link, we must have $N^\uparrow(t) \geq N^\downarrow(t)$ at all time steps. Assuming this is true for all time steps before t , show that the link transmission model formulas (10.63) and (10.65) ensure this condition holds at t as well.
 33. [42] Generalize the cell transmission model formula (10.55) to handle an arbitrary piecewise-linear fundamental diagram (not necessarily triangular or trapezoidal).
 34. [65] Generalize the link transmission model formulas (10.63) and (10.65) to handle an arbitrary piecewise-linear fundamental diagram.
 35. [35] Figure 10.35 represents the intersection of Lamar and Guadalupe, showing sending and receiving flows, saturation flows, turning movement proportions, and the signal timing plan (assume no lost time due to clearance intervals or startup delay). Note that the receiving flow on north-bound Lamar is quite low, because of congestion spilling back from a

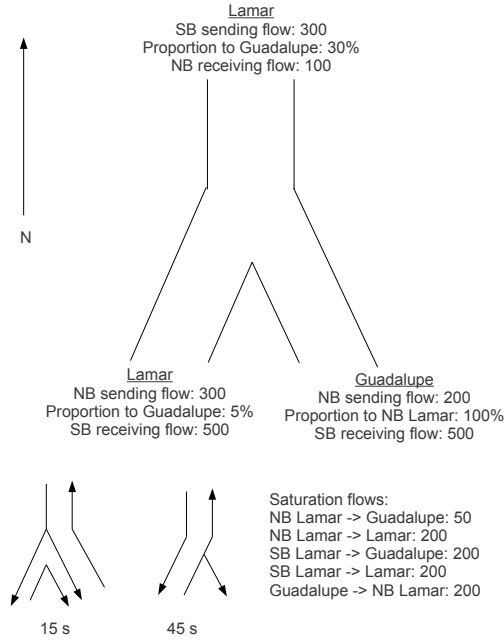


Figure 10.35: Intersection for Exercise 35

nearby signal just downstream. No U-turns are allowed, and drivers may not turn left from Guadalupe onto southbound Lamar.

- (a) Find the transition flows y_{ijk} for all five turning movements at the current time step, using the “smoothed signal” node model.
 - (b) The southbound receiving flow on Guadalupe is now reduced to 50 due to congestion further downstream. Find the updated transition flow rates for all turning movements.
36. [51] Extend the “basic signal” node model to account for turns on red, where a vehicle facing a red indication may make a turn in the direction nearest to them (usually right-on-red in countries that drive on the right, left-on-red in countries that drive on the left). Vehicles turning on red must yield to traffic which has a green indication.
 37. [53] Extend the “basic signal” node model to the case where there are turn lanes, and not all turn lanes from an approach have the same green time.
 38. [61] Extend the “basic signal” node model to account for permitted turns on green (a turning movement which has a green indication, but must

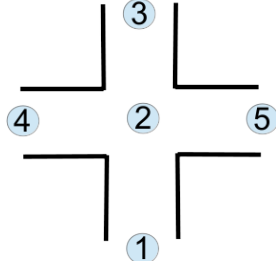


Figure 10.36: Intersection for Exercise 39.

yield to oncoming traffic).

39. [36] Consider the four-legged intersection shown in Figure 10.36. The set of turning movements is $\Xi = \{[1, 2, 3], [1, 2, 4], [3, 2, 1], [3, 2, 4], [5, 2, 4]\}$. For the current time step, the sending and receiving flow values are $S_{12} = 6$, $S_{32} = 6$, $S_{52} = 1$, $R_{12} = 14$, $R_{23} = 14$, and $R_{24} = 14$. Half of the drivers approaching from link (1,2) want to turn left, and half want to go straight. Half of the drivers approaching from link (3,2) want to turn right, and half want to go straight. Finally, the intersection geometry and signal timing are such that the capacities of the incoming links are $q_{12}^{max} = 200$, $q_{32}^{max} = 2000$, and $q_{52}^{max} = 20$.
Apply the “equal priorities” algorithm, and report the transition flows (y values) for each turning movement.
40. [36] Repeat Exercise 39, but with these values of the sending and receiving flows: $S_{12} = 18$, $S_{32} = 18$, $S_{52} = 3$, $R_{12} = 42$, $R_{23} = 42$, and $R_{24} = 12$.
41. [68] Modify the “partial stop control” node model to allow cases of absolute priority (if movement $[h, i, j]$ has absolute priority over $[h', i, j']$ at conflict point c , then $\alpha_{h'ij'}^c / \alpha_{hij}^c$ would be zero.) You will need to decide how priority will be granted for any combination of sending flows wishing to use a conflict point (including flows of equal absolute priority, and when not all sending flows are present), and your formulas can never divide by zero.
42. [89] Show that all of the node models in Section 10.6 satisfy all of the desiderata from Section 10.2.

Chapter 11

Time-Dependent Shortest Paths

This chapter discusses how travelers make choices when traveling in networks whose state varies over time. Two specific choices are discussed: how drivers choose a route when link costs are time-varying (Sections 11.1 and 11.2), and how drivers choose a departure time (Section 11.3). This chapter is the complement of the previous one. In network loading, we assumed that the travelers' choices were known, and we then determined the (time-varying) flow rates and congestion pattern throughout the network. In this chapter, we take this congestion pattern as known, and predict the choices that travelers would make given this congestion pattern. In particular, by taking the congestion level as fixed, we can focus the question on an *individual* traveler and do not need to worry about changes in congestion based on these choices just yet.

11.1 Time-Dependent Shortest Path Concepts

The *time-dependent shortest path* problem involves finding a path through a network of minimum cost, when the cost of links varies with time. As with the static shortest path problem, the “cost” of a link can include travel time, monetary costs, a combination of these, or any other disutility which can be added across links; and by a “shortest” path we mean one with least cost. If a traveler enters link (i, j) at time t , they will experience a cost of $c_{ij}(t)$. Unlike the static shortest path problem, however, we must always keep track of the travel time on a link even if the cost refers to a separate quantity: because the network state is dynamic, we must always know the time at which a traveler enters a link. The travel time experienced by a traveler entering link (i, j) at time t is denoted by $\tau_{ij}(t)$. This problem can be formulated either in discrete time (where t and $\tau_{ij}(t)$ are limited to being integer multiples of the timestep Δt) or in continuous time (where t can take any real value within a stated range).

There are several different variations of the time-dependent shortest path problem, all stemming from dynamic link costs and travel times. One can consider time-dependent shortest paths where *waiting* at intermediate nodes is allowed, or one can forbid it. In public transit networks, waiting at intermediate nodes is logical, but in road networks, one would not expect drivers to voluntarily stop and wait in the middle of a route. One can also restrict attention to problems where the link travel times satisfy the *first-in, first-out* (FIFO) property, where it is impossible for one to leave a link earlier by entering later, that is, for any link (i, j) and any distinct $t_1 < t_2$, we have

$$t_1 + \tau_{ij}(t_1) \leq t_2 + \tau_{ij}(t_2). \quad (11.1)$$

If this is true, more efficient algorithms can be developed. As an example, if the cost of a link is its travel time ($c_{ij}(t) = \tau_{ij}(t)$), there is no benefit to waiting in a FIFO network. In discrete time, equation (11.1) is equivalent to requiring

$$\tau_{ij}(t + \Delta t) - \tau_{ij}(t) \leq \Delta t, \quad (11.2)$$

that is, a link's travel time cannot decrease by more than Δt in one time step. In continuous time, if $\tau_{ij}(t)$ is a piecewise differentiable function, we must have

$$-\frac{d}{dt}\tau_{ij}(t) \leq -1, \quad (11.3)$$

everywhere that $\tau_{ij}(t)$ has a derivative, which expresses the same idea.

One can also distinguish time-dependent shortest path problems by whether the departure time is fixed, whether the arrival time is fixed, or neither. In the first case, the driver has already decided when they will leave, and want to find the route to the destination with minimum cost when leaving at that time, regardless of the arrival time at the destination. In the second case, the arrival time at the destination is known (perhaps the start of work), and the traveler wants to find the route with minimum cost arriving at the destination at that time (regardless of departure time). In the third case, both the departure and arrival times are flexible (as with many shopping trips), and the traveler wants to find a minimum cost route without regard to when they leave or arrive. In this way, we can model the departure time choice decision simultaneously with the route choice decision. Section 11.3 develops this approach further, showing how we can incorporate penalty “costs” associated with different departure and arrival times. Strictly speaking, one can imagine a fourth variant where *both* departure and arrival time are fixed, but this problem is not always well-posed; there may be no path from the origin to the destination with those exact departure and arrival times. In such cases, we can allow free departure and/or arrival times, but severely penalize departure/arrival times that differ greatly from the desired times.

In comparison with the static shortest path problem, the fixed departure time variant is like the one origin-to-all destinations shortest path problem, and the fixed arrival time variant is like the all origins-to-one destination shortest path problem.

To be specific, in this chapter, we will first focus on time-dependent shortest paths with fixed departure times and free arrival times, in Section 11.2, but all of these approaches can be adapted to solve the fixed arrival time/free departure time version without much difficulty. The origin r and departure time t_0 will be specified, and we will find the paths to all other nodes along which the sum of the costs c_{ij} is minimal, given the changes in the costs which will occur during travel. We will study two specific variants of the problem. In the first, time is continuous, but the network is assumed to follow the FIFO principle and the link costs must be the travel times. In the other variants, we do not require the FIFO assumption, and the link costs need not be the same as the travel times; but in exchange we will restrict ourselves to discrete time. Section 11.3 will then treat the case of free departure and arrival times, where only the origin r will be specified.

11.1.1 Time-expanded networks

For discrete-time shortest path problems, we can form what is known as the *time-expanded network*. This technique transforms a time-dependent shortest path problem into a static shortest path problem, that can be solved using algorithms for this simpler problem (see Section 2.4). However, the time-expanded network contains many more links and nodes than the original network, and can impose computational burdens.

If the time step is Δt , then the possible arrival times at any node are $0, \Delta t, 2\Delta t, \dots, T\Delta t$ where T is the time horizon under consideration. In the time-expanded network, we replicate each node $T + 1$ times, one for each possible arrival time. So, if the original network has n nodes, the time-expanded network has nT nodes. Each node in the time-expanded network is written in the form $i : t$, denoting the physical node i at the time interval $t\Delta t$. This labeling is shown in Figure 11.1. The original network is called the *physical network* when we need to distinguish it from the corresponding time-expanded network.

Arcs in the time-expanded network represent both the physical connection between links, as well as the time required to traverse the link. In particular, for each link (i, j) in the original network, and for each time interval $t\Delta t$, if the travel time is $\tau_{ij}(t\Delta t) = k\Delta t$, we create a link $(i : t, j : (t + k))$ in the time expanded network, assuming that $t + k \leq T$, and set the cost of this link equal to $c_{ij}(t)$. A little bit of care must be taken when links would arrive at a node later than the time horizon. If the time horizon is large enough, this should not be a significant issue. In this chapter, we will not create a time-expanded link if its head node falls outside of the time horizon. Exercises 4 and 5 introduce two other ways to treat boundary issues associated with the time horizon.

The advantage of the time-expanded network is that it reduces the time-dependent shortest path problem to the static one: solving the one-to-all static shortest path problem from node i at time t_0 corresponds exactly to solving the fixed-departure time-dependent shortest path problem. Furthermore, if all link travel times are positive, the time-expanded network is acyclic, with the time labels forming a natural topological order. Shortest paths on acyclic networks

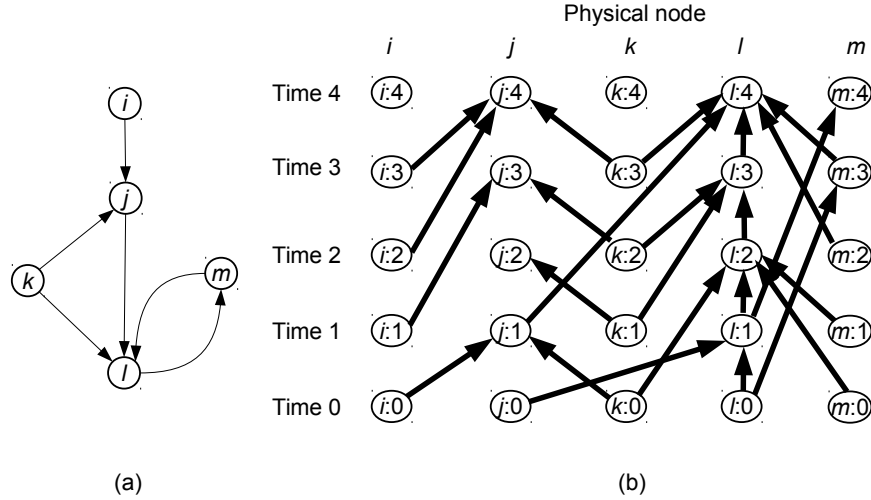


Figure 11.1: (a) Original network (b) Corresponding time-expanded network with five physical nodes and five time intervals.

can be solved rather quickly, so this is a significant advantage. Even if some link travel times are zero (as may occur with some artificial links or centroid connectors), the time-expanded network remains acyclic unless there is a *cycle* of zero-travel time links in the network; and if that is the case, it is often possible to collapse the zero-travel time cycle into a single node. Time-expanded networks can also unify some of the variants of the time-dependent shortest path problem. If waiting is allowed at a node i , we can represent that with links of the form $(i : t, i : (t + 1))$ connecting the same physical node to itself, a time step later. The FIFO principle in a time-dependent network means that two links connecting the same physical nodes will never cross (although they may terminate at the same node).

A disadvantage is that the number of time intervals T may be quite large. In dynamic traffic assignment, network loading often uses a time step on the order of a few seconds. If this same time step is used for time-dependent shortest paths, a typical planning period of a few hours means that T is in the thousands. Given a physical network of thousands of nodes, the time-expanded network can easily exceed a million nodes. Even though the time-expanded network is acyclic, this is a substantial increase in the underlying network size, increasing both computation time and amount of computer memory required. With clever implementations, it is often possible to avoid explicitly generating the entire time-expanded network, and only generate links and nodes as needed.

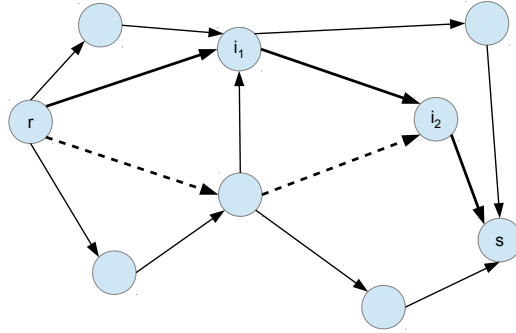


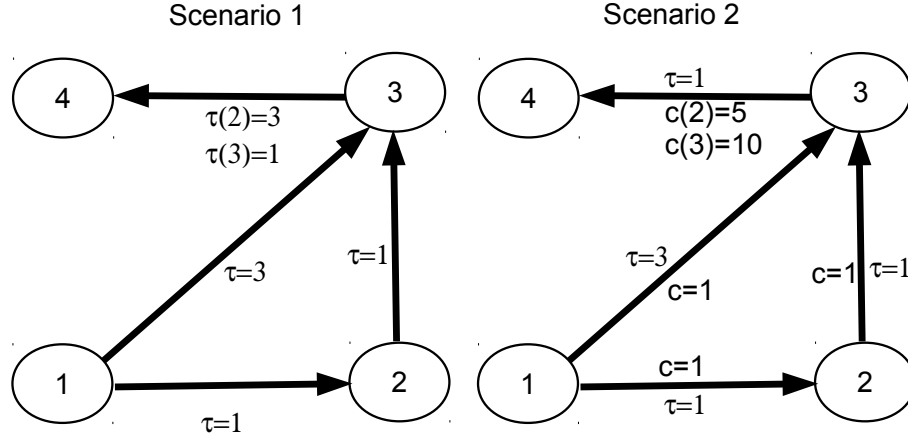
Figure 11.2: Bellman's principle illustrated. If the dashed path were a shorter route from s to i_2 , then it would also form a “shortcut” for the route from s to d .

