Structure and Dynamics of Information in Networks

David Kempe
Department of Computer Science
University of Southern California

July 6, 2011
Preface

The present notes are derived from a course taught at the University of Southern California. The focus of the course is on the mathematical and algorithmic theory underpinning the connections between networks and information. These connections take two predominant forms:

- Network structure itself encodes a lot of information. For instance, friendships between individuals let us draw inferences about shared interests or other attributes (location, gender, etc.). Similarly, hyperlinks between documents indicate similar topics, but can also be interested as endorsements, and thus hint at quality. The list of scenarios in which network structure helps us interpret information about individual nodes continues beyond this list, and will be explored in more detail throughout these notes.

- Networks also play a crucial role in disseminating or gathering information. This applies both to social networks, in which communication between individuals happens naturally, and computer networks, which are designed explicitly to facilitate the exchange of information and distributed computations. We will draw analogies between the two types of networks, and investigate the mathematical underpinnings of the diffusion of information over networks in the later chapters of these notes.

These notes are designed to accompany a one-semester graduate-level course in computer science. We assume a solid mathematical background, and familiarity with basic algorithmic techniques, including flows and cuts; the necessary topics are covered in textbooks by Kleinberg and Tardos [222] and Cormen et al. [101]. In particular, we will assume familiarity with Linear Programming (see, e.g., [88, 201]) and the concept of approximation algorithms [179, 329], as well as with basic notions of randomization in algorithms [258, 266]. A basic understanding of graph theory is also required (see, e.g., [50, 116]).

0.1 Further Reading

Several recent books cover topics overlapping with the content of these notes. “The Structure and Dynamics of Networks” [277] by Newman, Barabási and Watts collects a number of recent papers on network analysis and epidemic processes on networks, predominantly with a physics bent in the type of analysis. The collection “Network Analysis” edited by Brandes and Erlebach [62] covers many of the algorithmic and mathematical foundations of static network analysis.

The recent book “Social and Economic Networks” [190] by Matthew Jackson covers several of the same problems studied here, in particular the diffusion of innovations and modeling and properties of networks. Its focus tends to be on essentially uniformly random models of networks, and mean-field approximations in the style also performed in the physics literature. The book “Social Network Analysis” by Wasserman and Faust [333] covers the tools and concepts employed by sociologists in the analysis of social networks in great detail; the book “Social Network Analysis: A Handbook” by Scott [311] gives a more concise and basic introduction to these topics.

The recent book “Community Detection and Mining in Social Media” [322] covers many of the same topics as the present notes. The focus there is more on the data aspect (whereas the present notes focus more on proofs), so the two sets of notes are complementary.
0.2 Acknowledgments

The course I taught at the University of Southern California was strongly influenced by a similar course taught by Jon Kleinberg at Cornell University (CS685 at the time; now known as CS6850). The selection of topics, style of presentation, and organization owe a great deal to Jon’s insights and teaching styles. Frankly, Jon should be listed as a co-author; unfortunately, his well-known modesty is preventing him from accepting co-authorship.

The foundation of these notes were scribe notes produced by the students of the course CSCI 599 at USC in the Spring of 2005. The students in the class at the time were Farnoush Banaei-Kashani, Shishir Bharathi, Fang Bian, Yuriy Brun, Iftikhar Burhanuddin, Karthik Dantu, Ramki Gummadi, Shyam Kapadia, Salim Khan, Shiva Kintali, Nupur Kothari, Chansook Lim, Fragkiskos Papadopoulos, Animesh Pathak, Ranjit Raveendran, Dustin Reishus, Jyungryun Seo, Viral Shah, Affan Syed, Ashish Vaswani, and Hui Zhang. The notes have also strongly profited from interactions with students who took the class in the Spring of 2008. I would especially like to thank Miquel Camprodon for catching many mistakes in earlier versions of these notes. Nevertheless, there are surely many more mistakes and omissions in them, and I would be most grateful for any feedback.
## Contents

0.1 Further Reading ................................................................. 3  
0.2 Acknowledgments ............................................................. 4  

1 Background 7  
1.1 The early visionary ............................................................ 7  
1.2 Graph structure in the web ................................................... 8  
1.3 Further Reading ............................................................... 11  

2 Web Search using Links, and Spectral Analysis 13  
2.1 Challenges In Searching ....................................................... 13  
2.2 The HITS Algorithm ............................................................ 15  
2.3 PageRank ................................................................. 16  
2.4 Topic-sensitive PageRank ..................................................... 19  
2.5 Text-Based Information Retrieval .......................................... 20  
2.6 Further Reading ............................................................... 22  

3 Community Structure and Clusters in Graphs 25  
3.1 Dense Subgraphs ............................................................... 26  
3.2 Strong Communities ........................................................... 28  
3.3 Bipartite Subgraphs ............................................................ 31  
3.4 Modularity ................................................................. 33  
3.5 Correlation Clustering .......................................................... 35  
3.6 Further Reading ............................................................... 39  

4 Link-Based Classification 41  
4.1 Markov Random Fields and the Optimization Problem ................. 41  
4.2 A 2-approximation using LP-rounding ..................................... 42  
4.3 Uniform Labeling via Local Search ......................................... 45  
4.4 Further Reading ............................................................... 47  

5 Rank Aggregation and Meta-Search 49  
5.1 Rank Aggregation as Social Choice ....................................... 49  
5.2 Aggregation as Optimization .................................................. 52  
5.3 Further Reading ............................................................... 56  

6 Power-Law Distributions 59  
6.1 Preferential Attachment ........................................................ 60  
6.2 Heavy-Tailed Distributions via Optimization ............................. 61  
6.3 A Simple Model for Power Laws in Word Frequencies .................. 64  
6.4 Further Reading ............................................................... 65
7 Small-World Graphs and Decentralized Search
   7.1 Milgram’s Experiment (1967) ......................................................... 67
   7.2 Formalization, and Expander Graphs ............................................... 68
   7.3 The Watts-Strogatz Model ............................................................... 70
   7.4 Further Reading .............................................................................. 73

8 Diffusion of Information and Influence in Social Networks .............. 75
   8.1 The Bass Model for Innovation Diffusion .......................................... 75
   8.2 Schelling’s model for segregation ...................................................... 76
   8.3 Granovetter’s threshold model .......................................................... 76
   8.4 The Contagion model ........................................................................ 78
   8.5 A non-monotone infection model ...................................................... 80
   8.6 Causing a large spread of an infection ............................................... 82
   8.7 Further Reading .............................................................................. 87

9 Epidemic Algorithms ........................................................................... 91
   9.1 Spreading One Message .................................................................... 91
   9.2 Averaging using Gossip .................................................................... 94
   9.3 Further Reading .............................................................................. 98
Chapter 1

Background

1.1 The early visionary

This course is about the interplay between link network structure and information. There are two natural ways in which these can interact: The information can spread by using the edges of the network, or the network itself can be (part of) the information to be analyzed. Naturally, these two views are not mutually exclusive: often, we will want to analyze networks whose function it is to spread information. Similarly, the mere existence of a network to go with other information will often allow for information to spread as well. Lest this discussion be too abstract, let us begin by considering two of the most prominent types of networks in this course:

1. The World Wide Web (WWW) consists of web pages and the links between them. The links add significantly to the (textual and other) information content of the pages. They allow a focused navigation of the body of text. Furthermore, they indicate which pages may be about related topics, or more relevant, by virtue of heavy linking patterns. Thus, the WWW predominantly falls into the second category above. However, since the owners of web pages will also have access to the pages they link to, or pages further down chains of links, information tends to propagate along links as well, as pages copy important information from others.

2. Social Networks are made up of individuals and their relationships, such as friendships, collaborations, romantic relationships, etc. They serve many different functions, including social support, economic exchanges, and — quite crucially — the spread of information. Thus, they fall more into the first category. However, the mere fact that an individual forms links with others tends to imply that they might care about similar information, or share similar interests. Thus, the network structure also reveals a great deal of information about the individuals, beyond what could be discerned with standard socio-economic questions.

The great proliferation of networked bodies of information, and the much more ready availability of data on all kinds of networks, have recently led to a significant increase in interest in computational analysis of such networks. In many scenarios, such analysis, beyond its scientific interest, also has significant commercial or other benefit. For an example, look no further than the importance of web search, which we will discuss in more detail in Chapter 2.

Before embarking on our exploration of many of the algorithmic questions in this broad domain, we should revisit one of the visionary thinkers, who anticipated many of the recent developments and challenges well before computers were used pervasively, or most of the technologies were available.

In his visionary article “As We May Think” from 1945 [69], V. Bush is articulating the next great research challenge for a generation of scientists. The backdrop for the article was the success of the Manhattan Project, in which the development of the Atomic Bomb during World War II managed to unite and focus a large fraction of the scientists of the US on a common project, with significant scientific success. In his article,
Bush articulates his vision that scientists should continue to strive for solving the next big challenge in a similar collaborative effort. The big challenge he wants to see addressed is the vast amount of information available, and the lack of organization thereof. In order to deal with the overwhelming amount of information, Bush suggests having machines which would help with storing, organizing, and physically accessing such information.

Technologically, Bush suggests the use of microfilm for storage, head-mounted cameras for individuals which would let them take pictures easily, and scanner/photography technology. The core of his envisioned technological approach is called “Memex”, and could be considered the equivalent of a modern desk and workstation. It consists of a desk with multiple embedded monitors, which can display text, photos, etc. Books and other documents (pictures, data) are stored in the Memex, and can be easily accessed.

The visionary part is his view of creating associations between documents. When reading, say, a paper, a reader may need to look up a definition in another paper or book. At that point, having opened both documents to the relevant pages, a button would allow the reader to create an association between the two documents. Thus, Bush really anticipated hyperlinks here. His reasoning is that the rigid structure of a library index does not correspond to the associative type of thinking performed by humans. Bush anticipates not only creating such associations, but also sharing them with others, so that others will not have to perform the same type of work. He calls an annotated document a trail. Given the importance of such trails for others, he even anticipates a profession of trail blazers, whose job it is to read and process documents, and create associations helpful for others. Notice how virtually all of these visions have come to pass in the WWW: the WWW itself is a giant shared trail, and many sites (Yahoo!, Open Directory, Wikipedia) can be construed as trail blazers.

Bush already anticipates a number of important research challenges in this context. Besides the obvious technological ones, two important directions are privacy and automated processing. If trails, or personal information, are shared, then it becomes important to specify what types of information should not be leaked, and which ones are safe to share. Privacy in data mining has become a huge research area, albeit one which this course will not touch upon. The second question is how to automatically create associations, or mine the associations for useful information that was not obviously contained in the data originally. This question is becoming even more paramount as the amount of information in the WWW keeps increasing dramatically, to the point that the trails themselves are now so numerous that the reader may not be able to identify the most pertinent ones. These types of questions will be among the core topics of this course.

1.2 Graph structure in the web

The most prominent example of a large body of networked information is the World Wide Web. It quintessentially meets the focus of this course: the hyperlinks between pages containing textual information (or other formats) naturally enrich the information content, and allow us to draw additional inferences about the content, its importance, or relationships between multiple pages. Given the central importance of the web, we begin with an exploration of some of its basic graph-theoretic properties.

The presentation here closely follows the paper “Graph Structure in the Web” by Broder et al. [65]. It treats the World Wide Web as a directed graph, whose nodes are the static pages, and whose edges are the hyperlinks between them. This graph is also called the Web graph. The reported data are from a crawl in May 1999, and thus several orders of magnitude smaller than they would be today. Nevertheless, it is interesting to investigate some of the relative sizes of different parts of the Web.

1.2.1 Connectivity Structure

The crawl user by Broder et al. [65] contained ca. 203 million nodes and 1.5 billion edges. Of these, roughly 91% were contained in one large weakly connected component. The fact that there is only one such component is not surprising since there is only one World Wide Web. It would be rather surprising to have two completely disconnected large components. The exact percentage of pages in the large component should be taken with a grain of salt. After all, one may argue that small components not connected to or
from the “bulk” of the Web may have a much smaller chance of being discovered.

In looking at these numbers, one may suspect that a few highly connected “hubs” are really in charge of connecting all of the other nodes to each other. This suspicion turns out to be incorrect. Even when all nodes of degree more than 5 are removed, there is still a large component with more than 59 million pages. This shows that the connectivity structure is highly “decentralized”, and utilizes most of the Web.

Looking at strongly connected components (SCCs), the picture is somewhat different. There is still only one strongly connected component (SCC), containing approximately 56 million pages that are mutually reachable. Again, the existence of only one large SCC should not be surprising, for the same reason as the weakly connected component. In addition to the giant SCC, the strong connectivity structure contains the parts labeled IN, OUT, tendrils, and tubes in the diagram; they are all discussed below.

Figure 1.1: Connectivity of the web (diagram taken from [65])

Figure 1.1 gives an overview of the structure. IN denotes the set of all nodes from which the giant SCC is reachable, but which in turn are not reachable from SCC. Similarly, OUT is the set of all nodes reachable from SCC, but not able to reach SCC themselves. Each of them contains about 44 million pages.