11.1.2 Bellman's principle in time-dependent networks

The number of paths between any two points in a network can be very large, growing exponentially with the network size. Therefore, any approach based on enumerating all paths and comparing their cost will not scale well to realistic networks. So, efficient algorithms for finding shortest paths, whether time-dependent or not, must be more clever. The key insight is called *Bellman's principle*. In the static shortest path problem, Bellman's principle states that any segment of a shortest path between two nodes must itself be a shortest path between the endpoints of that segment. Otherwise, the shortest path between the endpoints of the segment could be spliced into the original path, reducing its cost (Figure 11.2).

In a time-dependent shortest path problem, the same general idea applies but we must be slightly more careful about how the principle is defined. In a FIFO network where link costs are equal to the travel time, the principle holds identically: since it is always better to arrive at nodes as soon as possible, any “shortcut” between two nodes in a path can be spliced into the full path, thereby reducing its total travel time. When the FIFO principle does not hold, or if link costs are different than travel time, this may not be true. Figure 11.3 shows two counterexamples. In the first, because the FIFO principle is violated and waiting is not allowed, the path segment $[1, 2, 3]$ has a shorter travel time between nodes 1 and 3 than path $[1, 3]$. However, when the link $(3, 4)$ is added to these paths, the path $[1, 2, 3, 4]$ has a higher travel time than path $[1, 3, 4]$ because link $(3, 4)$ is entered at different times, and so its travel time is different. In the second, path $[1, 3]$ has a lower cost than path $[1, 2, 3]$, even though path $[1, 3, 4]$ has higher cost than path $[1, 2, 3, 4]$. In the second case, this is true even though arriving at any node later can only increase cost, and the same result holds even if waiting is allowed.

In these cases, the correct approach is to apply Bellman's principle to the time-expanded network, an approach which also works in FIFO networks. Bell-



Costs and times are constant when no time index is given.

Figure 11.3: Two counterexamples to the naïve application of Bellman's principle in time-expanded networks.

man's principle *does* apply to the time-expanded network, corresponding to the following principle in the physical network: **Let π be a shortest path between nodes r and s when departing at time t_0 . If i and h are two nodes in this path, and if the arrival times at these nodes are t_i and t_h , respectively, then the segment of π between i and h must be a shortest path between these nodes when departing at t_i and arriving at t_h .** All of the algorithms in this chapter implicitly use this principle, by allowing us to construct shortest paths a single link at a time: if we already know the shortest path π_i from the origin at t_0 to some other node i at time t_i , then any other shortest path passing through node i at time t_i can be assumed to start with π_i .

11.2 Time-Dependent Shortest Path Algorithms

This section provides two algorithms for the time-dependent shortest path problem with fixed departure times. In the first, time can be modeled as either discrete or continuous, but the network must satisfy the FIFO principle and the cost of a link must be its travel time. In the second and third, time must be discrete, but FIFO need not hold and the cost of a link may take any value. The second algorithm finds the time-dependent shortest path from a single origin and departure time to all destinations, while the third algorithm finds the time-dependent shortest paths from a single origin and *all* departure times to a single

destination. There are many other possible variants of time-dependent shortest path algorithms, some of which are explored in the exercises — in particular, Exercise 6 asks you to develop a time-dependent shortest path algorithm for all origins and departure times simultaneously, which often arises in dynamic traffic assignment software. Nevertheless, the algorithms here should give the general flavor of how they function. Which one of these algorithms is best, or whether another variant is better, depends on the particular dynamic traffic assignment implementation. All of these algorithms use labels with similar (or even identical) names, but the meanings of these labels are slightly different in each.

11.2.1 FIFO networks

For this algorithm to apply, assume that the FIFO principle holds, in either its discrete or continuous form. Also assume that the cost of a link is simply its travel time, so $c_{ij}(t) = \tau_{ij}(t)$ for all links (i, j) and times t . We are also given the origin r and departure time t_0 . Because the FIFO principle holds, waiting at an intermediate node is never beneficial, and so it suffices to find the earliest time we can reach a node and ignore all later times. This allows us to adapt Dijkstra's algorithm (Section 2.4.3) for the static shortest path problem, a label-setting approach. Two labels are defined for each node: L_i gives the earliest possible arrival time to node i found so far, when departing origin r at time t_0 . The backnode label q_i gives the previous node in a path corresponding to this earliest known arrival time. By the dynamic version of Bellman's principle in FIFO networks, this suffices for reconstructing the shortest path from r to i . For nodes i where we have not yet calculated the earliest possible arrival time, we will set L_i to ∞ , and for nodes i where the backnode is meaningless (either because we have not yet found a path there, or because i is the origin), we set q_i to -1 . We also maintain a set of *finalized nodes* F . If node i belongs to F , this means that we are sure that we have found the earliest possible arrival time to i over all possible paths.

The algorithm functions as follows:

1. Initialize by setting $L_i \leftarrow \infty$ for all nodes i , except for the origin, where $L_r \leftarrow t_0$.
2. Initialize the set of finalized nodes $F \leftarrow \emptyset$, and the backnode vector $\mathbf{q} \leftarrow -\mathbf{1}$.
3. Choose an unfinalized node i with the lowest L_i value. (At the first iteration, this will be the origin r .)
4. Finalize node i : $F \leftarrow F \cup \{i\}$.
5. For each link $(i, j) \in \Gamma(i)$ such that $L_i + \tau_{ij}(L_i)$ is within the time horizon, perform the following steps:
 - (a) Update $L_j \leftarrow \min\{L_j, L_i + \tau_{ij}(L_i)\}$.

- (b) If L_j changed in the previous step, update $q_j \leftarrow i$.
6. If all nodes are finalized ($F = N$), terminate. Otherwise, return to step 3.

As an example of this algorithm, consider the network in Figure 11.4. The time-dependent travel times are shown in this figure. The FIFO assumption is satisfied: for links (1, 3) and (2, 4) the travel times are constant; for links (1, 2) and (3, 4) the travel times are increasing (so arriving earlier always means leaving earlier); and for link (2, 3) the travel time is decreasing, but at a slow enough rate that you cannot leave earlier by arriving later, cf. (11.3). Assume that the initial departure time from node 1 is at $t_0 = 2$, and that the time horizon is large enough that $L_i + \tau_{ij}(L_i)$ is always within the time horizon whenever step 5 is encountered. The steps of the algorithm are explained below, and summarized in Table 11.1. This table shows the state of the algorithm *just before* step 3 is performed, counting the first time through as iteration zero.

Initially, no nodes are finalized, all cost labels are initialized to ∞ (except for the origin, which is assigned 2, the departure time), and all backnode labels are initialized to -1 . The unfinalized node with the least L value is node 1, which is selected as the node i to scan. At the current time of 2, link (1, 2) has a travel time of 6, and link (1, 3) has a travel time of 10. Following these links would result in arrival at nodes 2 and 3 at times 8 and 12, respectively. Each of these is less than their current values of ∞ , so the cost and backnode labels are adjusted accordingly. At the next iteration, node 2 is the unfinalized node with the least L value, so $i = 2$. At time 8, link (2, 3) has a travel time of 1, and link (2, 4) has a travel time of 5. Following these links, one would arrive at nodes 3 and 4 at times 9 and 13, respectively. Both of these values are less than the current L values for these nodes, so their L and q labels are changed. At the next iteration, node 3 is the unfinalized node with the least L value, so $i = 3$. At this time, link (3, 4) would have a travel time of $4\frac{1}{2}$, and choosing it means arriving at node 4 at time $13\frac{1}{2}$. This is greater than the current value ($L_4 = 13$), so no labels are adjusted. Finally, node 4 is chosen as the only unfinalized node. Since it has no outgoing links ($\Gamma(4)$ is empty), there is nothing to do in step 5, and since all nodes are finalized the algorithm terminates.

At this point, we can trace back the shortest paths using the backnode labels: the shortest paths to nodes 2, 3, and 4 are $[1, 2]$, $[1, 2, 3]$, and $[1, 2, 4]$, respectively; and following these paths one arrives at the nodes at times 8, 9, and 13.

11.2.2 Discrete-time networks, one departure time

This algorithm applies in any discrete-time network, regardless of whether or not the FIFO principle holds, and regardless of whether $c_{ij}(t) = \tau_{ij}(t)$ or not. We are given the origin r and departure time t_0 , and work in the time-expanded network. Each node $i : t$ in the time-expanded network is associated with two labels. The label L_i^t denotes the cost of the shortest path from r to i known so far, when departing at time t_0 and arriving at time t . The backnode label q_i^t provides the previous node on a path corresponding to cost L_i^t . As before,

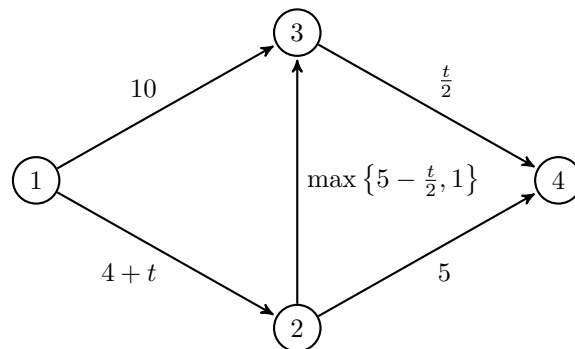


Figure 11.4: Network and travel time functions for demonstrating the FIFO time-dependent shortest path algorithm.

Table 11.1: FIFO time-dependent shortest path algorithm for the network in Figure 11.4, departing node 1 at $t_0 = 2$.

Iteration	F	i	L_1	L_2	L_3	L_4	q_1	q_2	q_3	q_4
0	\emptyset	—	2	∞	∞	∞	-1	-1	-1	-1
1	{1}	1	2	8	12	∞	-1	1	1	-1
2	{1, 2}	2	2	8	9	13	-1	1	2	2
3	{1, 2, 3}	3	2	8	9	13	-1	1	2	2
4	{1, 2, 3, 4}	4	2	8	9	13	-1	1	2	2

$L_i^t = \infty$ signifies that no path to i at time t is yet known, and $q_i^t = -1$ signifies that the backnode is meaningless. Since the FIFO principle may not hold, it is not always advantageous to arrive at a node as early as possible. To reflect this, we work in the time-expanded network, where we can naturally identify the best time to arrive at nodes. We presume that all links have strictly positive travel time, so that the time-expanded graph is acyclic and each link connects a node of earlier time to a node of later time. If there are zero-travel time links in the physical network, but no cycles of zero-travel time links, then step 4 is assumed to proceed in topological order by physical node. In this case, we can adapt an algorithm for finding static shortest paths (e.g., from Section 2.4), in acyclic networks by applying it to the time-expanded network in the following way:

1. Initialize $\mathbf{q} \leftarrow -\mathbf{1}$ and $\mathbf{L} \leftarrow \infty$.
2. Set $L_r^{t_0} \leftarrow 0$.
3. Initialize the current time to the departure time, $t \leftarrow t_0$
4. For each time-expanded node $i : t$ for which $L_i^t < \infty$, and for each time-expanded link $(i : t, j : t')$, perform the following steps:
 - (a) Set $L_j^{t'} = \min \{L_j^{t'}, L_i^t + c_{ij}(t)\}$.
 - (b) If $L_j^{t'}$ changed in the previous step, update $q_j^{t'} \leftarrow i : t$.
5. If $t = T$, then terminate. Otherwise, move to the next time step ($t \leftarrow t+1$) and return to step 4.

At the conclusion of this algorithm, we have the least-cost paths for each possible arrival time at each destination. To find the least-cost path to a particular destination s (at any arrival time), you can consult the L_s^t labels at all times t , and trace back the path for the arrival time t with the least L_s^t value.

This algorithm is demonstrated on the network in the *right* panel of Figure 11.3, and its progress is summarized in Table 11.2. The table shows the state of the algorithm *just before* step 4 is executed. For brevity, this table only reports L_i^t and q_i^t labels for nodes and arrival times which are reachable in the network (that is, i and t values for which $L_i^t < \infty$ at the end of the algorithm). All other cost and backnode labels are at ∞ and -1 throughout the entire duration of the algorithm.

Initially, all cost labels are set to ∞ and all backnode labels to -1 , except for the origin and departure time: $L_1^0 = 0$. The algorithm then sets $t = 0$, and scans over all physical nodes which are reachable at this time¹ Only node 1 can be reached at this time, and the possible links are $(1 : 0, 2 : 1)$ and $(1 : 0, 3 : 3)$.

¹Reachability is expressed by the condition $L_i^t < \infty$; a node which cannot be reached at this time will still have its label set to the initial value. Any node which *can* be reached at this time will have a finite L_i^t value, since step 4a will always reduce an infinite value.

Table 11.2: Discrete time-dependent shortest path algorithm for the network in the right panel of Figure 11.3, departing node 1 at $t_0 = 0$.

t	L_1^0	L_2^1	L_3^2	L_3^3	L_4^3	L_4^4	q_1^0	q_2^1	q_3^2	q_3^3	q_4^3	q_4^4
0	0	∞	∞	∞	∞	∞	-1	-1	-1	-1	-1	-1
1	0	1	∞	1	∞	∞	-1	1	-1	1	-1	-1
2	0	1	2	1	∞	∞	-1	1	2	1	-1	-1
3	0	1	2	1	7	∞	-1	1	2	1	3	-1
4	0	1	2	1	7	11	-1	1	2	1	3	3

Following either link incurs a cost of 1, which is lower than the (infinite) values of L_2^1 and L_3^3 , so the cost and backnode labels are updated.

The algorithm then sets $t = 1$. Only node 2 is reachable at this time, and the only link is $(2 : 1, 3 : 2)$. Following this link incurs a cost of 1; in addition to the cost of 1 already involved in reaching node 2, this gives a cost of 2 for arriving at node 3 at time 2. The cost and time labels for $3 : 2$ are updated. Since the costs and times are different, notice that arriving at node 3 at a later time (3 vs. 2) incurs a lower cost (1 vs. 2). This is why we need to track labels for different arrival times, unlike the algorithm in the previous section.

The next time step is $t = 2$. Only node 3 is reachable at this time (from the path $[1, 2, 3]$), and the only link is $(3 : 2, 4 : 3)$. Following this link incurs a cost of 5, resulting in a total cost of 7, and the labels for $4 : 3$ are updated. Time $t = 3$ is next, and again only node 3 is reachable at this time — but from the path $[1, 3]$. The only link is $(3 : 3, 4 : 4)$, and following this link incurs a cost of 10, for a total cost of 11. Labels are updated for $4 : 4$. There are no further label changes in the algorithm (all nodes have already been scanned at all reachable times), and it terminates as soon as t is increased to the time horizon.

After termination, the L_3^t labels show that we can reach node 3 either with a cost of 1 (arriving at time 3) or a cost of 2 (arriving at time 2). The least-cost path thus arrives at time 3, and it is $[1, 3]$. The L_4^t labels show that we can reach node 4 either with a cost of 7 (arriving at time 3) or a cost of 11 (arriving at time 4). The least-cost path arrives at time 3, and it is $[1, 2, 3, 4]$. (The least-cost path to node 2 is $[1, 2]$, since there is only one possible arrival time there). Notice that the naïve form of Bellman's principle is not satisfied: the least-cost path to node 2 is *not* a subset of the least-cost path to node 3. This was why we needed to keep track of different possible arrival times to nodes — the least-cost path to node 3 *arriving at time 2* is indeed a subset of the least-cost path to node 4 *arriving at time 3*.

11.3 Departure Time Choice

Travelers often have some flexibility when choosing their departure or arrival times, and may choose to leave earlier or later in order to minimize cost. For example, commuters with flexible work hours may want to time their commutes

to avoid congestion, or when dynamic congestion charges are lower. Both flexible departures and arrivals can be incorporated into the shortest path algorithms described in the previous section — in fact, you may have already noticed that the algorithms in Sections 11.2.1 and 11.2.2 allow flexible arrival times. This section explores these choices more systematically, showing how the departure time can also be made flexible, and ways to represent different departure time behaviors.

11.3.1 Artificial origins and destinations

In the time-expanded network, departure and arrival time choice can be modeled by adding artificial “super-origin” and “super-destination” nodes which reflect the start and end of a trip without regard to the time. A super-origin is connected to time-expanded nodes which correspond to allowable departure times, and a super-destination is connected to time-expanded nodes which correspond to allowable arrival times. Initially, we will assign these links a cost of zero, which means the traveler is indifferent among any of these departure or arrival times. Section 11.3.2 describes how arrival and departure time preferences can be modeled.

Figure 11.5 shows how super-origins and super-destinations can be added to the physical network of Figure 11.1(a), assuming that nodes i and k are origins, and nodes j and l are destinations. The figure is drawn as if departures were allowed only for times 0, 1, and 2, but arrivals are allowed at any time. The network is similar to the time-expanded network in Figure 11.1(b), but now includes four artificial nodes, and artificial links corresponding to allowable departure and arrival times.

Once these artificial links and nodes are added to the network, the algorithm from Section 11.2.2 can be applied directly, with only very minor changes. In what follows, r is the origin and the algorithm finds least-cost paths to all destinations, for any allowable departure and arrival time:

1. Initialize $\mathbf{q} \leftarrow -\mathbf{1}$ and $\mathbf{L} \leftarrow \infty$ for all nodes in the time-expanded network.
2. Set $L_r \leftarrow 0$ for the super-origin r , and for each artificial link $(r, r : t)$ set $L_r^t \leftarrow c_{r,depart}(t)$ and $q_r^t \leftarrow r$, where $c_{r,depart}(t)$ is the cost on the artificial link for departing node r at time t .
3. Initialize the current time t to the earliest possible departure time (the lowest t index for which an artificial link $(r, r : t)$ exists).
4. For each time-expanded node $i : t$ for which $L_i^t < \infty$, and for each time-expanded link $(i : t, j : t')$, perform the following steps:
 - (a) Set $L_j^{t'} = \min \{L_j^{t'}, L_i^t + c_{ij}(t)\}$.
 - (b) If $L_j^{t'}$ changed in the previous step, update $q_j^{t'} \leftarrow i : t$.

Likewise, for each artificial link $(i : t, i)$ reaching a super-destination node i with cost $c_{i,arrive}(t)$, perform the following steps:

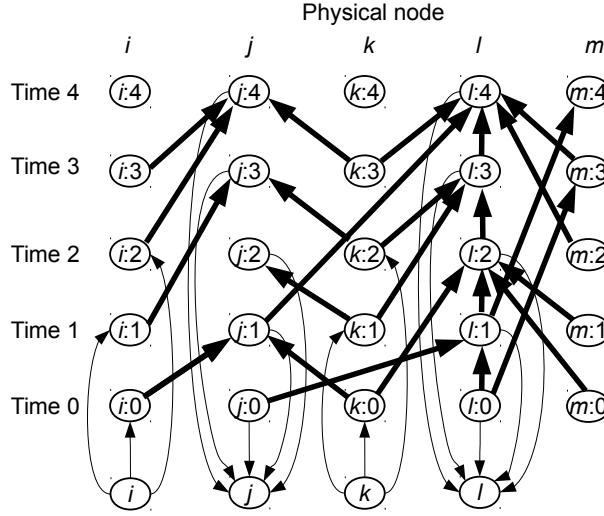


Figure 11.5: Time-expanded network with artificial origins and destinations. Artificial links shown with light arrows for clarity.

- (a) Set $L_i = \min \{L_i, L_i^t + c_{i, arrive}(t)\}$.
 - (b) If L_i changed in the previous step, update $q_i \leftarrow i : t$.
5. If $t = T$, then terminate. Otherwise, move to the next time step ($t \leftarrow t+1$) and return to step 4.

The algorithm initializes labels differently in step 2; step 3 starts at the earliest possible departure time rather than the fixed time t_0 ; and step 4 is expanded to update labels both at adjacent time-expanded nodes and super-destinations. All other steps work in the same way.

To demonstrate this algorithm, consider the network in Figure 11.6, where the time horizon is $T = 20$ and the destination is node 4, and where the cost of a link is equal to its travel time. Table 11.3 shows the cost and forward node labels at the conclusion of the algorithm. Each iteration of the algorithm generates one row of this table, starting with $T = 1$ and working up to $T = 20$. Whenever $L_i^t = \infty$ is seen in Table 11.3, there is no way to reach the destination node within the time horizon, if leaving node i at time t .

Table 11.3 also shows the labels for the super-origin and super-destination, below the labels for the time-expanded nodes. The backnode label for the super-destination tells us that the least-cost path arrives at node 4 at time 6; the q_4^6 label tells us the least-cost path there comes through node 3 at time 1; q_3^1 tells us the least-cost path there comes through node 1 at time 0, and q_1^0 brings us to the super-origin. Therefore, we should depart the origin at time 0, and follow the path $[1, 3, 4]$ to arrive at the destination at time 6, with a total cost of $c_4 = 6$.

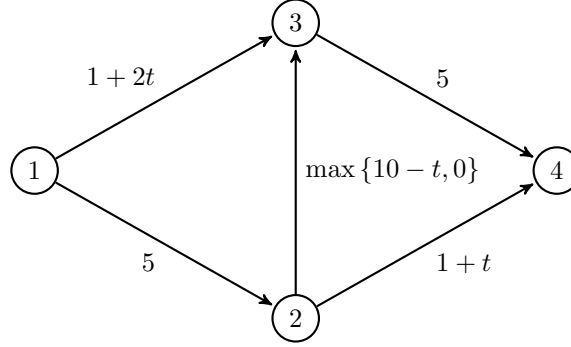


Figure 11.6: Network and travel time functions for demonstrating the all departure times time-dependent shortest path algorithm.

Table 11.3: Labels at conclusion of time-dependent shortest path algorithm with departure time choice, to node 4 in Figure 11.6.

t	L_1^t	L_2^t	L_3^t	L_4^t	q_1^t	q_2^t	q_3^t	q_4^t
20	0	5	5	10	1	1 : 15	2 : 20	3 : 15
19	0	5	5	10	1	1 : 14	2 : 19	3 : 14
18	0	5	5	10	1	1 : 13	2 : 18	3 : 13
17	0	5	5	10	1	1 : 12	2 : 17	3 : 12
16	0	5	5	10	1	1 : 11	2 : 16	3 : 11
15	0	5	5	10	1	1 : 10	2 : 15	3 : 10
14	0	5	5	∞	1	1 : 9	2 : 14	-1
13	0	5	5	12	1	1 : 8	2 : 13	2 : 6
12	0	5	5	10	1	1 : 7	2 : 12	3 : 7
11	0	5	5	11	1	1 : 6	2 : 11	2 : 5
10	0	5	5	∞	1	1 : 5	2 : 10	-1
9	0	5	∞	8	1	1 : 4	-1	3 : 4
8	0	5	∞	∞	1	1 : 3	-1	-1
7	0	5	5	∞	1	1 : 2	1 : 2	-1
6	0	5	∞	6	1	1 : 1	-1	3 : 1
5	0	5	∞	∞	1	1 : 0	-1	-1
4	0	∞	3	∞	1	-1	1 : 1	-1
3	0	∞	∞	∞	1	-1	-1	-1
2	0	∞	∞	∞	1	-1	-1	-1
1	0	∞	1	∞	1	-1	1 : 0	-1
0	0	∞	∞	∞	1	-1	-1	-1

(L, q) labels at super-origin 1: $(0, -1)$

(L, q) labels at super-destination 4: $(6, 4 : 6)$

11.3.2 Arrival and departure time preferences

While there is often flexibility in departure or arrival time, usually travelers are not completely indifferent about when they depart or arrive. For instance, there may be a well-defined arrival deadline (start of work, or check-in time before a flight), and a strong desire to arrive before this deadline. Departing extremely early would ensure arriving before the deadline, but carries opportunity costs (by departing later, the traveler would have more time to do other things). Both situations can be modeled by attaching a cost to when a traveler departs the origin, and to when they arrive at the destination.

A common way to model arrival costs is with the *schedule delay* concept. In this model, travelers have a preferred arrival time t^* at the destination, and arriving either earlier or later than t^* is undesirable. The most general formulation involves a function $f(t)$ denoting the cost (or disutility) of arriving at the destination at time t . This function is typically convex and has a minimum at t^* . One example of such a function is

$$f(t) = \alpha[t^* - t]^+ + \beta[t - t^*]^+, \quad (11.4)$$

where $[\cdot]^+ = \max\{\cdot, 0\}$ expresses the positive part of the quantity in brackets. In (11.4), the first bracketed term then represents the amount by which the traveler arrived early, compared to the preferred time, and the second bracketed term represents the amount by which the traveler arrived late. The coefficients α and β then weight these terms and convert them to cost units; generally $\alpha < \beta$ to reflect the fact that arriving early by a certain amount of time, while undesirable, is usually not as bad as arriving late by that same amount of time.

Another possible function is nonlinear, taking a form such as

$$f(t) = \alpha([t^* - t]^+)^2 + \beta([t - t^*]^+)^2. \quad (11.5)$$

In this function, the penalty associated with early or late arrival grows faster and faster the farther the arrival time from the target. This may occur if, for instance, being ten minutes late is more than ten times as bad as being one minute late. Special cases of these functions arise when $\alpha = \beta$ (the function becomes symmetric), or when $\alpha = 0$ (there is no penalty for early arrival, but only for late arrival). This function is also differentiable everywhere, in contrast to equation (11.4) which is not differentiable at t^* . For certain algorithms this may be advantageous.

To each of these schedule delay functions $f(t)$, one can add the cost of the path arriving at time t (the sum of the link costs c_{ij} along the way) to yield the total cost of travel. We assume that travelers will choose both the departure time and the path to minimize this sum. The algorithm from Section 11.3.1 can be used to find both this ideal departure time and the path. The only change is that artificial links connecting time-expanded destination nodes $(s : t)$ to superdestinations s now have a cost of $f(t)$, rather than zero. At termination, the departure time t_0^* minimizing $L_r^{t_0^*}$ corresponds to the least total cost, and the forwardnode labels trace out the path.

To demonstrate this algorithm, again consider the network in Figure 11.6, but with the arrival time penalty function

$$f(t) = [10 - t]^+ + 2[t - 10]^+, \quad (11.6)$$

which suggests that the traveler wishes to arrive at time 10, and that late arrival is twice as costly as early arrival. As before, the time horizon is $T = 20$ and the destination is node 4, and the cost of each link is equal to its travel time. Table 11.4 shows the cost and backnode labels at the conclusion of the algorithm, which runs in exactly the same way as before except that the artificial destination links $(4 : t, 4)$ have a cost equal to $f(t)$. At the conclusion of the algorithm, we can identify the total cost on the shortest path from node 1 to node 4, now including the penalty for arriving early or late at the destination.

Table 11.4 also shows the labels for the super-origin and super-destination, below the labels for the time-expanded nodes. The backnode label for the super-destination tells us that the least-cost path arrives at node 4 at time 9; the q_4^9 label tells us the least-cost path there comes through node 3 at time 4; q_3^4 tells us the least-cost path there comes through node 1 at time 1, and q_1^1 brings us to the super-origin. Therefore, we should depart the origin at time 1, and follow the path $[1, 3, 4]$ to arrive at the destination at time 9, with a total cost of $c_4 = 9$. Of this cost, 8 units are due to travel time (difference between arrival and departure times), and 1 unit is due to the arrival time penalty from equation (11.6) with $t = 9$.

Departing earlier, at $t = 0$, on the same path would reduce the travel time to 6, but increase the early arrival penalty to 4. The total cost of leaving at $t = 0$ is thus higher than departing one time step later. Leaving at $t = 2$ and following the same path increases the travel time cost to 10. Since this means arriving at time 12, there is a late penalty cost of 4 added, resulting in a total travel cost of 14. By comparing all possible paths and departure times, you can verify that it is impossible to have a total cost less than 9.

Considering costs associated with departure time, rather than arrival time, is done in essentially the same way, by assigning a nonzero cost to the artificial links connecting the super-origin r to the time-expanded nodes $r(t)$. It is thus possible to have only departure time penalties, only arrival time penalties, both, or neither, depending on whether the artificial origin links and artificial destination links have nonzero costs.

11.4 Dynamic A^*

Because the time-expanded network transforms the time-dependent shortest problem into the classical, static shortest path problem from Section 2.4, any of the algorithms described there can be applied. The time-expanded network is acyclic, so it is fastest to use the algorithm from Section 2.4.1 — and indeed this is all that the algorithm in Section 11.2.2 is, using the time labels as a topological order.

Table 11.4: Labels at conclusion of arrival-penalty time-dependent shortest path, to node 4 in Figure 11.6.

t	L_1^t	L_2^t	L_3^t	L_4^t	q_1^t	q_2^t	q_3^t	q_4^t
20	0	5	5	10	1	1 : 15	2 : 20	3 : 15
19	0	5	5	10	1	1 : 14	2 : 19	3 : 14
18	0	5	5	10	1	1 : 13	2 : 18	3 : 13
17	0	5	5	10	1	1 : 12	2 : 17	3 : 12
16	0	5	5	10	1	1 : 11	2 : 16	3 : 11
15	0	5	5	10	1	1 : 10	2 : 15	3 : 10
14	0	5	5	∞	1	1 : 9	2 : 14	-1
13	0	5	5	12	1	1 : 8	2 : 13	2 : 6
12	0	5	5	10	1	1 : 7	2 : 12	3 : 7
11	0	5	5	11	1	1 : 6	2 : 11	2 : 5
10	0	5	5	∞	1	1 : 5	2 : 10	-1
9	0	5	∞	8	1	1 : 4	-1	3 : 4
8	0	5	∞	∞	1	1 : 3	-1	-1
7	0	5	5	∞	1	1 : 2	1 : 2	-1
6	0	5	∞	6	1	1 : 1	-1	3 : 1
5	0	5	∞	∞	1	1 : 0	-1	-1
4	0	∞	3	∞	1	-1	1 : 1	-1
3	0	∞	∞	∞	1	-1	-1	-1
2	0	∞	∞	∞	1	-1	-1	-1
1	0	∞	1	∞	1	-1	1 : 0	-1
0	0	∞	∞	∞	1	-1	-1	-1

(L, q) labels at super-origin 1: $(0, -1)$

(L, q) labels at super-destination 4: $(9, 4 : 9)$

Table 11.5: Link costs for Exercises 1 and 10.

Entry time	Link 1	Link 2	Link 3	Link 4	Link 5
0	5	1	3	1	1
1	6	1	3	6	1
2	4	2	3	5	2
3	5	3	3	4	4
4	7	5	3	3	2
5	8	8	3	2	1

Section 2.4.4 also presented the A^* algorithm, which provides a single path from one origin to one destination, rather than all shortest paths from one origin to all destinations, or all origins to one destination. By focusing on a single origin and destination, A^* can often find a shortest path much faster than a one origin-to-all destinations algorithm. The tradeoff is that the algorithm has to be repeated many times, once for every OD pair, rather than once for every origin or destination. In static assignment, one-to-all or all-to-one algorithms are preferred because, in many cases, running A^* for each OD pair takes more time than running a one-to-all algorithm for each origin.

In dynamic traffic assignment with fixed departure times, however, the number of “origins” in the time-expanded network is multiplied by the number of departure times. If one were to write a full time-dependent OD matrix, the number of entries in this matrix is very large: in a network with 1000 centroids and 1000 time steps, there are 1 *billion* entries, one for every origin, every destination, and departure time. This is much larger than the number of vehicles that will be assigned, so almost every entry in this matrix will be zero. In such cases, A^* can work much better, only being applied to origins, destinations, and departure times with a positive entry in the matrix.

As discussed in Section 2.4.4, an effective estimate for A^* in traffic assignment problems is to use the free-flow travel costs. As a preprocessing step at the start of traffic assignment, you can use an all-to-one static shortest path algorithm to find the least-cost travel cost g_j^s from every node j to every destination s at free flow. For the remainder of the traffic assignment algorithm, you can then use g_j^s as the estimates for A^* . This is quite effective in practical networks.

11.5 Exercises

1. [13] Table 11.5 shows time-dependent costs on five links, for different entry times. Which links have costs satisfying the FIFO principle?
2. [23] Prove that waiting is never beneficial in a FIFO network where link costs are equal to travel time.

Table 11.6: Backnode labels for Exercise 9.