Hanging off IN and OUT are tendrils containing nodes that are reachable from portions of IN, or that can reach portions of OUT, without passage through SCC. Tubes are made up of nodes on paths from IN to OUT, but not passing through the giant SCC. Tubes and tendrils together comprise another 44 million pages. Finally, there are some completely disconnected strongly connected components, whose contribution to the total size of the Web is rather negligible.

The depiction of the strong connectivity structure in Figure 1.1 resembles a bow tie, and the connectivity structure has therefore often been referred to as “bow tie structure” of the web. Notice, however, that every directed graph has this structure when drawn with respect to one particular strongly connected component. The bow tie structure really carries no information about the Web specifically; what is more interesting is the relative sizes of the different components.

Notice that the size of some components, most notably IN and the tendrils of OUT, are more likely to be underestimated than the others. The reason is that these components suffer from a “seeding” problem: how to discover a new page if it is not reachable from an already crawled one? Some solutions to this problem are manual submission of web pages and random generation of names, but nevertheless, one would expect that these two components are more undersampled than those which are by definition reachable from known nodes.

Despite these sampling issues, the above structure and relative sizes can be expected to be quite stable,
even to missed pages. The reason is that any paths from OUT to IN, from SCC to IN, or from OUT to SCC, would already have been discovered if they existed. So these components will not be contracted.

An interesting consequence of the numbers depicted above is the fraction of node pairs that can reach each other. If we select two nodes at random, here will be a directed path from one node to the other with probability roughly 30%. This can be computed by a simple case distinction over the different components in which the two nodes could land. For such a path exists either when both nodes are in SCC, or one is in IN and the other in SCC or OUT, or the first one in SCC and the second in OUT. (There are some additional, but much less frequent, cases involving tendrils and tubes.) Looking at the sizes of the components then shows that the probability of such selections is about 30%, contrary to the once popular belief that almost all node pairs are connected to each other via directed paths.

### 1.2.2 Path Lengths

Once we have established the fraction of page pairs which are connected to each other, a second interesting question becomes how many hops separate an average or worst-case pair of pages. These can be computed using invocations of BFS from all start nodes, in time $O(nm)$. Unfortunately, given the orders of magnitude of $m$ and $n$, this running time is too slow. No faster algorithm for calculating the diameter of a graph is currently known.

Broder et al. [65] instead randomly sample node pairs, and measure their distance. This provides a lower bound on the diameter (since the maximum distance may have eluded the algorithm), and a fairly accurate estimate of the average distance. Using these techniques, Broder et al. gave a lower bound of 28 on the diameter of SCC, and a lower bound of 503 on the diameter of the WWW. The much larger bound on the entire WWW is likely due to long chains of pages (chapters of a book, for instance), which would likely be either in the OUT component, or inside tendrils. The average distance between sampled pairs of nodes in the SCC was found to be roughly 16 if the edge direction was considered and roughly 7 if the edge direction was disregarded.

Given that time $O(nm)$ is too slow for extremely large graphs for computing the diameter, an interesting approximation algorithms question would be for which types of graphs random sampling gives a good lower bound on the diameter of the graph, i.e., trying to find a structural characterization of graphs for which many node pairs have nearly the maximum distance.

### 1.2.3 Degree Distributions

Another interesting statistic is the degree distribution of web pages. How many links go into or out of the “typical” page? Broder et al. [65] denote by $p_i$ the number of pages with $i$ incoming links, and by $q_i$ the number of pages with $i$ outgoing links. Their analysis finds that $p_i \sim i^{-2.1}$, and $q_i \sim i^{-2.72}$, for several orders of magnitude of $i$. Thus, both the indegree and outdegree seem to follow a power law.

**Remark 1.1** A power law is a degree distribution of the form $p_i = C \cdot i^{-\alpha}$, for some constants $C, \alpha > 0$. It can be recognized “heuristically” by drawing the distribution in a log-log plot. Then, $\log p_i = \log C - \alpha \log i$, so the power law appears as a straight line with slope $-\alpha$ and intercept $\log C$. In order to be meaningful, a power law should normally manifest itself over at least three orders of magnitude. We will discuss power laws in much more detail in Chapter 6.

How interesting is it to observe a power law degree distribution? By itself, a power law does not tell us much about the structure of the graph: there are still many completely different graph structures which would all have the same degree distribution. In particular, the term “power law graph” does not really carry much meaning.

However, one sense in which the power law observation is somewhat interesting is that it clearly rules out the simplest graph model — namely the Erdős-Rényi $G(n,p)$ model [129] as an accurate description of the web graph. In an Erdős-Rényi graph of $n$ nodes, each edge is present independently with probability $p$. The result is that the degree distributions are Binomial. In particular, the degrees will be sharply concentrated around their expectation $pn$. Erdős-Rényi graphs are often used as default models, absent a
better understanding of the graph structure, and the degree distribution observed by Broder et al. rule out that $G(n, p)$ could be an accurate model.

We should also stress here that the interesting part of a power law is that it captures a fairly large number (only polynomial decay) of nodes with high values of $i$, i.e., large in-degree or out-degree. Not too surprisingly, the indegrees have a heavier tail, since it requires less effort to receive a significant number of incoming links than to create a significant number of outgoing ones.

1.3 Further Reading

The early history of the WWW, including Bush’s paper and many other position documents and milestones, are described in a short history of the WWW from 1945–1995, available online [99]. For a personal account of the early days and invention of the Web from one of its inventors, Tim Berners-Lee, see his book “Weaving the Web” [37].

Another early paper analyzing graph properties of the Web is [220]. In fact, it also presents some graph models discussed later, web search algorithms, and algorithms for community detection. Beyond degree distributions and path lengths, it analyzes the frequencies of certain dense subgraphs, and connectivity properties of local subgraphs.

Some of the mathematical issues that will be explored in this course are described nicely in a survey of research directions accompanying a talk given by David Donoho [121].
Chapter 2

Web Search using Links, and Spectral Analysis

In this chapter, we investigate the use of web graph structure to improve web search. This application cuts straight to the heart of this course, in that the network structure augmenting the textual information contained in web pages can be used to significantly improve the performance of information extraction algorithms. Our focus, after a brief discussion of the issues encountered in web search in general, will be on two algorithms: Hits and PageRank. The latter underlies the successful search engine Google.

Both of the algorithms are based on computing eigenvectors of certain matrices representing the Web graph. The reader unfamiliar with (or rusty on) basic linear algebra might find this a good time to consult standard linear algebra texts such as Strang’s [319] or Horn and Johnson’s [185]. The eigenvector based algorithms in PageRank and Hits also bear some similarity with techniques used in analyzing textual data, specifically Latent Semantic Analysis. We therefore turn our attention to that topic in Section 2.5.