Entry time	Node 1	Node 2	Node 3	Node 4
6	-1	4	3	2
5	-1	4	3	-1
4	-1	-1	-1	-1
3	-1	4	-1	-1
2	-1	4	-1	-1
1	-1	-1	-1	-1

3. [34] Prove that there is an acyclic time-dependent shortest path in a FIFO network where link costs are equal to travel time.
4. [53] In this chapter, we assumed it is impossible to travel beyond the time horizon, by not creating time-expanded links if they arrive at a downstream node after \bar{T} . Another alternative is to assume that travel beyond \bar{T} is permitted, but that travel times and costs stop changing after that point and take constant values. (Perhaps free-flow times after the peak period is over.)
 - (a) Modify the time-expanded network concept to handle this assumption. (The network should remain finite.)
 - (b) Modify the algorithm in Section 11.2.2 to work in this setting.
5. [63] Another way to handle the time horizon is to assume that the link travel times and costs are *periodic* with length \bar{T} . (For instance, \bar{T} may be 24 hours, and so entering a link at $t = 25$ hours would be the same as at $t = 1$ hour.) First repeat Exercise 4 with this assumption. Then prove that the modified algorithm you create will converge to the correct time-dependent shortest paths for any possible departure time.
6. [31] Modify the algorithm in Section 11.3.1 to become an all-to-one algorithm, that finds least-cost paths from all origins and all departure times to one destination (arrival at any time is permitted).
7. [34] Show that the classical form of Bellman's principle (Section 2.4 holds in a FIFO, time-dependent network where the link costs are equal to the link travel time.)
8. [32] Find the time-dependent shortest path when departing node 1 in the network shown in the *left* panel of Figure 11.3, departing at time 0.
9. [23] Tables 11.6 and 11.7 show the backnode and cost labels for a time-dependent shortest path problem, where the destination is node 4, and the time horizon is 6. (Figure 11.7 shows the network topology.) What is the shortest path from node 1 to node 4, when departing at time 1?

Table 11.7: Cost labels for Exercise 9.

Entry time	Node 1	Node 2	Node 3	Node 4
6	0	1	5	5
5	0	2	4	∞
4	0	∞	∞	∞
3	0	2	∞	∞
2	0	1	∞	∞
1	0	∞	∞	∞

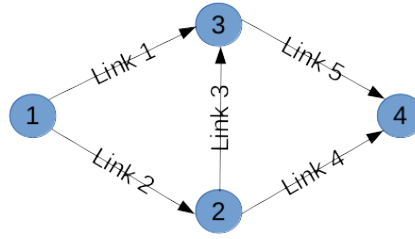


Figure 11.7: Network for Exercises 9 and 10.

10. [22] In the same network as Exercise 9, fill in the forwardnode and cost labels for time 0. The link costs are as in Table 11.5, and waiting is not allowed at nodes.
11. [45] Consider the network in Figure 11.8.
- Verify that the travel times satisfy the FIFO principle.
 - Find the shortest paths between nodes 1 and 4 when departing at $t = 0$, $t = 2$, and $t = 10$.
 - For what departure times would the travel times on paths $[1, 2, 4]$ and $[1, 3, 4]$ be equal?

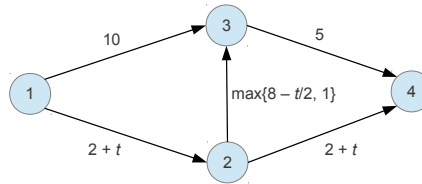


Figure 11.8: Network for Exercise 11.

12. [23] Solve the time-dependent shortest path problem on the network in Figure 11.6 using the algorithm in Section 11.2.1. Comment on the amount of work needed to solve the problem using this algorithm, compared to the way it was solved in the text.
13. [51] In the schedule delay equation (11.4), we typically assume $0 \leq \alpha \leq 1 \leq \beta$ for peak-hour commute trips. What counterintuitive behavior would occur if any of these three inequalities were violated?
14. [12] Find the optimal departure times and paths from nodes 2 and 3 in the network in Figure 11.6, with equation (11.6) as the arrival time penalty. (You can answer this question directly from Table 11.3.)
15. [36] Find the optimal departure times and paths from nodes 1, 2, and 3 in the network in Figure 11.6, if the arrival time penalty function is changed to $f(t) = ([12 - t]^+)^2 + 2([t - 12]^+)^2$.
16. [0] According to the penalty function in Exercise 15, what is the desired arrival time?
17. [88] Implement all of the algorithms described in this chapter, and test them on transportation networks with different characteristics (number of origins, destinations, nodes, links, time intervals, ratio of links to nodes, with or without waiting, FIFO or non-FIFO, etc.). What algorithms perform best under what circumstances? For circumstances most resembling real-world transportation networks, what performs best?

Chapter 12

Dynamic User Equilibrium

The previous two chapters presented separate perspectives on dynamic traffic modeling. Chapter 10 described the network loading problem, in which driver behavior (as represented by path choices and possibly departure times) was known, and where we sought to represent the resulting traffic and congestion patterns on the network. Chapter 11 described the time-dependent shortest path and departure time choice problems, in which the network state (as represented by travel times and costs) was known, and where we sought to identify how drivers would behave. This chapter synthesizes these two perspectives through the concept of dynamic user equilibrium, defined as driver choices and network traffic which are mutually consistent, given the network loading and driver behavior assumptions. The dynamic user equilibrium principle, and how the network loading and driver behavior models can be connected, are the subjects of Section 12.1. Section 12.2 then presents several algorithms that can be used to solve for dynamic equilibria. Such algorithms are almost always heuristics, since realistic network loading models are not amenable to exact analysis, and indeed Section 12.3 shows that dynamic user equilibrium need not exist; and if it exists, it need not be unique. This section also provides examples to show that the dynamic user equilibrium solution need not minimize total travel time in a network, and that providing additional capacity to a network can increase total travel time, a dynamic analogue of the Braess paradox.

12.1 Towards Dynamic User Equilibrium

This section has two major goals, aimed at reconciling the network loading problem of Chapter 10 and the time-dependent shortest path algorithms of Chapter 11. In particular, we will need to solve these problems sequentially, and repeatedly, as shown in Figure 12.1. The first goal of this section is to provide the mechanics to link the output of each problem to the input needed by the other. Following these, we will be equipped to define dynamic user equilibrium formally. Fixed point and variational inequality formulations are

given as well.

12.1.1 Flow Representation

The first question to address is how to represent the choices of all the travelers on the network. In describing network loading, Chapter 10 took an aggregate approach. All of the link models in that chapter can be seen as fluid models, where vehicles are infinitely divisible rather than discrete entities. However, the algorithms in Chapter 11 took a disaggregate approach, and found paths and departure times for a single vehicle. There are several ways to reconcile these perspectives.

The method adopted in this chapter is to store the number of vehicles departing on each path during each time interval. Denote by h_t^π the number of vehicles which start traveling on path π during the t -th interval. These values do not need to be integers; all that matters is that they are nonnegative, and that all paths connecting the same origin and destination sum to the total demand between these zones for each time interval. The h_t^π values can be collected into a single matrix H , whose dimensions are equal to the number of paths and time intervals. With this notation, the feasible set can be written as

$$\bar{H} = \left\{ H \in \mathbb{R}_+^{\bar{T} \times |\Pi|} : \sum_{\pi \in \Pi_{rs}} h_t^\pi = d_t^{rs} \quad \forall (r, s) \in Z^2, t \right\}. \quad (12.1)$$

The main advantage of this notation is that the behavior is clear: by tracking some auxiliary variables in the network loading (as described in Section 12.1.3, at any point in time we can see exactly which vehicles are on which links, and can trace these to the paths the vehicles must follow). These paths can be connected directly with the paths found in a time-dependent shortest path algorithm. A disadvantage of this approach is that the number of paths grows exponentially with the network size. In practice, “column generation” schemes are popular, in which paths π are only identified when found by a shortest path algorithm, and h_t^π values only need be calculated and stored for paths to which travelers have been assigned.

This is not the only possible way to represent travel choices. A link-based representation requires fewer variables. For each turning movement $[h, i, j]$, each destination s , and each time interval t , let $\alpha_{hij,s}^t$ denote the proportion of travelers arriving at node i , via approach (h, i) , during the t -th time interval, who will exit onto link (i, j) *en route* to destination s . This approach mimics the flow splitting rules often found in traffic microsimulation software. The number of variables required by this representation grows with the network size, but not at the exponential rate required by a path-based approach. One can show that the path-based and link-based representation of route choice are equivalent in the sense that there exist α values which represent the same network state as any feasible set of H values, regardless of the network loading model, and that one can identify the H values corresponding to a given set of α values (see Exercise 1). The primary disadvantage of this representation is that the

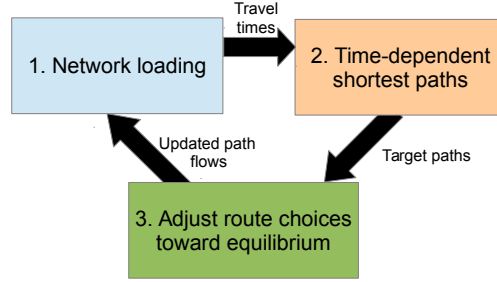


Figure 12.1: Iterative dynamic traffic assignment framework, with details from previous chapters.

behavior is less clear: one cannot trace the path of any vehicle throughout the network deterministically (although one can do so stochastically, making a turn at each node by treating the α values as probabilities, identifying the arrival times at successive nodes using the procedure in Section 12.1.2).

Both of these methods adopt the aggregate, continuous-flow perspective of Chapter 10. Yet another way to represent travel choices is to adopt the individual perspective of the algorithms in Chapter 11, and explicitly model each vehicle as a discrete agent with one specific path. This is the most behavioral way to represent choices: each individual vehicle is assigned one path, with no divisibility or fractional flows. Downsides of this approach are scalability — if the demand doubles but the network topology is unchanged, the individual vehicle approach would require twice as many variables, whereas the continuous methods would not require any more variables — and increased difficulty in implementing the network loading algorithms. It is certainly possible to implement discrete versions of merges, diverges, the cell transmission model, and so forth, but one must be careful about rounding. For instance, if the time step is chosen so that the capacity of a link is less than half a vehicle per timestep, consistently rounding numbers to the nearest integer would mean no vehicles can ever exit. Stochastic rounding, or accumulating a continuous flow value which can then be rounded, can address these difficulties.

12.1.2 Travel time calculation

Once network loading is complete, a time-dependent shortest path algorithm can be applied to determine the cost-minimizing routes for travelers. To do so, we need the travel times $\tau_{ij}(t)$ on each link (i, j) for travelers entering at each time t . The network loading models do not provide this directly, but rather give the cumulative counts N at the upstream and downstream ends of each link, at each time interval. However, this information will suffice for calculating the travel times. For link (i, j) at time t , the upstream count $N_{ij}^{\uparrow}(t)$ gives the total number of vehicles which have entered the link by time t , while the downstream count $N_{ij}^{\downarrow}(t)$ gives the total number of vehicles which left the link by time t .

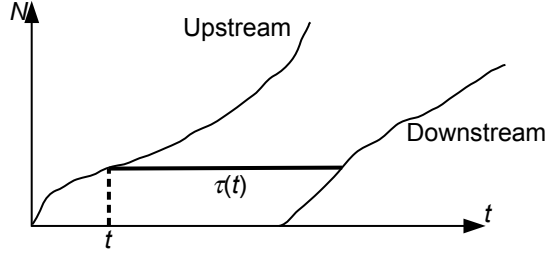


Figure 12.2: Obtaining travel times from plots of the cumulative counts N^\uparrow and N^\downarrow .

Assume for a moment that $N_{ij}^\uparrow(t)$ and $N_{ij}^\downarrow(t)$ are strictly increasing, continuous functions of t . If this is the case, we can define inverse functions $T_{ij}^\uparrow(n)$ and $T_{ij}^\downarrow(n)$, respectively giving the times when the n -th vehicle entered the link and left the link. The travel time for the n -th vehicle is then the difference between these: $T_{ij}^\downarrow(n) - T_{ij}^\uparrow(n)$. Graphically, this can be seen as the horizontal difference between the upstream and downstream N curves. (Figure 12.2). Then, to find the travel time for a vehicle entering the link at time t , we simply evaluate this difference for $n = N_{ij}^\uparrow(t)$:

$$\tau_{ij}(t) = T_{ij}^\downarrow(N_{ij}^\uparrow(t)) - T_{ij}^\uparrow(N_{ij}^\uparrow(t)) = T_{ij}^\downarrow(N_{ij}^\uparrow(t)) - t \quad (12.2)$$

since T_{ij}^\uparrow and N_{ij}^\uparrow are inverse functions.

A little bit of care must be taken because $N_{ij}^\uparrow(t)$ and $N_{ij}^\downarrow(t)$ are not strictly increasing functions of time unless there is always a positive inflow and outflow rate for link (i, j) . Furthermore, we often introduce a time discretization. For both of these reasons, the inverse functions $T_{ij}^\uparrow(n)$ and $T_{ij}^\downarrow(n)$ may not be well-defined. We thus need to modify equation (12.2) in a few ways:

- We must ensure that $\tau_{ij}(t)$ is always at least equal to the free-flow travel time on the link. The danger is illustrated in Figure 12.3, where no vehicles enter or leave the link for an extended period of time. The horizontal distance between the N^\uparrow and N^\downarrow curves at their closest point is small (the dashed line in the figure), but this does not reflect the actual travel time of any vehicle. In reality, a vehicle entering the link when it is completely empty, and when there is no downstream bottleneck, would experience free-flow conditions on the link.
- If the link has no outflow for an interval of time, then $N_{ij}^\downarrow(t)$ will be constant over that interval. This frequently happens with traffic signals. In this case, there are multiple values of time where $N_{ij}^\downarrow(t) = n$. The correct way to resolve this is to define $T_{ij}^\downarrow(n)$ to be the *earliest* time for

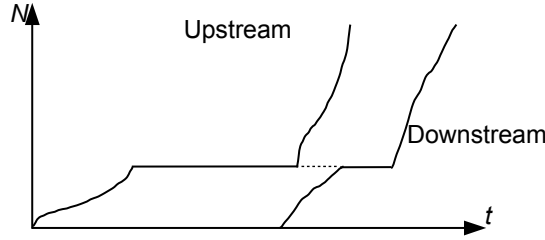


Figure 12.3: Cumulative counts are flat when there is no inflow or outflow.

which $N_{ij}^\downarrow(t) = n$:

$$T_{ij}^\downarrow(n) = \min_t \left\{ t : N_{ij}^\downarrow(t) = n \right\}. \quad (12.3)$$

- In discrete time, the time at which the n -th vehicle departs may not line up with a multiple of Δt , so there may be no known point where $N_{ij}^\downarrow(t)$ is exactly equal to n . In this case, it is appropriate to interpolate between the last time point where $N_{ij}^\downarrow(t) < n$, and the first time point where $N_{ij}^\downarrow(t) \geq n$.

With these modifications to how T_{ij}^\downarrow is calculated, the formula (12.2) can be used to calculate the travel times on each link and arrival time.

The travel time on a path π for a traveler departing at time t , denoted $C^\pi(t)$, can then be calculated sequentially. If the path is $\pi = [r, i_1, i_2, \dots, s]$, then the traveler departs origin r at time t , and arrives at node i_1 at time $t + \tau_{ri_1}(t)$. The travel time on link (i_1, i_2) is then $\tau_{i_1 i_2}(t + \tau_{ri_1}(t))$, so the traveler arrives at i_2 at time $t + \tau_{ri_1}(t) + \tau_{i_1 i_2}(t + \tau_{ri_1}(t))$, and so forth. Writing out this formula can be a bit cumbersome, but calculating it in practice is quite simple: it is nothing more than accumulating the travel times of the links in the path, keeping track of the time at which each link is entered.

12.1.3 Determining splitting proportions

The network loading requires knowledge of when and where vehicles enter the network, and the splitting proportions p at diverges and general intersections. In dynamic traffic assignment, this is reflected indirectly, through the variables $h^\pi(t)$, denoting the number of travelers departing path π at time t . We do not specifically give the proportions p as a function of time, because we do not know when any traveler will reach any node before actually performing the network loading. Path choice is a behavioral parameter associated with travelers, so it is easier to reconcile path choices with behavior if they are expressed this way, rather than by simply specifying the turning fractions at each node.