2.1 Challenges In Searching

Writing a search engine involves many challenges. The first type are technical: writing a crawler that deals with the many types of formats, and — more challenging — all the types of errors human authors of web pages introduce, is itself a complicated task. In addition, crawlers also have to avoid overloading web servers and deal with connection timeouts. A good crawler needs to repeat crawls frequently, to avoid the problem of having stale or outdated information available in searches. This is particularly important for searches about recent or upcoming events, to which the correct answers changes frequently.

A more recent important trend is the development of the “deep web”: information which is not readily available in web pages, but rather stored in databases, and accessible only through query forms. It is not clear how such information could (or should) be accessible to search engines, and how it can be returned in response to queries.

Designing the actual search engine to deal with the large numbers of queries encountered in practice is also a significant engineering challenge, and solving it successfully has played as large a part as algorithmic insights in the success of current search engines. However, in this course, we will focus more on the algorithmic challenges inherent in answering different types of queries. For this purpose, it is useful to have a brief look at the different types of queries commonly posed to search engines. This list is neither exclusive nor exhaustive, but highlights some of the issues a web search engine designer will have to take into account.

1. Specific queries, such as “Does Netscape support the JDK 1.1 code-signing API?” As we will discuss in detail below, the difficulty in handling specific queries is scarcity: there are very few or no pages which contain the required information or keywords, and it is often difficult to determine the identity of these pages.
2. **Broad-topic queries**, such as “Find information about the Java programming language.” Here, we have the opposite problem: the number of pages that could potentially be relevant is far too large for a human user to digest. We will need to devise techniques to narrow down or rank the set of pages returned.

3. **Similar-page queries**, such as “Find pages ‘similar’ to java.sun.com.” The problem here is to determine in what respects the pages are supposed to be similar to the present one. This requires, at least implicitly, forming an estimate of the topic or content of a page.

4. **Current events.** Here, the challenge is that the data base needs to be very up to date, and hence the crawling needs to be regular. For such dynamic changes, one might fear that link structure will not change fast enough to help improve accuracy.

5. **Navigational queries**, such as “Home page of the Sun corporation.” Such a query essentially replaces the formation of a bookmark as navigational help. Therefore, it is important that the best fit be clearly identified and listed first.

6. **(Personalized) Opinion**, such as “The best hotel under $100 in Los Angeles.” Here in particular, the search engine might even need to know more about the opinions of the searcher, but also needs to aggregate opinions expressed in pages to form some kind of consensus.

In dealing with these different types of queries, the first broad class of algorithmic challenges is the “classic” information retrieval (IR) problem of extracting useful information from a large body of text or other sources. For instance, different authors use different formats (html, pdf, xml, jpg, etc.), different authoring styles, or different terminology to refer to the same idea or object. Extracting the meaning or intent of text has been a traditionally hard problem, and is a research area of its own. We will discuss a very basic technique in Section 2.5, but will otherwise not make any attempts to cover this area in this class.

In the context of web search, the traditional IR problem is further complicated by the fact that the authors of the individual pages may have goals vastly differing from those of the search engine or the searcher. Many web sites are actively misrepresenting their content or importance in order to attract more visitors. Common techniques include adding interesting words even when they do not relate to the page’s content (called content spamming), or linking to a page from other pages such as wikis or blogs (called link spamming), in order to make the page appear more relevant. As a result, there has been a lot of recent research work on dealing with such web spam, and more generally, the area of adversarial information retrieval has been very active.

Here, we will spend more time exploring solutions to the problems caused by specific and broad queries. A sparse query may have only few (or no) results, even though the searcher’s intent would result in relevant answers. A broad query might have thousands or millions of results, and it becomes imperative to extract the ones most likely to be useful. Both commonly occur because the user — not knowing which pages will answer his query — does not know which particular keywords to enter to have the engine return these queries.

The problem of sparse queries can be at least partially solved using traditional IR techniques, such as augmenting a search by known synonyms, or other words frequently used by other searchers together with the search terms. However, the link structure provides perhaps much more useful additional information. For instance, imagine that the search is for “car manufacturers”. The searcher’s intent is probably to find the websites of Toyota, Honda, etc. However, some or many of these sites may not contain the exact search terms. To deal with this issue, we can notice that there will likely be many websites containing the phrase “car manufacturers”, and linking to the sites of actual car manufacturers. Thus, using information contained in linking (or linked-to) pages, we can infer additional information about a web site.

Link structure may also help us deal with the abundance problem: the number of pages containing a query can be extremely large, and must be ordered by relevance. To estimate a page’s relevance, we can consider the number and nature of links it receives, while positing that incoming links implicitly express some notion of endorsement. Thus, other things being equal, pages with more links could appear more relevant. These ideas essentially underlie most of the successful search engines used today, and we will elaborate on the underlying mathematics in more detail in the next sections.
2.2 The HITS Algorithm

The HITS algorithm was suggested by Kleinberg [216], and is one of the first algorithms to use link structure in order to deal with the abundance problem. It has two components: first, it constructs a focused subgraph of the web; second, it uses “hub” and “authority” weights to identify particularly relevant pages.

The first component of HITS deals with the problem mentioned above in the context of sparse queries: that some of the most relevant pages might not contain the query terms. To deal with this issue, the algorithm first identifies a core set of approximately 200 pages using standard text-based search. This set is expanded by adding all the pages linked to from the core set, and some pages linking to the core. Since many pages could conceivably link to one page in the core set, the number of in-links is restricted to about 50, which are chosen heuristically. In total, this generates a focused induced subgraph \( G \) of \( n \approx 3000 \) nodes.

The second component of HITS aims to deal with the issue of broad searches, and attempts to find the most “relevant” pages in \( G \). It is based on the observation that a page could be relevant for two reasons: (1) it could be a good authority, containing genuine information about the topic, or (2) it could be a good hub, collecting links to many good authorities.

The distinction, and the introduction of hubs, requires a bit more justification. Why are we not simply considering authorities and their linking behavior to each other? In fact, we will consider this approach in the context of the PageRank algorithm (which forms the core of Google’s ranking system) in Section 2.3. One justification for considering hubs in addition to authorities is that much of the World Wide Web forms a competitive environment, in which multiple authoritative pages on the same topic may not link to each other. For example, in the context of “car manufacturer” web pages discussed previously, we would not expect the most authoritative pages — such as Honda, Toyota, or Ford — to recommend each other to a customer. Any connection between these pages will have to be indirect, by virtue of other pages considering all of them authoritative.

To identify which authorities are good, we use the observation that good authorities should be pointed to by (many) good hubs. In turn, a hub is good if it points to (many) good authorities. The tacit assumption here is that links constitute an endorsement. The definition of what constitutes a good hub or authority is circular. However, we can translate it into an actual algorithm as follows.

Each page \( v \) has authority weight \( a_v \) and hub weight \( h_v \), initialized to 1. In any one iteration, these are updated, by having authorities inherit authority weights from the hubs that point to them, and hubs inherit hub weight from the authorities they point to. Thus, one iteration updates

\[
\begin{align*}
a'_v &= \sum_{u \rightarrow v} h_u & \text{for all } v \\
h'_v &= \sum_{v \rightarrow u} a'_u & \text{for all } v.
\end{align*}
\]

Notice that after one iteration, the authority weight of \( v \) is the number of pages pointing to \( v \), which is intuitively related to the authority of a page. Further iterations refine this notion by giving more weight to hubs that are deemed more relevant.

If we let \( a, h \in \mathbb{R}^n \) denote the vectors of all authority and hub weights, we can express the same operations in terms of matrices as follows. Let \( B \) be the adjacency matrix of \( G \), i.e., \( B_{ij} = 1 \) iff there is an edge \( i \rightarrow j \). Then, we can rewrite the update rules as

\[
\begin{align*}
a' &= B^T \cdot h \\
h' &= B \cdot a'.
\end{align*}
\]

Hence, the update rule for authority weights can be written as \( a' = (B^T B) \cdot a \), and for hub weights as \( h' = (B B^T) \cdot h \). The matrix \( B^T B \) is called the co-citation matrix. We can understand it intuitively as counting the number of pages pointing to both \( i \) and \( j \).

To ensure that this update rule converges to a solution, we need to normalize the weights after each update, for instance such that \( \sum a'^2_v = 1 \) and \( \sum h'^2_v = 1 \) (otherwise, the values would keep growing).
We will revisit the issue of whether the normalized rule converges in a moment. However, we first want to explore what it would converge to, if at all. Convergence means that subsequent iterations do not cause any changes, so \( a, h \) would have to be fixed points, i.e., \( a = \beta a \) and \( h = \gamma h \), where \( \beta, \gamma \) are the respective normalization constants ensuring that \( \sum (a_i')^2 = 1 \) and \( \sum (h_i')^2 = 1 \).

Substituting the expressions we derived above gives us that \( \beta \cdot a = (B^T B) \cdot a \) and \( \gamma \cdot h = (BB^T) \cdot h \). Thus, the authority weights are exactly an eigenvector of \( B^T B \) with eigenvalue \( \beta \), and the hub weights are exactly an eigenvector of \( BB^T \) with eigenvalue \( \gamma \).

Now, we can return to the issue of convergence of the iterative process (along with the question which of the eigenvectors are returned). Since the matrix \( B^T B \) is symmetric, \( \mathbb{R}^n \) has a basis of orthonormal eigenvectors \( \omega_1, \omega_2, \ldots, \omega_n \) of \( B^T B \). Let \( \lambda_1, \lambda_2, \ldots, \lambda_n \) be the corresponding eigenvalues with \( |\lambda_1| \geq |\lambda_2| \geq \cdots \geq |\lambda_n| \). If the graph corresponding to \( B^T B \) is connected, then by the Perron-Frobenius Theorem for non-negative matrices, \( |\lambda_1| > |\lambda_2| \), \( \omega_1 \) has all non-negative entries, and each \( \omega_i \) for \( i > 1 \) has at least one negative entry. Because the eigenvectors form a basis, we can write the starting vector of authority weights as \( a_0 = \sum_{i=1}^n \alpha_i \omega_i \), and the previous properties imply that \( \alpha_1 \neq 0 \). The next iteration of authority weights then gives
\[
\begin{align*}
a_1 &= (B^T B) \cdot a_0 = \sum_{i=1}^n \alpha_i (B^T B) \omega_i = \sum_{i=1}^n \alpha_i \lambda_i \omega_i. \\

\end{align*}
\]

By induction, we can now show that for any \( k \),
\[
\begin{align*}
a_k &= \sum_{i=1}^n \alpha_i \lambda_i^k \omega_i. \\

\end{align*}
\]

Hence, as \( k \to \infty \), the normalized authority weights converge to \( a_\infty = \omega_1 \). Therefore, authority weights are just the first eigenvector of \( B^T B \). By an identical argument, the hub weights converge to the first eigenvector of \( BB^T \).

In retrospect, we can thus think of the authority weight computation as using the power iteration method (see, e.g., \([165]\)) to compute the first eigenvector of \( B^T B \). Of course, one could instead use different techniques to compute the eigenvector. Usually, power iteration is quite fast. However, as we can see from Equation (2.1), its convergence rate depends on the spectral gap \( |\lambda_1| - |\lambda_2| \). In fact, if \( |\lambda_1| = |\lambda_2| \), the method may not converge at all. On the other extreme, if the spectral gap is bounded by a constant, independent of the number of nodes \( n \), then the power iteration will converge exponentially fast. (To see this, simply use triangle inequality for the difference between \( a_k \) and \( \omega_1 \).)

### 2.2.1 Extensions

Having been so successful with interpreting the top eigenvector of the co-citation matrix and its transpose, we may be interested in interpreting its other eigenvectors as well. It turns out that they can be considered as identifying interesting hub-authority structure in subgraphs of \( G \), and will prove useful in identifying more specific topic communities within \( G \). We will return to this issue in Section 3.3.1.

The approach of HITS can be easily extended to help with finding pages related to a given page. Dean and Henzinger [109] suggest simply building the focused subgraph starting at the given page \( u \). All pages pointing to \( u \), and pointed to by \( u \), are included (up to a given size limit). In addition, all pages pointed to by parents of \( u \), and pointed to children of \( u \), are also included, up to a degree constraint per page. ([109] suggest heuristics for deciding which links to retain if the degree of a page is too large.) Once the focused subgraph has been determined in this way, one simply computes the authority scores and outputs the top authorities. Naturally, looking at web graph communities containing the page \( u \) is another way to uncover related pages (see, e.g., Section 3.3.1). We will investigate communities in more detail in Chapter 3.

### 2.3 PageRank

In the previous section, we looked at and analyzed the HITS algorithm, based on Hubs and Authorities. HITS was designed with a scenario in mind where authorities might not cite each other. Hence, HITS uses
the idea of conferring authority indirectly through hubs. In scenarios where authorities do cite each other, such as academic research work, it may be more appropriate to have direct conferral of authority between nodes. This is the idea behind the PageRank [63] algorithm.

Which algorithm is “better” may be hard to evaluate, and will depend on the query, among others. It raises the more general question of how one can even evaluate the performance of an algorithm when its objective function is not clearly defined.

2.3.1 A first approach

Based on our intuition above, we would like to say that a page has high PageRank if it is pointed to by many pages of high PageRank. On the other hand, if a page points to many other pages, it presumably will not confer a lot of authority to all of them, but rather divide its authority evenly. This suggests the following update rule, starting from an arbitrary vector, such as \( p_i = 1/n \) for all \( i \): the new PageRank is

\[
p'_i = \frac{1}{d^+(j)} \cdot p_j,
\]

where \( d^+(j) \) is the out-degree of \( j \).