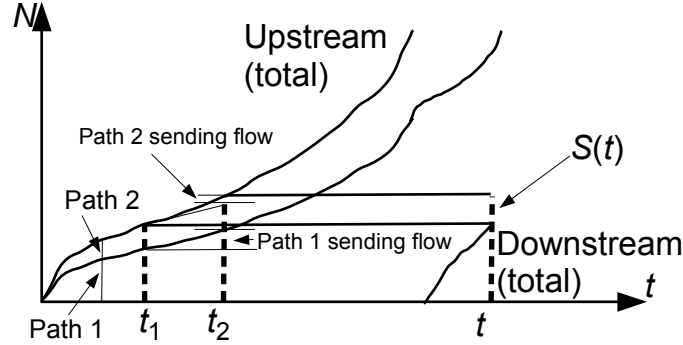


Figure 12.4: Disaggregating the sending flow by path.

That said, the network loading does in fact need $p_{hij}(t)$ values for each turning movement $[h, i, j]$ at a diverge or general intersection, for each time period t . These are obtained by examining the vehicles comprising the sending flow $S_{hi}(t)$, and calculating the fraction of these vehicles whose path includes link (i, j) as the next link. One way to do this is to disaggregate the cumulative count values $N^\uparrow(t)$ and $N^\downarrow(t)$ calculated at the upstream and downstream ends of each link. For each path π in the network, and for every link (h, i) and time t , define $N_{hi,\pi}^\uparrow(t)$ and $N_{hi,\pi}^\downarrow(t)$ to be the total number of vehicles using path π which have respectively entered and left link (h, i) by time t . Clearly we have

$$N_{hi}^\uparrow(t) = \sum_{\pi \in \Pi} N_{hi,\pi}^\uparrow(t) \quad N_{hi}^\downarrow(t) = \sum_{\pi \in \Pi} N_{hi,\pi}^\downarrow(t) \quad \forall (h, i) \in A, t. \quad (12.4)$$

Then, the sending flow for each link and time interval can be disaggregated in the same way, with $S_{hi,\pi}(t)$ defined as the number of vehicles in the sending flow which are using path π . Assume that we are calculating sending flow for link (h, i) at time t , and have determined $S_{hi}(t)$. At this point in time, the total number of vehicles which has left this link is $N_{hi}^\downarrow(t)$. Therefore, the vehicles in the sending flow are numbered in the range $N_{hi}^\downarrow(t)$ to $N_{hi}^\downarrow(t) + S_{hi}(t)$. Using the inverse functions (12.2), the times at which these vehicles entered link (h, i) are in the range $T^\uparrow(N_{hi}^\downarrow(t))$ to $T^\uparrow(N_{hi}^\downarrow(t) + S_{hi}(t))$. Denote these two times by t_1 and t_2 , respectively. Then the disaggregate sending flow $S_{hi,\pi}(t)$ is the number of vehicles on path π which entered link (h, i) between t_1 and t_2 , that is,

$$S_{hi,\pi}(t) = N_{hi,\pi}^\uparrow(t_2) - N_{hi,\pi}^\uparrow(t_1). \quad (12.5)$$

This procedure is illustrated in Figure 12.4.

The proportion of travelers in the sending flow S_{hi} wishing to turn onto link (i, j) can now be calculated. Let $\Pi_{[h,i,j]}(t)$ be the set of paths which arrive at

node (h, i) at time t and immediately continue onto link (i, j) . Then

$$p_{hij}(t) = \frac{\sum_{\pi \in \Pi_{[h,i,j]}(t)} S_{hi,\pi}(t)}{S_{hi}(t)} \quad (12.6)$$

for use in diverge or general intersection node models. Then, after a node model produces the actual flows $y_{hij}(t)$ between links, the disaggregated cumulative counts must be updated as well. Let $y_{hi}(t) = \sum_{(i,j) \in \Gamma(i)} y_{hij}(t)$ be the total flow which leaves link (h, i) during time t . Then all vehicles which entered the link between t_1 and $T_{hi}^\uparrow(N^\downarrow(t)_{hi} + y_{hi})$ can leave the link (use t_3 to denote this latter time), so we update the downstream aggregate counts to

$$N_{hi,\pi}^\downarrow(t + \Delta t) = N_{hi,\pi}^\uparrow(t_3) \quad (12.7)$$

and the upstream counts to

$$N_{ij,\pi}^\uparrow(t + \Delta t) = N_{hi,\pi}^\downarrow(t + \Delta t). \quad (12.8)$$

where (h, i) is the link immediately upstream of (i, j) on path π .

This formula is an approximation, since it assumes that the vehicles on different paths in the sending flow are uniformly distributed. In reality, vehicles towards the start of the sending flow may have a different mix of paths than vehicles towards the end of the sending flow. This means that the p value (calculated based on the entire sending flow) and the actual vehicles which are moved (which are from the start of the sending flow) may not be entirely consistent. In practice, this violation is typically small. If one wishes to obtain an exact formula, an iterative process can be used to ensure that the proportions p are consistent with the proportion of vehicles between t_1 and t_3 .

12.1.4 Principle of dynamic user equilibrium

With the procedures in the previous two subsections, we can now iteratively perform network loading and find time-dependent shortest paths. The dynamic traffic assignment problem is to find a mutually consistent solution between these two models. Let H be a matrix of path flows indicating the number of vehicles departing on each path at each time, and let T be a matrix of path travel times, indicating the travel time for each path and departure time. Then the network loading can be concisely expressed by

$$T = \mathcal{N}(H) \quad (12.9)$$

where \mathcal{N} is a function encompassing whatever network loading is performed, mapping the path flows to path travel times. Likewise, let the assumed behavioral rule be denoted by \mathcal{B} , which indicates the allowable path flow matrices H when the travel times are T . For instance, one may assume that departure times are fixed, but only minimum-travel time paths must be used; that minimum-travel time paths must be used but the departure times must minimize total

cost (travel time plus schedule delay); or some other principle. In general, there may be multiple H matrices which satisfy this principle, as when multiple paths have the same, minimal travel time, in which case \mathcal{B} actually denotes a *set* of H matrices. So, we can express the behavioral consistency rule by

$$H \in \mathcal{B}(T). \quad (12.10)$$

We are now in a position to define the solution to the dynamic traffic assignment problem as a fixed point problem: *the path flow matrix H is a dynamic user equilibrium if*

$$H \in \mathcal{B}(\mathcal{N}(H)). \quad (12.11)$$

That is, the path flow matrix must be consistent with the driver behavior assumptions, when the travel times are obtained from that same path flow matrix. The most common behavioral rule is that departure times are fixed, but only minimum-travel time paths will be used. In this case the principle can be stated more intuitively as *all used paths connecting the same origin to the same destination at the same departure time must have equal and minimal travel time*. If departure times can be varied to minimize total cost, then the principle can be stated as *all used paths connecting the same origin to the same destination have the same total cost, regardless of departure time*.

The dynamic user equilibrium solution can also be stated as a solution to a variational inequality. For the case of fixed departure times, the set of feasible path flow matrices is given by

$$\bar{H} = \left\{ H \in \mathbb{R}_+^{\bar{T} \times |\Pi|} : \sum_{\pi \in \Pi_{rs}} h_t^\pi = d_t^{rs} \quad \forall (r, s) \in Z^2, t \right\}, \quad (12.12)$$

that is, matrices with nonnegative entries where the sum of all path flows connecting the origin r to destination s at departure time t is the corresponding value in the time-dependent OD matrix d_t^{rs} . Then, the dynamic user equilibrium solution \hat{H} satisfies

$$\mathcal{N}(\hat{H}) \cdot (\hat{H} - H) \leq 0 \quad \forall H \in \bar{H}, \quad (12.13)$$

where the product \cdot is the Frobenius product, obtained by treating the matrices as vectors and calculating their dot product.

In the case of departure time choice, we can define $S = \mathcal{S}(H)$ to be the matrix of total costs for each path and departure time, where \mathcal{S} is obtained by composing the schedule delay function (11.4) with the network loading mapping \mathcal{N} . In this case, the set of feasible path flow matrices is given by

$$\hat{H} = \left\{ H \in \mathbb{R}_+^{\bar{T} \times |\Pi|} : \sum_{\pi \in \Pi_{rs}} \sum_t h_t^\pi = d^{rs} \quad \forall (r, s) \in Z^2 \right\}, \quad (12.14)$$

where d^{rs} is the OD matrix giving total flows between each origin and destination throughout the analysis period. The variational inequality in this case is

$$\mathcal{S}(\hat{H}) \cdot (\hat{H} - H) \leq 0 \quad \forall H \in \hat{H}. \quad (12.15)$$

Although these fixed point and variational inequality formulations encode the equilibrium principle, and concretely specify the problem at hand, little else can be proven concretely except in the case of simple link and node models (such as point queues). The vast majority of theoretical results on variational inequalities and fixed point problems relies on regularity conditions, such as continuity or monotonicity of the mappings \mathcal{N} or \mathcal{H} . For realistic link models and node models, these mappings are not continuous. For example, queue spillback imposes discontinuities on the network loading mapping. This means that we usually cannot prove that dynamic user equilibrium always exists or is unique, and indeed later in this chapter we present counterexamples where this need not be the case. However, the variational inequality can be used to define gap measures that can measure the degree of compliance with the equilibrium principle. In practice, whether the increased realism of a dynamic traffic model outweighs these disadvantages is case-dependent, and these considerations enter into how the appropriate modeling tools are chosen.

12.2 Solving for Dynamic Equilibrium

This section describes how dynamic user equilibrium solutions can be found. All of the algorithms in this section follow the same general framework, based on iterating between network loading, finding time-dependent shortest paths, and a path updating procedure, and they share common termination criteria. Three specific path updating procedures are discussed here: the convex combinations method, simplicial decomposition, and gradient projection. These are all heuristics, without making stronger assumptions on the properties of the network loading. This is demonstrated more fully in the next section, where we show that dynamic traffic assignment does not always share the neat solution properties of static assignment. These methods are described only for the case of fixed departure times, but it is not difficult to extend any of them to the case of departure time choice with schedule delay.

12.2.1 General framework

The general framework for almost all dynamic traffic assignment algorithms involves iterating three steps: network loading, as discussed in Chapter 10, time-dependent shortest paths, as discussed in Chapter 11, and updating the path flow matrix H , as discussed here. The algorithm takes the following form:

1. Initialize the path flow matrix H to a feasible value in \bar{H} .
2. Perform the network loading using path flows H , obtaining path travel times $T = \mathcal{N}(H)$.
3. Identify the time-dependent shortest paths.
4. Check termination criteria; if sufficiently close to dynamic user equilibrium, stop.

5. Update the path flow matrix H and return to step 2.

A few steps in this algorithm warrant further explanation. The choice of the initial solution is arbitrary. One reasonable choice is to identify one or more low travel-time paths between each OD pair based on free-flow travel times, and to divide the demand in the OD matrix between them. The effort and time involved in identifying the initial solution should be balanced against the benefit obtained by a more careful choice, as opposed to starting with a simpler choice and just running additional iterations.

The termination criteria in step 4 can take several forms. The most theoretically sound termination criterion is a gap measure, comparing the path flows H and travel times T with the assumed behavioral rule for travelers. Such measures are preferred to simpler criteria, such as stopping when the path flows do not change much from iteration to iteration — if the latter occurs, there is no way to tell whether this is because we are near an equilibrium solution, or because the algorithm is unable to make any more progress and has gotten stuck near a non-equilibrium solution. For the case of fixed demand, the behavioral rule is that travelers must take a least travel-time path between their origin and destination, for their departure time. After performing the time-dependent shortest path step, we can calculate the travel time on such a path for each OD pair (r, s) and departure time t ; call this value $\tau_{rs,t}^*$. If all travelers in the network were on such paths, then the total travel time experienced by all travelers in the network would be given by

$$SPTT = \sum_{(r,s) \in Z^2} \sum_t d_t^{rs} \tau_{rs,t}^*. \quad (12.16)$$

This is often known as the *shortest path travel time*. This is contrasted with the *total system travel time*, which is the actual total travel time spent by all travelers on the network:

$$TSTT = \sum_{\pi \in \Pi} \sum_t h_{\pi}^t \tau_{\pi}^t = H \cdot T. \quad (12.17)$$

Clearly $TSTT \geq SPTT$, and $TSTT = SPTT$ only if all travelers are on least-time paths, which corresponds to dynamic user equilibria. Therefore, the gap between these values is a quantitative measure of how close the path flows H are to satisfying the equilibrium principle. Two common gap values involve normalizing the difference $TSTT - SPTT$ in two ways. The *relative gap* γ normalizes this difference by SPTT,

$$\gamma = \frac{TSTT - SPTT}{SPTT}, \quad (12.18)$$

while the *average excess cost* normalizes this difference by the total number of travelers on the network,

$$AEC = \frac{TSTT - SPTT}{\sum_{(r,s) \in Z^2} \sum_t d_t^{rs}}. \quad (12.19)$$

The algorithm can be terminated whenever either of these measures is sufficiently small. An advantage of the average excess cost is that it is measured in time units, and has the intuitive interpretation as being the average difference between a traveler's actual travel time, and the shortest path travel time available to them. Similar expressions can be derived for departure time choice, with schedule delay.

Finally, the last step (updating the path flow matrix) requires the most care, and this section presents three alternatives. The convex combinations method is the simplest and most economical in terms of computer memory. The simplicial decomposition method converges faster, by storing previous solutions and using them to find a better improvement direction for H . The third method is gradient projection, which involves computation of an approximate derivative of path travel times with respect to path flow. This derivative enables the use of Newton's method to identify the amount of flow to shift between paths.

12.2.2 Convex combinations method

The convex combinations method is the simplest way to update the path choice matrix. In this method, after computing time-dependent shortest paths, one identifies a "target matrix" H^* which would give the path flows travelers would choose if the path travel times were held constant at their current values T . For the case of fixed departure times, for each OD pair (r, s) and departure time t , all of the demand d_t^{rs} is loaded onto its time-dependent shortest path in H^* , and all other path flows are set to zero. For the case of departure time choice, for each OD pair (r, s) the departure time t^* and corresponding path with minimal schedule delay is identified, and all d^{rs} travelers associated with this OD pair are assigned to that path and departure time, with all other path flows set to zero. This is often called an *all-or-nothing assignment*, because it involves choosing a single alternative to assign an entire class of travelers.

The basic form of the convex combinations method updates H using the formula

$$H \leftarrow \lambda H^* + (1 - \lambda)H, \quad (12.20)$$

where $\lambda \in [0, 1]$ is a step size showing how far to move in the direction of the target. If $\lambda = 0$, then the path flow matrix is unchanged, whereas if $\lambda = 1$, the path flow matrix is set equal to the target. Intermediate values produce a solution which is a weighted average of the current solution H and the target H^* . There are a number of ways to choose the step sizes λ . The simplest alternative is to use a decreasing sequence of values, choosing for the i -th iteration the step size $\lambda_i = 1/(i + 1)$, so the step sizes form the sequence $1/2, 1/3, 1/4$, and so on. It is common to use sequences such that $\sum \lambda_i = +\infty$ and $\sum \lambda_i^2 < \infty$, because we do not want the sequence of H values to converge too quickly (or else it may stop short of the equilibrium), but if the step size drops too slowly, the algorithm is likely to "overshoot" the equilibrium and oscillate.

More sophisticated variations are possible as well. Rather than using a fixed step size, λ can be chosen to minimize, say, the relative gap or average excess

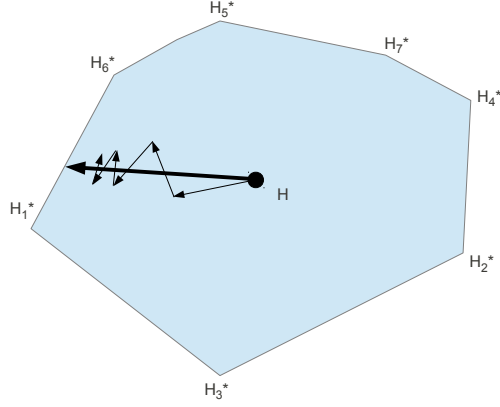


Figure 12.5: Simplicial decomposition can move towards non-corner points.

cost. While this likely decreases the number of iterations required to reach a small gap, the computation at each iteration is increased. In particular, each λ value tested requires an entire network loading step to determine the corresponding gap, which can be a significant time investment. It is also possible to vary the λ value for different path flows and departure times. Some popular heuristics are to use larger λ values for OD pairs which are further from equilibrium, or to vary λ for different departure times — since travel times for later departure times depend on the path choices for earlier departure times, it may not make sense to invest much effort in equilibrating later departure times until earlier ones have stabilized.