We can immediately see that at any fixpoint of this update rule, the solution has to be the solution to a linear system in \( n \) variables. We could also converge to the fixpoint by iteratively applying the update rule.

A more concise way of writing the update rule can be obtained as follows. Let \( M \) be the square matrix with rows and columns corresponding to web pages, and \( M_{ij} = \frac{1}{d^+(j)} \) if \( i \to j \) and 0 otherwise. With \( p_k \) denoting the vector of PageRank weights after \( k \) iterations, the update rule can be written as

\[
p_{k+1} = M^T \cdot p_k.
\]

Notice that the operation has the useful property of “Mass Conservation”: the sum of the PageRank weights of all pages is always exactly 1, as it is only redistributed along outlinks.

2.3.2 A Problem and Solution

The above solution has a problem: it may not converge, and when it does, the result may not be what we are looking for. Specifically, if a web page has no outlinks (except, say, one to itself), then it will never pass on PageRank weight to any other nodes, but it will receive PageRank weight. Hence, all of the weight will collect at sinks, or, more generally, in the sink strongly connected components. All other nodes will have PageRank 0, which is certainly not corresponding to our intuition. In addition, which sinks will end up with all the weight will depend on the starting vector, so the PageRanks are not independent of the starting assignment, which is another undesirable property. Finally, in the case of a cycle of two nodes, the weight will oscillate between those nodes, but not converge.

There is a simple way to fix all of those problems at once. To motivate the approach, we notice that we can view \( M \) as the matrix of transition probabilities of a random walk on the web graph. When the random walk is at a node \( i \), it chooses uniformly among all outgoing links, and follows that link. Hence, the first attempt above corresponds to computing the probabilities of being in a given state \( i \) in the limit of infinitely many steps. We can now modify the random walk as follows: for some small \( \epsilon \approx 1/7 \), with probability \( (1-\epsilon) \), the new Markov Chain does exactly the same as the old random walk. With the remaining probability \( \epsilon \), the new Markov Chain chooses a uniformly random vertex and jumps to it. This random process can intuitively be motivated by the model of a random surfer who gets bored and jumps to a uniformly random page with probability \( \epsilon \) in each step.

Remark 2.1 While this random surfer model provides a reasonable motivation for the Markov Chain approach, using it in computation should not be taken as implying that the model is an accurate representation of actual surfers. Nor should it be assumed that modeling the behavior of a human surfer accurately would necessarily result in a good search engine. After all, search engines are supposed to augment human search capacity. It is better to consider the Random Surfer model as a decent intuition for a technical hack necessary to make a Markov Chain ergodic.

In terms of matrices, we can express the new process as follows. Let \( 1_{p \times q} \) denote the \( p \times q \) matrix with all ones. The update step from the distribution \( p_k \) (at some time \( k \)) to \( p_{k+1} \) can now be written as follows:

\[
p_{k+1} = (1-\epsilon) M^T \cdot p_k + \epsilon \cdot \frac{1}{n} 1_{n \times 1} = ((1-\epsilon) M^T + \epsilon \cdot \frac{1}{n} 1_{n \times n}) \cdot p_k =: M' \cdot p_k.
\]
To show that this new version of PageRank converges, we can use the following well-known theorem about Markov Chains.

**Theorem 2.2** If the Markov Chain $M$ is irreducible and aperiodic (the gcd of all cycle lengths is 1), then it converges to a unique stationary probability $\pi$, i.e., $\pi = M^T \cdot \pi$.

The term “irreducible” means that the corresponding graph (of non-zero entries of $M$) is strongly connected. Certainly, our matrix $M'$ is both irreducible and aperiodic, as we added a jump probability of $\epsilon/n$ between any pair of vertices, i.e., the graph is now complete. (Notice that for $M$ that are not aperiodic, oscillations might occur, as in the example mentioned above of a cycle of length 2. Similarly, the probabilities of the random walk may depend on where the random walk started. This also applies if $M$ is not irreducible. On the other hand, if $M$ is irreducible and aperiodic, then the probabilities are independent of the starting vector.)

**Remark 2.3** For students who have seen this before, a very simple Markov Chain coupling argument (about the simplest possible) shows that the mixing time of this chain is $1/\epsilon$, i.e., a constant. Simply couple two copies of the chain as follows: for both chains, use the same random choice for whether to jump randomly or follow an edge. If following an edge in both, make the choices independently (unless the chains are already in the same state, in which case make the same choice). If jumping uniformly, couple the choices to have both chains jump to the same node. It is easy to see that each chain individually is a faithful copy of the original Markov Chain $M'$. Also, the expected time until the chains jump uniformly is $1/\epsilon$, and after that point, both chains will always be in the same state. Thus, the expected time to couple is constant.

As a result of the constant mixing time, the convergence is exponentially fast: within $O(\log \frac{1}{\delta})$ iterations, the error is at most $\delta$. Regarding computation, we are now dealing with a very large and dense matrix. However, the fact that most entries are equal to $\epsilon/n$ allows for efficient matrix computation techniques to be applied nevertheless.

The stationary probability $\pi$ of the matrix $M'$ is again the top eigenvector, much like for the HITS algorithm. Thus, in principle, we could apply techniques other than the original update rule (which again corresponds to the power iteration). However, given the exponentially fast convergence of power iterations, there is really no need for other techniques. Of course, in practice, to perform computations of this size, several computational tricks need to be applied.

We have not yet described how the PageRank algorithm uses the values $p = \pi$. At query time, the set of candidate pages is narrowed down by using text-based matches, as well as several other heuristics. The other pages are then simply sorted by decreasing $p_i$ values. (In practice, Google uses significantly more heuristics to determine the ordering, but the above is a first approximation of the original approach in the actual search engine.)

The attentive reader will have noticed two differences between the descriptions of PageRank and HITS: using the adjacency vs. co-citation matrix, and using the entire web graph vs. a focused subgraph. These two issues are really orthogonal, and one can consider combining the “focused subgraph” of Section 2.2 with computing the PageRank measure at query time. The result would likely yield more relevant query results than plain PageRank. However, it would require recomputing the corresponding PageRank values at query time, which is not feasible for practical applications.