12.2.3 Simplicial decomposition

The simplicial decomposition algorithm extends the convex combinations method by remembering the target matrices H^* from earlier iterations. This requires more computer memory, but storing these previous target matrices leads to a more efficient update of the path flow vector. In particular, choosing an all-or-nothing assignment as a target matrix essentially restricts the choice of search direction to the corner points of the feasible region. As a result, when the equilibrium solution is reached, the convex combinations method will “zig-zag,” taking a number of short, oblique steps rather than a more direct step towards the equilibrium problem. (Figure 12.5). Simplicial decomposition is able to do so by combining multiple target points into the search direction. In this algorithm, the set \mathcal{H} is used to store all of the target path flow matrices found thus far.

A second notion in simplicial decomposition is that of a “restricted equilibrium.” Given a set $\mathcal{H} = \{H_1^*, H_2^*, \dots, H_k^*\}$, we say that a time-dependent path flow matrix \hat{H} is a *restricted equilibrium relative to \mathcal{H}* if it solves the variational

inequality

$$\mathcal{N}(\hat{H}) \cdot (\hat{H} - H_i^*) \leq 0 \quad \forall H_i^* \in \mathcal{H}. \quad (12.21)$$

This is different from the variational inequality (12.13) because the only possible choices for H are the matrices in the set \mathcal{H} , rather than *any* of the feasible H matrices satisfying (12.12). Effectively, H is a restricted equilibrium if none of the targets in \mathcal{H} lead to improving directions in the sense that the total system travel time would be reduced by moving to some $H_i^* \in \mathcal{H}$ while fixing the travel times at their current values.

At a high level, simplicial decomposition works by iterating between adding a new target matrix to \mathcal{H} , and then finding a restricted equilibrium using the current matrices in \mathcal{H} . In practice, it is too expensive to exactly find a restricted equilibrium at each iteration. Instead, several “inner iteration” steps are taken to move towards a restricted equilibrium with the current set \mathcal{H} before looking to add another target. In each inner iteration, the current solution H is adjusted to $H + \mu \Delta H$, where μ is a step size and ΔH is a direction which moves toward restricted equilibrium. One good choice for this direction is

$$\Delta H = \frac{\sum_{H_i^* \in \mathcal{H}} [\mathcal{N}(H) \cdot (H - H_i^*)]^+ (H_i^* - H)}{\sum_{H_i^* \in \mathcal{H}} [\mathcal{N}(H) \cdot (H - H_i^*)]^+} \quad (12.22)$$

This rather intimidating-looking formula is actually quite simple. It is nothing more than a weighted average of the directions $H_i^* - H$ (potential moves toward each target in \mathcal{H}), where the weight for each potential direction is the extent to which it improves upon the current solution: $[\mathcal{N}(H) \cdot (H - H_i^*)]$ is the reduction in total system travel time obtained by moving from H to H_i^* while holding travel times constant. If this term is negative, there is no need to move in that direction, so the weight is simply set to zero. The denominator is simply the sum of the weights, which serves as a normalizing factor.

The step size μ is chosen through trial-and-error. One potential strategy is to iteratively test μ values in some sequence (say, $1, 1/2, 1/4, \dots$) until we have found a solution acceptably closer to restricted equilibrium than H .¹ “Acceptably closer” can be calculated using the *restricted average excess cost*

$$AEC' = \frac{\mathcal{N}(H) \cdot H - \min_{H_i^* \in \mathcal{H}} \{\mathcal{N}(H) \cdot H_i^*\}}{\sum_{(r,s) \in Z^2} \sum_t d_t^{rs}} \quad (12.23)$$

which is similar to the average excess cost, but instead of using the shortest path travel time uses the best available target vector in \mathcal{H} .

Unlike the version of simplicial decomposition used for the static traffic assignment problem (Chapter 7), there is no guarantee that a sufficiently small choice of μ will result in a reduction of restricted average excess cost. However, in practice, this rule seems to work acceptably.

Putting all of this together, for simplicial decomposition, step 5 of the dynamic traffic assignment algorithm in Section 12.2.1 involves performing all of these steps as a “subproblem”:

¹Those familiar with nonlinear optimization may see parallels between this and the Armijo rule.

5. **Subproblem:** Find an approximate restricted equilibrium H using only the vectors in \mathcal{H} .
- (a) Find the improvement direction ΔH using equation (12.22).
 - (b) Update $H \leftarrow H + \mu \Delta H$, with μ sufficiently small (to reduce AEC').
 - (c) Perform network loading to update travel times.
 - (d) Return to step (a) of subproblem unless AEC' is small enough.

Furthermore, the set \mathcal{H} must be managed. It is initialized to be empty, and whenever time-dependent shortest paths are found, a new all-or-nothing assignment (the one which would have been chosen as the sole target in the convex combinations method) is added to \mathcal{H} .

12.2.4 Gradient projection

The gradient project method updates the path flow matrix H by using Newton's method. Consider first the simpler problem of trying to shift flows between two paths π_1 and π_2 (for the same departure time) to equalize their travel times. Write $\tau_1(h_1, h_2)$ and $\tau_2(h_1, h_2)$ to denote the travel times on the two paths as a function of the flows h_1 and h_2 on the two paths. If we were to shift Δh vehicles from path 1 to path 2, the difference in travel times between the paths is given by

$$g(\Delta h) = \tau_1(h_1 - \Delta h, h_2 + \Delta h) - \tau_2(h_1 - \Delta h, h_2 + \Delta h). \quad (12.24)$$

We want to choose Δh so that the difference $g(\Delta h)$ between the path travel times is equal to zero.

Using one step of Newton's method, an approximate value of Δh is given by

$$\Delta h = -\frac{g(0)}{g'(0)}. \quad (12.25)$$

The numerator $g(0)$ is simply the difference in travel times between the two paths before any flow is shifted. To calculate the denominator, we need to know how the difference in travel times will change as flow is shifted from π_1 to π_2 . Computing the derivative of (12.24) with the help of the chain rule, we have

$$g'(0) = \left(\frac{\partial t_1}{\partial h_2} + \frac{\partial t_2}{\partial h_1} \right) - \left(\frac{\partial t_1}{\partial h_1} + \frac{\partial t_2}{\partial h_2} \right). \quad (12.26)$$

This formula is not easy to evaluate exactly, but we can make a few approximations. The first term in (12.26) reflects the impact the flow on one path has on the *other* path's travel time, while the second term reflects the impact the flow on one path has on *its own* travel time. Typically, we would expect the second effect will be larger in magnitude than the first effect (although exceptions do exist). So, as a first step we can approximate the derivative as

$$g'(0) \approx -\left(\frac{\partial t_1}{\partial h_1} + \frac{\partial t_2}{\partial h_2} \right). \quad (12.27)$$

The next task is to calculate the derivative of a path travel time with respect to flow along this path. Unlike in static traffic assignment, there is no closed-form expression mapping flows to travel times, but rather a network loading procedure must be used. For networks involving triangular fundamental diagrams (which can include the point queue model, cf. Section 10.5.3), the derivative of the travel time on a single link (i, j) with respect to flow entering at time t can be divided into two cases. In the first case, suppose that the vehicle entering link (i, j) at time t exits at $t' = t + \tau_{ij}(t)$, and that the link is demand-constrained at that time (that is, all of the sending flow $S_{ij}(t')$ is able to move). In this case, even if a marginal unit of flow is added at this time, the link will remain demand-constrained, and all of the flow will still be able to move. No additional delay would accrue, so the derivative of the link travel time is

$$\frac{d\tau_{ij}(t)}{dh} = 0. \quad (12.28)$$

In the second case, suppose that at time t' flows leaving link (i, j) are constrained either by the capacity of the link or by the receiving flow of a downstream link. In this case, not all of the flow is able to depart the link, and a queue has formed at the downstream end of (i, j) . The clearance rate for this queue is given by $y_{ij}(t')$, so the incremental delay added by one more vehicle joining the queue is $1/y_{ij}(t')$, and

$$\frac{\tau_{ij}(t)}{dh} = \frac{1}{y_{ij}(t')}. \quad (12.29)$$

(The denominator of this expression cannot be zero, since t' is the time at which a vehicle is leaving the link.)

Therefore, we can calculate the derivative of an entire path's travel time inductively. Assume that $\pi = [r, i_1, i_2, \dots, s]$, and that t_i gives the travel time for arriving at each node i in the path. Then dt_π/dh_π is obtained by summing expressions (12.28) and (12.29) for the uncongested and congested links in this path, taking care to use the correct time indices:

$$\frac{dt_\pi}{dh_\pi} \approx \sum_{(i,j) \in \pi} \frac{1}{y_{ij}(t_j)} [y_{ij}(t_j) < S_{ij}(t_j)], \quad (12.30)$$

where the brackets are an indicator function equal to one if the statement in brackets is true, and zero if false.

We are almost ready to give the formula for (12.27), but we can make one more improvement to the approximation. Two paths may share a certain number of links in common. Define the *divergence node* of two paths to be node where the paths diverge for the first time. (Figure 12.6). Prior to the divergence node, the two paths include the same links, so there will be no effect of shifting flow from one path to the other. Therefore, the sum in (12.30) need only be taken for links beyond the divergence node. Even if the paths rejoin at a later point, they may do so at different times, so we cannot say that there is no effect

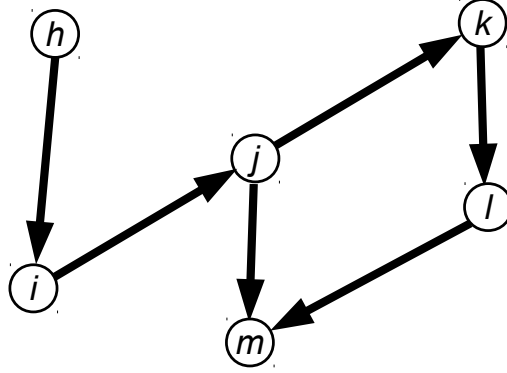


Figure 12.6: Node j is the divergence node of these two paths.

of shifting flow between common links downstream of a divergence node. So, if $d(\pi_1, \pi_2)$ is the divergence node, then we have

$$\frac{dt_\pi}{dh_\pi} \approx \sum_{(i,j) \in \pi: (i,j) > d(\pi_1, \pi_2)} \frac{1}{y_{ij}(t_j)} [y_{ij}(t_j) < S_{ij}(t_j)], \quad (12.31)$$

where the “ $>$ ” notation for links indicates links downstream of a node in a path. This expression can then be used in (12.27) and (12.25) to find the approximate amount of flow which needs to be shifted from π_1 to π_2 to equalize their costs.

The procedure for updating H can now be described as follows:

5. **Subproblem:** Perform Newton updates between all paths and the target path found.
 - (a) For each path π with positive flow, let ρ be the least-travel time path connecting the same OD pair and departure time.
 - (b) Calculate Δh using equations (12.31), (12.27), and (12.25), using π as the first path and ρ as the second.
 - (c) If $\Delta h < h_\pi$, then update $h_\rho \leftarrow h_\rho + \Delta h$ and $h_\pi \leftarrow h_\pi - \Delta h$. Otherwise, set $h_\rho \leftarrow h_\rho + h_\pi$ and $h_\pi \leftarrow 0$.

As stated, the path travel times and derivatives are not updated after each shift is performed. Doing so would increase accuracy, but greatly increase computation time since a complete network loading would have to be performed.

12.3 Properties of Dynamic Equilibria

The dynamic traffic assignment problem is much less well-behaved than the static traffic assignment problem, where a unique equilibrium provably exists

under natural conditions. This is because the more realistic network loading procedures used in dynamic traffic assignment lack the regularity properties which make static assignment more convenient mathematically. Recall that the dynamic user equilibrium path flows solve the variational inequality

$$\mathcal{N}(\hat{H})(\hat{H} - H) \leq 0 \quad \forall H \in \bar{H}. \quad (12.32)$$

Showing existence of a solution to this variational inequality would typically rely on showing that \mathcal{N} is continuous. However, when queue spillback is modeled, the travel times are discontinuous in the path flow variables. Showing uniqueness of a solution to the variational inequality typically involves showing that \mathcal{N} has some flavor of monotonicity, but common node models can be shown to violate monotonicity even when we restrict attention to networks involving only a single diverge and merge. Thus, the failure of existence or uniqueness results is directly traced to the features that make dynamic traffic assignment a more realistic model. As an analyst, you should be aware of this tradeoff.

This section catalogs a few examples of dynamic traffic assignment problems highlighting some unusual or counterintuitive results. We begin with an example where the choice of node model results in no dynamic user equilibrium solution existing. Two examples of equilibrium nonuniqueness are given, first where multiple user equilibrium solutions exist, and second where literally *all* feasible path flow matrices are user equilibrium solutions, even though they vary widely in total system travel time. We conclude with a dynamic equivalent to the Braess paradox which was developed by Daganzo, in which increasing the capacity on a link can make network-wide conditions arbitrarily worse, because of queue spillback.

12.3.1 Existence: competition at merges

The network in Figure 12.7 has two origin-destination pairs (A to B, and C to D). Aside from these four centroids, there are two nodes representing intersections. Turns are not allowed at these intersections, so $\Xi(1) = \{[A, 1, 2], [C, 1, D]\}$ and $\Xi(2) = \{[1, 2, B], [C, 2, D]\}$. As a result, each OD pair has only two routes available to it: $[A, B]$ and $[A, 1, 2, B]$ for A to B, and $[C, 1, D]$ and $[C, 2, D]$ for C to D. For the sake of convenience, call these four routes “top,” “bottom,” “left,” and “right,” respectively. The figure shows the *travel times* with each link. Two links in the network also carry a toll, expressed in the same units as travel time, and drivers choose routes to minimize the sum of travel time and toll. Capacities and jam densities are large enough that congestion will never arise, and all links will remain at free flow.

Nodes 1 and 2 have a distinctive node model, representing absolute priority of one approach over the other. For Node 1, all flow on the movement $[A, 1, 2]$ must yield to flow on $[C, 1, D]$. This can be expressed with the following node

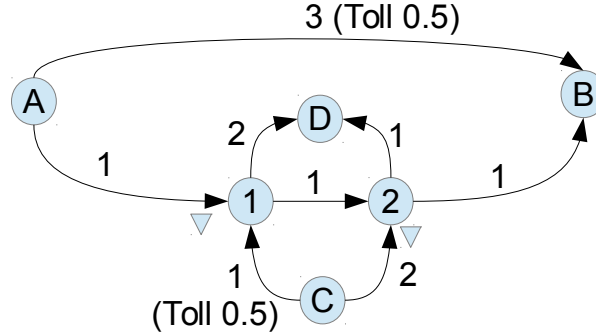


Figure 12.7: Network for example with no equilibrium solution. Link travel times shown.

model:

$$y_{C1D} = \min \{S_{C1}, R_{1D}\} \quad (12.33)$$

$$y_{A12} = \begin{cases} 0 & \text{if } y_{C1D} > 0 \\ \min \{S_{A1}, R_{12}\} & \text{otherwise} \end{cases} \quad (12.34)$$

Likewise, for Node 2, all flow on the movement $[C, 2, D]$ must yield to flow on $[1, 2, B]$. That is, traffic from C to D has priority at node 1; and traffic from A to B has priority at node 2. These relationships are indicated on the figure with triangles next to the approach that must yield.

Each OD pair has one unit of demand, that departs during the same time interval. We now show that there is no assignment of demand to paths that satisfies the equilibrium principle, requiring any used path to have minimal travel time. We begin first by showing that any equilibrium solution cannot split demand among multiple paths; each OD pair must place all of its flow on just one path or the other. With the given demand, the cost on the bottom path is either 3 (if there is no flow on the left path from C to D) or 4 (if there is), in either case, the cost is different from that on the top path (3.5). All demand will thus use either the top path or the bottom path, whichever is lower. Similarly, for the OD pair from C to D, the cost on the right path is either 3 or 4, whereas that on the left is either 3.5 or 4.5. Both paths cannot have equal travel time, so at most one path can be used.

Thus, there are only four solutions that can possibly be equilibria: all flow from A to B must either be on the top path or the bottom path; and all flow from C to D must either be on the left path or the right path. First consider the case when all flow from A to B is on the top path. The left path would then have a cost of 3.5, and the right path a cost of 3. So all flow from A to B must be on the right path. But then the travel times on the top and bottom paths are

3.5 and 3, respectively, so this solution cannot be an equilibrium (the travelers from A to B would prefer the bottom path, contradicting our assumption).

So now consider the other possibility, when all flow from A to B is on the bottom path. The travel times on the left and right paths are now 3.5 and 4, so all travelers from C to D would pick the left path. But then the travel times on the top and bottom paths are 3.5 and 4, respectively, so this solution is not an equilibrium either.

Therefore, there is *no* assignment of vehicles to paths that satisfies the principle of user equilibrium. This network is essentially the Ginger/Harold “matching pennies” game from Section 1.3 encoded into dynamic traffic assignment. Game theorists often resolve this type of situation by proposing a *mixed-strategy equilibrium*, in which the players randomize their actions. One can show that a mixed-strategy equilibrium exists in this network (Exercise 8), but this equilibrium concept is usually not applied to transportation planning. This is for both computational reasons (when there are many players, as in practical planning networks, computing mixed-strategy equilibria is hard) and for modeling reasons (typically travelers do not randomize their paths with the intent of “outsmarting” other travelers).

Mathematically, the reason no equilibrium exists is because the node model specified by (12.33) and (12.34) is not continuous in the sending and receiving flows. This is similar to static traffic assignment, where equilibrium existence cannot be guaranteed if the link performance functions are not continuous. In dynamic traffic assignment, when building node models that capture absolute priority (as at yields, two-way stops, or signals with permitted turns), ensuring continuity is difficult. In part, this is the motivation for using ratios of α values in Chapter 10 to reflect priorities at merges and general intersections.

But even if the link and node models are chosen to be continuous, dynamic equilibrium need not exist. Repeating the arguments made in Chapter 6.2 would also require the travel time calculations to be continuous in the link cumulative entrances and exits N^\uparrow and N^\downarrow . The use of a minimum to disambiguate multiple possible values of travel times in Equation (12.3) can mean that small changes in a link’s cumulative entries or exits could cause a large change in the travel time (imagine how the length of the dashed line in Figure 12.3 would change if one of the corner points were shifted slightly). Continuity could be assured only if the cumulative counts were *strictly* increasing, that is, if vehicles were constantly flowing into and out of each link. Therefore, although one can construct dynamic traffic assignment models which guarantee existence of equilibria, these require quite strong and significant restrictions on the network loading and travel time calculation procedures.

12.3.2 Uniqueness: Nie’s merge

The network in Figure 12.8 consists of a single diverge and merge. One can visualize this network as a stylized version of a freeway lane drop (Figure 12.9), where the reduction from two lanes to one lane reduces the capacity from 80 vehicles per minute to 40 vehicles per minute, and where drivers either merge

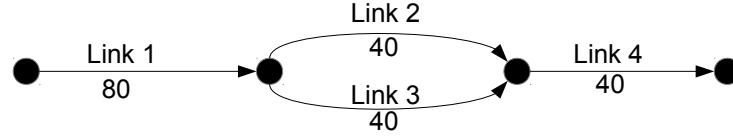


Figure 12.8: Nie's merge network. All capacity values in vehicles per minute.

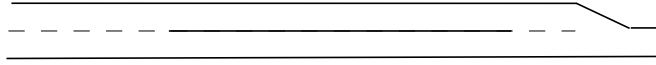


Figure 12.9: A physical interpretation of Nie's merge.

early (choosing the top lane at the diverge point) or late (waiting until the lane drop itself). The inflow rate from upstream is 80 vehicles per minute. Since the capacity of the downstream exit is only 40 vehicles per minute, the excess vehicles will form one or more queues in this network.

To simplify the analysis, assume that the time horizon is short enough that none of these queues will grow to encompass the entire length of the link. With this assumption, queue spillback can be ignored, and we can focus on the issue of route choice. Furthermore, under these assumptions, S_1 is always 80 vehicles per minute, and R_2 , R_3 , and R_4 are always 40 vehicles per minute. We work with a timestep of one minute, and will express all flow quantities in vehicles per minute.²

In this network, there is only one choice of routes (the top or bottom link at the diverge). We will restrict our attention to path flow matrices H where the proportions of vehicles choosing the top and bottom routes are the same for all departure times, and show that there are multiple equilibrium solutions even when limiting our attention to such H matrices. There may be still more equilibria where the proportion of vehicles choosing each path varies over time.

Therefore, the ratio h_{124}/h_{134} is constant for all time intervals, which means that the splitting fractions p_{12} and p_{13} at the diverge point are also constant, and equal to $h_{124}/40$ and $h_{134}/40$, respectively. Using the diverge and merge models from Section 10.2.3, we can analyze the queue lengths and travel times on each link as a function of these splitting fractions.

It turns out that three distinct dynamic user equilibria exist:

²This is larger than what is typically used in practice, but simplifies the calculations and does not affect the conclusions of this example.

Equilibrium I: If $p_{12} = 1$ and $p_{13} = 0$, then the diverge formula (10.17) gives the proportion of flow which can move as

$$\phi = \min \left\{ 1, \frac{40}{80}, \frac{40}{0} \right\} = \frac{1}{2}. \quad (12.35)$$

Therefore, the transition flows at each time step are $y_{12} = 40$ and $y_{13} = 0$. Since $y_{12} < S_{12}$, a queue will form on the upstream link, and its sending flow will remain at $S_1 = 80$. Therefore, once the first vehicles reach the merge, we will have $S_2 = 40$ and $S_3 = 0$. Applying these proportions, together with $R_4 = 40$, the merge formula gives $y_{24} = 40$ and $y_{34} = 0$. There will be no queue on link (2,4), so both link (2,4) and link (3,4) are at free-flow. This solution is an equilibrium: the two paths through the network only differ by the choice of link 2 or link 3, and both of these have the same travel time since they are both at free-flow. In physical terms, this corresponds to all drivers choosing to merge early; the queue forms upstream of the point where everyone chooses to merge, and there is no congestion downstream.

Equilibrium II: If $p_{12} = 0$ and $p_{13} = 1$, the solution is exactly symmetric to that of Equilibrium I. A queue will again form on the upstream link, now because all drivers are waiting to take the bottom link. In physical terms, this corresponds to all drivers choosing to merge early, but to the lane which is about to drop. They then all merge back when the lane actually ends. This solution does not seem especially plausible, but it does satisfy the equilibrium condition.

Equilibrium III: If $p_{12} = p_{13} = 1/2$, drivers wish to split equally between the top and bottom links. The proportion of flow which can move at the diverge is

$$\phi = \min \left\{ 1, \frac{40}{40}, \frac{40}{40} \right\} = 1, \quad (12.36)$$

so all vehicles can move and there is no queue at the diverge: $y_{12} = y_{13} = 40$. Once these vehicles reach the merge, we will have $S_2 = S_3 = 40$ and $R_4 = 40$. The merge formula (10.12) then gives $y_{24} = y_{34} = 20$, so queues will form on both merge links. However, since the inflow and outflow rates of links 2 and 3 are identical, the queues will have identical lengths, and so the travel times on these links will again be identical. Therefore, this solution satisfies the principle of dynamic user equilibrium. In physical terms, this is the case when no drivers change lanes until the lane actually ends, and queueing occurs at the merge point.

By examining the intermediate solutions, we can calculate the travel times on both paths for different values of the fraction p_{12} . (We do not need to state the values of the other fraction, since $p_{13} = 1 - p_{12}$.) This is shown in Figure 12.10, and the three equilibria correspond to the crossing points of the paths. Note that all three of the equilibria share the same equilibrium travel time. Although

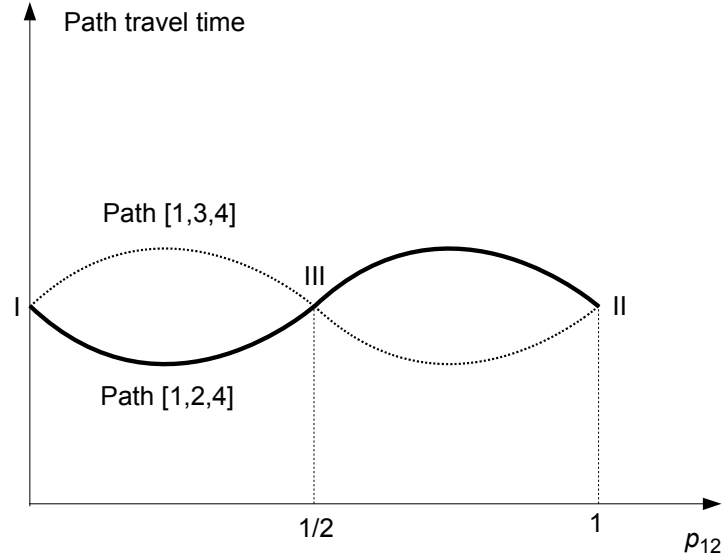


Figure 12.10: Travel time on top and bottom paths as a function of splitting fraction, with three equilibrium solutions marked.

the queues in Equilibrium III are only half as long as those in Equilibria I and II, being split between two links, the outflow rates of these queues are also only half as great (20 veh/min instead of 40 veh/min).

This example shows that the dynamic user equilibrium solution is not unique, even in an extremely simple network. This nonuniqueness has practical consequences. The effect of a potential improvement to link 2 or 3 will depend crucially on the number of travelers on the link, which varies widely in all three equilibria. One criterion for distinguishing which of these equilibria is more likely is *stability*, which explores what would happen if the equilibria were slightly perturbed. As shown in Figure 12.10, if one begins at Equilibrium I and perturbs p_{13} to a small value (reducing p_{12} by the same amount), we see that the travel time on the top path increases, whereas the travel time on the bottom path decreases. It may seem odd that increasing flow on the bottom path decreases its travel time — what is happening is that congestion at the diverge decreases (lowering the travel time of both paths), but congestion forms on the top link at the merge (increasing the travel time just of the top path). Superimposing these effects produces the result in Figure 12.10. As a result, travelers will switch from the (slower) top path to the (faster) bottom one, moving us even further away from Equilibrium I. Therefore, Equilibrium I is not stable. The same analysis holds for Equilibrium II.

Equilibrium III, on the other hand, is stable. If a few travelers switch from the top to the bottom path, the travel time on the top path decreases and that on the bottom path increases. Therefore, travelers will tend to switch back to the

top path, restoring the equilibrium. The same holds if travelers switch from the bottom path to the top path. This gives us reason to believe that Equilibrium III is more likely to occur in practice than Equilibrium I or II. This type of analysis is much more complicated in larger networks, and for the most part is completely unexplored. Coming to a better understanding of the implications of nonuniqueness in large networks, as well as techniques for addressing this, is an important research question.

12.3.3 Nie's merge redux

Modifying the merge network from the previous section, we can produce an even more pathological result. We now eliminate the lane drop entirely, but preserve the diverge/merge network, by setting the capacity of the downstream link to 80 vehicles per minute, the same as the upstream link. This might represent a roadway section where lane changing is prohibited, perhaps in a work zone or with a high-occupancy/toll lane. Now, at the merge point $S_2 + S_3 \leq R_4$ no matter what the flow pattern is on the network. This means that the merge will always be freely-flowing, and there will be no queues at the merge. Links 2, 3, and 4 will be at free-flow regardless of the path flows H .

However, queues can still occur at the *diverge*, where the analysis is the same as before. If $p_{12} = p_{13} = 1/2$, no queues will form at the diverge point, whereas for any other values of p_{12} and p_{13} formula (12.35) will predict queues on the upstream link, reaching their maximum length if $p_{12} = 1$ or $p_{13} = 1$. However, *all of these solutions satisfy the principle of equilibrium* because the only delay occurs on the upstream link, which is common to both paths. No matter what the values of p are, there is zero relative gap or average excess cost, since no traveler can reduce his or her travel time by choosing a different route. Choosing a different route could influence the travel time of those further upstream, but under the assumption that drivers act only to minimize their own travel time, this influence on other travelers is of no concern.

Figure 12.11 shows the queue lengths and total system travel time as p varies. The system-optimal solution is unique: if $p_{12} = p_{13} = 1/2$ there are no queues in the system and all vehicles experience free-flow travel times. For any other values of p_{12} and p_{13} , a queue will form at the diverge and some delay is experienced. Yet *all* possible p values satisfy the equilibrium principle, since the delay is upstream of the diverge and drivers cannot choose another route to minimize their own travel time. Furthermore, as the time horizon grows longer, the difference in total travel time between the worst of the equilibria (either $p_{12} = 0$ or $p_{12} = 1$) and the system optimum solution can grow arbitrarily large. This is in contrast with the bounded “price of anarchy” which can often be found for static traffic assignment (cf. Section 6.4.3).

This effect less of a practical concern than the original version of Nie's merge — there are good behavioral reasons to doubt that a significant imbalance of travelers will choose one alternative over another identical one. One argument is from entropy principles (Section 6.2.2): if travelers have the same behavior assumption, it is unlikely they would all choose one route over another with equal

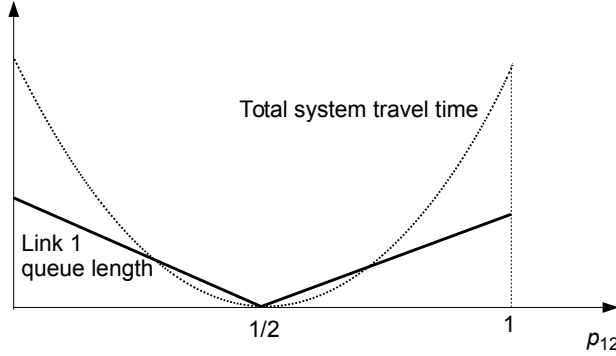


Figure 12.11: Queue length and total travel time for different splitting proportions p in the modified Nie’s merge network. Note that all p values represent equilibria.

travel time. Furthermore, the assumption of a triangular fundamental diagram implies that travel speeds remain at free-flow for all subcritical densities. In practice the speed will drop slightly due to variations in preferred speeds and difficulties in overtaking at higher density, so drivers would likely prefer the route chosen by fewer travelers.

However, from the standpoint of *modeling* the fact that all feasible solutions are equilibria poses significant challenges. Literally any solution will have zero gap, and if an all-or-nothing assignment is chosen as the initial solution (as is sometimes done in implementations), dynamic traffic assignment software will report that a perfect equilibrium has been reached. This example shows that initial solutions for dynamic traffic assignment should be carefully chosen, perhaps by spreading vehicles over multiple paths, or breaking ties stochastically in shortest path algorithms to avoid assigning all vehicles to the same path in the first all-or-nothing assignment.

12.3.4 Efficiency: Daganzo’s paradox

The previous section showed that the worst user equilibrium solution can be arbitrarily worse than the system optimal solution, in terms of total travel time. However, there was still one user equilibrium solution that was also system optimal (the case where $p_{12} = p_{13} = 1/2$). This need not be in the case. Here we present an example where the *only* user equilibrium solution can be arbitrarily worse than system optimum. Furthermore, it is “paradoxical” in the sense that increasing the capacity of the only congested link on the network can make the problem worse, and that reducing the capacity on this link can improve system conditions! In this sense, it is a dynamic equivalent of the Braess paradox from Section 5.3. Many have criticized the Braess paradox on the grounds that the link performance functions used in static assignment are unrealistic. In the example shown below, queue spillback (a feature unique to dynamic network

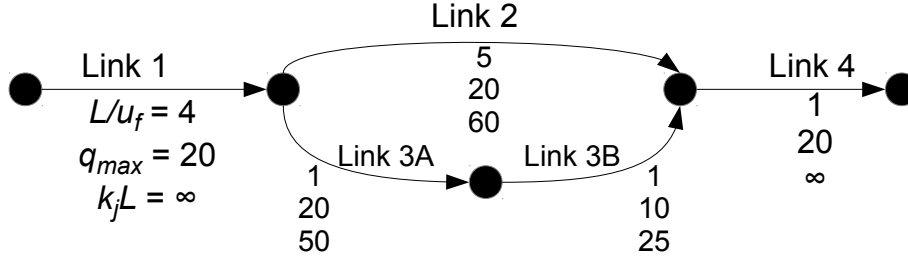


Figure 12.12: Network for Daganzo's paradox.

loading) is actually the critical factor in the paradox.