If an algorithm such as HITS or PageRank does use focused graphs, there are several more natural heuristics by which one can improve the quality of search results. For instance, different weights can be associated with different links, depending on how much relevant text (e.g., query terms) appears close to the link’s anchor. Similarly, different weights can be assigned to different pages based on content-based similarity to the query topic. Both Chakrabarti et al. [73] and Bharat and Henzinger [41] suggest various heuristics to use the content of pages to alter the matrix, and obtain improved search results. Chakrabarti et al. [73] found that most of the relevant text of links occurs between 25 bytes before and 50 bytes after the anchor, and suggest taking such weights into account. Bharat and Henzinger [41] suggest using the first set of textual matches as a “comparison point” for the expanded focused graph, and prune or weigh down pages which are very different.
2.4 Topic-sensitive PageRank

As discussed above, computing topic-specific PageRank values for each query at query time is too time-consuming to work on the scale required by search engines nowadays. On the other hand, pre-computing the PageRank values for all queries and storing them is too memory-intensive. Using the same PageRank values regardless of the specific query can yield unsatisfactory results.

A middle path has been proposed by Haveliwala [177]. The idea is to precompute PageRank values for a few “landmark” topics, and then express queries as a combination of these topics. Calculating the PageRank values from the pre-computed ones may then be significantly easier than recomputation from scratch.

Assume that there are $T$ topics, and each topic $t$ is characterized by its restart probability vector $f_t$, i.e., the probability distribution (over pages) with which a new random page is chosen when the Markov Chain jumps. The idea of using topic-specific restart probabilities was proposed by Rafiei and Mendelzon, who suggest jumping to a uniformly random page containing the search term [295]. (The context of their work is computing topic-specific PageRank scores to find out which topics a page is particularly known for.) The corresponding update rule is then $p’ = (1 - \epsilon)M^T \cdot p + \epsilon \cdot f_t$. By defining the matrix $F_t = [f_t f_t \ldots f_t]$, we can express this new rule as

$$p’ = (1 - \epsilon)M^T \cdot p + \epsilon \cdot f_t =: M_t \cdot p.$$  

The corresponding PageRank values for topics $t$ are now the stationary probabilities $\pi_t$ of $M_t$.

From the PageRank values for topics $t$, we can compute those for queries $q$ as follows. We can consider a query as a convex combination of topics, i.e. $q = \sum t \gamma_t \cdot t$ (where $\sum t \gamma_t = 1$). Then, we can identify with $q$ the reset vector $f_q = \sum t \gamma_t f_t$. The interesting observation is that the corresponding stationary probabilities $\pi_{f_q}$ can be obtained as convex combinations of the probabilities for the landmark topics:

**Lemma 2.4** For each query $q$ with $f_q = \sum t \gamma_t f_t$, we have $\pi_{f_q} = \sum t \gamma_t \cdot \pi_t$.

**Proof.** We will show that the vector on the right is stationary for $M_q$. As the stationary distribution is unique, this proves that it is equal to $\pi_q$. In order to do so, we use the fact that each $\pi_t$ is stationary for its corresponding $M_t$, and then use the linearity of matrix-vector multiplication and summation:

$$\sum t \gamma_t \cdot \pi_t = \sum t \gamma_t \cdot ((1 - \epsilon)M^T \cdot \pi_t + \epsilon \cdot f_t) = (1 - \epsilon)M^T \sum t \gamma_t \cdot \pi_t + \epsilon \sum t \gamma_t f_t = (1 - \epsilon)M^T \sum t \gamma_t \cdot \pi_t + \epsilon f_q = M_q \cdot \sum t \gamma_t \cdot \pi_t.$$  

Hence, the PageRanks of linear combinations of topics can be efficiently computed from precomputed topic PageRanks.

2.4.1 Single-Word Queries

To go even further, one could try to pre-compute the PageRanks for every single-word query. At first, this may seem very daunting, as the number of words is far in excess of 100000, and hence, it appears as though the storage requirement would be larger than $10^5 \cdot 10^9 = 10^{14}$ PageRank values. However, more careful indexing may reduce this requirement significantly, as most words do not appear in most pages. A simple back-of-the-envelope calculation observed by Domingos and Richardson [120] goes as follows. Let $w$ denote a word, and $i$ a page. Further, let $p_w$ be the number of pages containing $w$, and $s_i$ the number of words contained in page $i$, and $x_{w,i} = 1$ if word $w$ appears in page $i$. Then, the total required index size is

$$\sum_w p_w = \sum_{w,i} x_{w,i} = \sum_i s_i.$$  

Notice that the average page contains only about a few hundred words, so the last sum is only about a few hundred times the number $n$ of pages in the web. While this is not yet quite feasible, it is not too far removed from current technology.
For multi-word queries, the situation is a lot more bleak: the number of distinct pairs that appear in web pages is much higher, and the same kind of simplistic analysis does not work: the sum now becomes \( \sum_i s_i^2 \), which may be much larger. Devising good precomputation techniques for multi-word queries seems like a promising direction for further research.

### 2.5 Text-Based Information Retrieval

We have so far hinted several times at the importance (and difficulty) of text-based information retrieval. Clearly, it plays a key role in web search. It also turns out that some of the core techniques in information retrieval use ideas similar to the eigenvector-based ones we discussed for HITS and PageRank.

The most naive search strategies for text-based information retrieval would just return pages that contain the query term(s). There are two problems with this approach: synonyms (different words having the same meaning) and polysemy (the same word having more than one meaning). Synonyms may lead the search engine to miss relevant pages, because the exact query terms may not appear in all relevant pages. Polysemy may lead the search engine to return irrelevant pages; the pages may contain the search term, but in a different context and with a different meaning than the user intended.

Consider the following table. There are 6 pages, and 6 words occurring on the pages. Page 1 contains words 1, 2, and 3, and so on. Imagine searching these pages for word 3. A naive search engine would simply return pages 1, 2, and 4, because these pages all contain the query word. However, we might argue that page 3 is also relevant: while it does not contain the exact query term, it is very similar to pages 1 and 2 that do contain the query term. Similarly, page 4 may not be considered (as) relevant: while it does contain the query word, it is very similar to pages 5 and 6 that do not contain the query word.

<table>
<thead>
<tr>
<th></th>
<th>Word 1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
<th>6</th>
</tr>
</thead>
<tbody>
<tr>
<td>Page 1</td>
<td>x</td>
<td>x</td>
<td>x</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td></td>
<td></td>
<td>x</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td></td>
<td>x</td>
<td></td>
<td>x</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td></td>
<td></td>
<td>x</td>
<td>x</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td></td>
<td>x</td>
<td>x</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td></td>
<td>x</td>
<td></td>
<td>x</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td></td>
<td>x</td>
<td></td>
<td>x</td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

To put these observations on a more rigorous and general footing, we can use techniques from Spectral Analysis of Data [23], which is also known as Latent Semantic Analysis [111] in the Information Retrieval community. We consider the table as a matrix, where the cells with an ‘x’ are 1, and the cells without an ‘x’ are 0. This matrix is called the *term-document matrix*.