See the network in Figure 12.12, where time is measured in units of time steps Δt . Like the networks in the two previous sections, it consists of a single merge and diverge. However, the free-flow travel times and capacities on the top and bottom links are now different: the top route is longer, but has a higher capacity, while the bottom route is shorter at free-flow, but has a bottleneck limiting the throughput on this route — link 3B has only half the capacity of 3A. We will use the spatial queue model of Section 10.1.3 to propagate traffic flow, although the same results would be obtained with an LWR-based model or anything else which captures queue spillback and delay. The input demand is constant, at 20 vehicles per time step.

The capacity on the top route is high enough to accommodate all of the demand; if all of this demand were to be assigned to this route, the travel time would be 10 minutes per vehicle. Assigning all vehicles to the top route is neither the user equilibrium nor the system optimum solution, but it does give an upper bound on the average vehicle travel time in the system optimal assignment — it is possible to do better than this if we assign some vehicles to the bottom, shorter route.

To derive the user equilibrium solution, notice that initially all vehicles will choose the bottom path, since it has the lower travel time at free flow. A queue will start forming on link 3A, since the output capacity is only 10 vehicles per time step (because of the series node model from Section 10.2.1, and the capacity of link 3B) and vehicles are entering at double this rate. As the queue grows, the travel time on link 3A will increase as well. With the spatial queue model and these inflow and outflow rates, you can show that the travel time for the n -th vehicle entering the link is

$$1 + \frac{n}{q_{in}} \left(\frac{q_{in}}{q_{out}} - 1 \right) = 1 + \frac{n}{20} \quad (12.37)$$

as long as the queue has not spilled back (see Exercise 10.3).

Based on equation (12.37), when the 60th vehicle enters the link, it will spend four time units on link 3A, and thus its travel time across the entire bottom path would be the same as if it had chosen the top path. At this point there are 40 vehicles on link 3A, as can be seen by drawing the N^\uparrow and N^\downarrow

curves for this link. From this point on, vehicles will split between the top and bottom paths to maintain equal travel times.

So far, so good; the first 60 vehicles that enter the network have a travel time of less than 10 minutes, and all the rest have a travel time of exactly 10 minutes. Now see what happens if we increase the capacity on link 3B, with the intent of alleviating congestion by improving the bottleneck capacity. If the capacity on 3B increases from 10 to 12, the story stays the same, except that it is the first 90 vehicles that have a travel time of less than 10 minutes. We can see this by setting (12.37) equal to four (the time needed to equalize travel times on the top and bottom paths), but with $q_{in} = 12$ instead of 10, and solving for n . Network conditions indeed have improved. But if the capacity increases still further, to 15, then equation (12.37) tells us that it is only after 180 vehicles have entered the system that travelers would start splitting between the top and bottom links. By tracing the N^\uparrow and N^\downarrow curves, at this point in time 120 vehicles will have exited link 3A, meaning the queue length would be 60 vehicles. **But the jam density of the bottom link only allows it to hold 50 vehicles.** The queue will thus spill upstream of the diverge node, and in this scenario, *no* vehicles will opt to take the top path. By the time a driver reaches the diverge point, the number of vehicles on the bottom link is 50, giving a travel time on 3A of $3\frac{1}{3}$ minutes, and a total travel time of $9\frac{1}{3}$ minutes from origin to destination. This is less than the travel time on the top path, so drivers prefer waiting in the queue to taking the bypass route.

As a result, *all* drivers will choose the bottom path, and the queue on the origin centroid connector will grow without bound, as will the travel times experienced by vehicles entering the network later and later. By increasing the length of time vehicles enter the network at this rate, we can make the delays as large as we like.

We thus see that even with dynamic network loading, increasing the capacity on the only bottleneck link can make average travel times worse — and in fact arbitrarily worse than the system optimum solution, where the average travel time is less than 10 minutes. The reason for this phenomenon is the interaction between queue spillback and selfish route choice: in the latter scenario it would be better for the system for some drivers to choose the top route, even though it would worsen their individual travel time.

Of course, if the capacity on 3B were increased even further, all the way to 20 vehicles per time step, delays would drop again since there would be no bottleneck at all. Exercise 11 asks you to plot the average vehicle delay in this network as the capacity on link 3B varies from 0 to 25 vehicles per time step.

12.3.5 Conclusion

The purpose of these examples is to show that dynamic user equilibrium is complex. Guaranteeing existence or uniqueness of dynamic equilibrium requires making strong assumptions on traffic flow propagation. However, for some practical applications, using a more realistic traffic flow model is more important than mathematical properties of the resulting model. Many people find comfort

in the fact that we can at best solve for an equilibrium approximately, and thus dismiss the question of whether an equilibrium “truly” exists as akin to asking how many angels can dance on the head of a pin.

We emphasize that existence and uniqueness are not simply mathematical abstractions, and that they have significant implications for practice: if an equilibrium does not exist, should we really be ranking projects based on equilibrium analysis? If multiple equilibria exist, what should we plan for? Can we even find them all? At the same time, we acknowledge that using static equilibrium to sidestep these difficulties is often unacceptable. For many applications, the assumptions in link performance functions are simply too unrealistic. Such is the nature of mathematical modeling in engineering practice. No tool is right for every task, but rather experienced practitioners know how to match the available tools to the job at hand. By reading this book, we hope that you have gained the insight to understand the advantages and disadvantages of both static and dynamic traffic assignment models, and to make educated decisions about the right tool for a particular project.

Finally, more research is needed on these topics, understanding the significance of equilibrium nonexistence and nonuniqueness, particularly in practical networks and not just in “toy” networks such as the ones in this section. As researchers, we would be delighted for you to contribute to work in this field.

12.4 Exercises

1. [37] (*Equivalence of link-based and path-based flow representations*). Given splitting proportions $\alpha_{hij,s}^t$ for each destination s , time interval t , and turning movement $[h, i, j]$, show how “equivalent” path flow values h_t^π can be found. Then, if given path flows h_t^π , show how “equivalent” $\alpha_{hij,s}^t$ values can be found. (“Equivalent” means that the link cumulative counts N^\uparrow and N^\downarrow would be the same for all time steps after performing network loading, possibly with a small error due to time discretization that would shrink to zero as $\Delta t \rightarrow 0$.)
2. [37] Consider the four-link network shown in Figure 12.13, and perform network loading using the cell transmission model. (Each link is one cell long.) During the first time interval, 10 vehicles enter Link 1 on the top path; during the second time interval, 5 vehicles enter Link 1 on the top path and 5 on the bottom path; and during the third time interval, 10 vehicles enter Link 1 on the bottom path. No other vehicles enter the network.

Interpolating as necessary, what time do the first and last vehicles on the top path exit cell 4? the first and last vehicles on the bottom path? What is the derivative of the travel time on the top path for a vehicle leaving at the start of interval 2?
3. [12] In the course of the convex combinations algorithm, assume that the H and H^* matrices are as below, and $\lambda = 1/3$. What is the new H

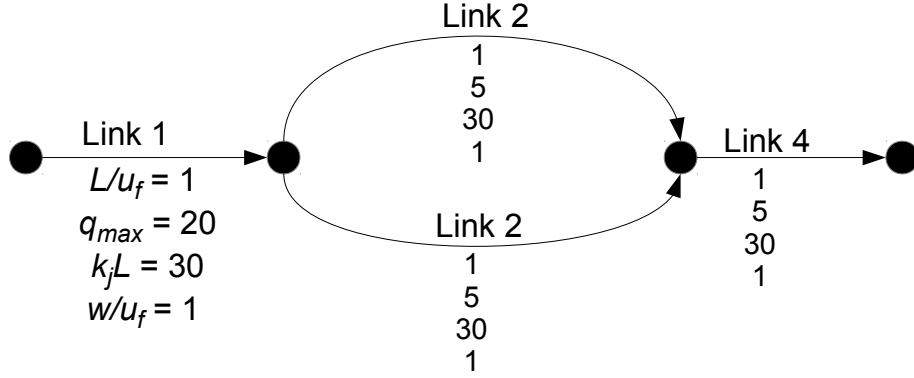


Figure 12.13: Network for Exercise 2.

matrix?

$$H = \begin{bmatrix} 6 & 12 \\ 30 & 24 \end{bmatrix} \quad H^* = \begin{bmatrix} 18 & 0 \\ 0 & 54 \end{bmatrix} \quad (12.38)$$

4. [23] Consider a network with only one origin-destination pair connected by four paths, with three departure time intervals. Departure times are fixed. At some point in the simplicial decomposition algorithm, \mathcal{H} contains the following three matrices:

$$\begin{bmatrix} 20 & 0 & 0 & 0 \\ 0 & 10 & 0 & 0 \\ 0 & 0 & 30 & 0 \end{bmatrix} \begin{bmatrix} 0 & 20 & 0 & 0 \\ 0 & 10 & 0 & 0 \\ 0 & 30 & 0 & 0 \end{bmatrix} \begin{bmatrix} 20 & 0 & 0 & 0 \\ 0 & 0 & 0 & 10 \\ 0 & 0 & 30 & 0 \end{bmatrix} \quad (12.39)$$

and the current path flow and travel time matrices are:

$$H = \begin{bmatrix} 14 & 6 & 0 & 0 \\ 0 & 8 & 0 & 2 \\ 0 & 9 & 21 & 0 \end{bmatrix} \quad T(H) = \begin{bmatrix} 20 & 20 & 24 & 27 \\ 30 & 34 & 37 & 40 \\ 44 & 35 & 36 & 40 \end{bmatrix} \quad (12.40)$$

What are the unrestricted and restricted average excess costs of the current solution? What is the matrix ΔH ?

5. [35] Assume that there is a single OD pair, two paths, and three departure times; 15 vehicles depart during the first, 10 during the second, and 5 during the third. Let $H_{\pi t}$ denote the number of vehicles departing on path π at time t , and that the path travel times are related to the path

flows by the following equations:³

$$\begin{aligned}T_{11} &= 15 \\T_{21} &= 13 + 3H_{11} + H_{21} \\T_{12} &= 15 + H_{11}^2 \\T_{22} &= 13 + H_{11} + 3H_{12} + H_{21}^2 + H_{22} \\T_{13} &= 15 + H_{12}^2 + \frac{1}{2}H_{11}^2 \\T_{23} &= 13 + H_{12} + 3H_{13} + H_{22}^2 + H_{23}\end{aligned}$$

Find the path flows obtained from three iterations of the convex combinations method with step sizes $\lambda_1 = 1/2$, $\lambda_2 = 1/4$, and $\lambda_3 = 1/6$ (so you should find a total of four H^* matrices, counting the initial matrix, and take three weighted averages). Your initial matrix should be the shortest paths with zero flow. What is the resulting AEC? *Break any ties in favor of path 1.*

6. [35] Using the same network and demand as in Exercise 5, now assume that you are solving the same problem with simplicial decomposition, and at some stage \mathcal{H} contains the following two matrices:

$$\begin{bmatrix} 15 & 0 \\ 0 & 10 \\ 5 & 0 \end{bmatrix} \quad \begin{bmatrix} 15 & 0 \\ 10 & 0 \\ 0 & 5 \end{bmatrix}$$

and that the current solution is

$$H = \begin{bmatrix} 15 & 0 \\ 4 & 6 \\ 3 & 2 \end{bmatrix}$$

- What is the average excess cost of the current solution?
 - What is the restricted average excess cost?
 - What is the search direction from H ?
 - Give the updated H matrix and new restricted average excess cost after taking a step with $\mu = 0.1$.
 - If we terminate the subproblem and return to the master algorithm, what matrix (if any) do we add to \mathcal{H} ?
7. [35] Again using the network and demand from Exercise 5, now start with the initial solution

$$H = \begin{bmatrix} 10 & 5 \\ 5 & 5 \\ 0 & 5 \end{bmatrix}$$

³In practice these would come from performing network loading and calculating path travel times as described in this chapter; these functions are provided here for your convenience.

Using the gradient projection method, identify a new H and T matrix. Try both exact and quasi-Newton steps, and see which gives the greater reduction in AEC .

8. [46] Find the mixed-strategy equilibrium in the network in Figure 12.7. (For each of the two vehicles, indicate the probability it will choose each of the two paths available to it; these probabilities should be such that the *expected* travel times are the same for both options.)
9. [35] Consider a network with one origin, one destination, two possible departure time intervals (1 and 2), and two paths (A and B). A total of 16 vehicles depart during time interval 1, and 4 depart during time interval 2. Let t_τ^π refer to the travel time on the π -th path when departing at the τ -th time interval, with h_τ^π defined similarly for path flows. Suppose the travel times are related to the path flows as follows:

$$\begin{aligned} t_1^A &= 0.5h_1^A + 5h_2^A + 6 \\ t_1^B &= 0.5h_1^B + 3h_2^B + 10 \\ t_2^A &= 0.3h_1^A + 0.6h_2^A + 0.8 \\ t_2^B &= 0.2h_1^B + 0.4h_2^B + 2 \end{aligned}$$

and consider the following four path flow matrices:

$$H_1 = \begin{bmatrix} 0 & 16 \\ 4 & 0 \end{bmatrix} \quad H_2 = \begin{bmatrix} 12 & 4 \\ 2 & 2 \end{bmatrix} \quad H_3 = \begin{bmatrix} 8 & 8 \\ 2 & 2 \end{bmatrix} \quad H_4 = \begin{bmatrix} 16 & 0 \\ 0 & 4 \end{bmatrix}$$

(M. Netter, 1972)⁴

- (a) Which of the above path flow matrices satisfy the principle of user equilibrium?
- (b) Which of the above path flow matrices represent *efficient* equilibria?
- (c) Perturb H_3 in the following way: adjust h_2^A from 2 to 2.1 (so h_2^B becomes 1.9), changing all of the travel times. Then, adjust h_1^A and h_1^B until $t_1^A = t_1^B$ (restoring equilibrium for the first departure time). Then, adjust h_2^A and h_2^B until $t_2^A = t_2^B$ (restoring equilibrium for the second departure time), and so on, alternating between the two departure times, until you reach a new equilibrium solution. Which equilibrium solution do you end up at? Is H_3 stable?
10. [33] Consider an instance of the Daganzo paradox network (Figure 12.12), but with travel times $\tau_2 = 10$, $\tau_{3A} = 4$ and $\tau_{3B} = 1$, capacities of $c_{3B} = 10$ and $c_1 = c_2 = c_{3A} = c_4 = 50$ vehicles per time step, and an inflow rate of $Q = 50$ vehicles per time step. Assume that the maximum number of vehicles which can fit on link 3A is $(k_j L)_{3A} = 300$ vehicles.

⁴If you are wondering how the travel time on a path can be influenced by drivers leaving at a later time, recall that FIFO violations can occur due to phenomena such as express lanes opening, allowing “later” vehicles to overtake “earlier” ones and delay them in the network.

- (a) How many vehicles should be on link 3A at any time to maintain equilibrium?
 - (b) Does this queue fit on the link?
 - (c) The capacity on 3A is now increased to $c_{3A} = 25$. At the new equilibrium, have drivers' travel times increased, decreased, or stayed the same?
 - (d) The capacity is now increased to $c_{3A} = 50$. At the new equilibrium, have drivers' travel times increased, decreased, or stayed the same (relative to $c_{3A} = 20$)?
11. [55] In the network of Figure 12.12, find the equilibrium solution in terms of the capacity on link 3B. Plot the average travel time as this capacity ranges between 0 and 25 vehicles per time step. Assume a time horizon of 20 time steps. (You may do this either by deriving a closed-form expression for the average travel time, or by using one of the algorithms in this chapter to solve for the approximate equilibrium solution.)