In general, this matrix can be used for representing a variety of data sets, where rows index objects in the data set, columns index attributes of those objects, and the \([i,j]\) entry of the matrix represents the value of the \(j^{th}\) attribute in the \(i^{th}\) object. Some examples of interest are where both rows and columns refer to web sites and the \([i,j]\) entry indicates that site \(i\) has a link to the site \(j\); another is that rows index individuals, columns index products, and the \([i,j]\) entry indicates whether individual \(i\) is interested in, or has previously purchased, product \(j\).

The main tool in extracting the latent structure from the matrix \(A\) is the *singular-value decomposition* (see, e.g., [185]):

**Lemma 2.5** Let \(A\) be an \(m \times n\) matrix (e.g., document-term matrix). Then, \(A\) can be written as \(A = U \cdot \Sigma \cdot V^T\), where \(U\) is an \(m \times k\) orthonormal matrix, \(V\) is an \(n \times k\) orthonormal matrix (\(k = \text{rank}(A)\)), and

\[
\Sigma = \begin{bmatrix}
\sigma_1 & 0 & 0 & 0 \\
0 & \sigma_2 & 0 & 0 \\
0 & 0 & \ddots & 0 \\
0 & 0 & 0 & \sigma_k
\end{bmatrix}
\]
is a $k \times k$ diagonal matrix with $\sigma_1 \geq \sigma_2 \geq \cdots \geq \sigma_k > 0$. This is called the singular value decomposition (SVD) of $A$, and the $\sigma_i$ are called the singular values.

Let us look more closely at the entries of $A$. We can write the entry $a_{i,j} = \sum_{\ell = 1}^{k} u_{i,\ell} \cdot \sigma_\ell \cdot v_{j,\ell}$. We can consider each $\ell = 1, \ldots, k$ to be a “concept”. Then, we can interpret $u_{i,\ell}$ as “how much is page $i$ about concept $\ell$”, and $v_{j,\ell}$ as “how much does word $j$ belong in concept $\ell$”, and $\sigma_\ell$ as the “importance” of concept $\ell$. Note that in representing $A$ in this way, we are making the implicit assumption that concepts behave linearly: the frequency with which a word occurs on a page is the sum over all concepts, and there are no superlinear amplifications, or sublinear reductions.

Row $i$ of $U$ can be seen as a $k$-dimensional description of page $i$. Concept $\ell = 1, \ldots, k$ corresponds to a column of $U$ (or a column of $V$), so the $k$-dimensional description of $i$ expresses it in terms of the concepts. Notice that the same column in $U$ and $V$ correspond to the same concept. Also, notice that since $U$ and $V$ are orthonormal, concepts are orthogonal.

Following our intuition, we would expect “similar” pages to have similar concept vectors (again, this is predicated on the assumption that concepts behave linearly). We can then measure the similarity of two pages by comparing their concept vectors. If there is a small angle between the vectors (or a large inner product), then the two pages are about the same (or very similar) concepts. However, up to this point, we are still simply using all the data in the matrix $A$: merely rephrasing it in terms of the SVD does not yet give improved search results.

The important issue which we wanted to address was that the entries of the term-document matrix were not derived from an ideal generation process materializing the concepts in the form of terms. Rather, real matrices have “errors”. More formally, if the world were about $k \ll \min(m,n)$ concepts, then an “ideal world” matrix would have rank $k$. In the real world, web pages do not conform to our ideal concepts. People write web pages, and their individual tendencies and writing styles vary. This will cause $A$ to have rank (almost) $\min(n,m)$. The concepts $\ell = k+1, \ldots, \min(n,m)$ are “error-correcting” concepts explaining peculiarities of $A$. Hence, to get at the “true” contents of pages, we would like to prune out the entries derived from those concepts.

How do we do this? First, we determine the “right” $k$. Having computed the SVD $A = U \cdot \Sigma \cdot V^T$, we “blank out” all concepts for $\ell > k$, by defining

$$
\Sigma_k = \begin{bmatrix} 
\sigma_1 & 0 & \cdots & 0 \\
0 & \ddots & 0 \\
0 & 0 & \sigma_k & 0 \\
\vdots & 0 & 0 & \ddots \\
0 & \cdots & 0 & 0 
\end{bmatrix}
$$

and $A_k = U \cdot \Sigma_k \cdot V^T$. In effect, that gets rid of all concepts for $\ell > k$, and thus only retains $U_k = [u_1 \ u_2 \ \cdots \ u_k]$ and $V_k = [v_1 \ v_2 \ \cdots \ v_k]$. (The $\mathbf{u}_i$ and $\mathbf{v}_i$ are the columns of $U$ and $V$, respectively.) The resulting matrix $A_k$ is still an $m \times n$ matrix, but it now has rank $k$. The following lemma (see, e.g., [185]) shows that it approximates $A$ well.

**Lemma 2.6** The matrix $A_k$ is the best rank-$k$ approximation to $A$: it minimizes $\|A - B\|_2$ over matrices $B$ of rank $k$, where $\|A - B\|_2 = \max_{\|x\|_2 = 1} \| (A - B) \cdot x \|_2$.

Given a document-term matrix $A$ and a query $q$, we would like to find the pages in $A$ that are most relevant to $q$. First, we compute $A_k$ (for some choice of $k$). We consider the space $\mathbb{R}^k$ as the “concept space”, into which we can map both words and pages based on their rows in $U$ and $V$. Specifically, each page $i$ is identified with the $i^{th}$ row of $U_k$, and each word $j$ with the $j^{th}$ row of $V_k$. To find pages relevant to the query, we simply output the pages closest to $q$ in concept space.
How do we choose a good $k$? A good $k$ is one with $\sigma_k \gg \sigma_{k+1}$. If no such $k$ exists, then $A$ cannot be approximated well with low rank, although we will still want to choose some $k$ — it just does not come with good guarantees. In the end, choosing the right $k$ is somewhat of a heuristic.

The result of this approach is that we can (hopefully) extract meaningful concepts, and thus identify two pages as similar, and relevant to a query, even if the terms occurring in them are quite different. Conversely, even if a search term occurs in a document, we may be able to filter it out as irrelevant if other terms force it to lie far away in concept space from the point at which our query is located.

Similar ideas can be applied to several other scenarios:

- Collaborative Filtering (or Recommendation Systems), i.e., the reconstruction of the missing data items. Finding similarities between users (again, in some “concept space”) lets us predict which unseen items a user might like.

- Web site ranking.

- Shopping Predictions.

Singular-value decomposition has found similar applications in a number of domains. In different fields, it (or close variants) is often also called Factor Analysis, Latent Semantic Analysis, or Principal Component Analysis.

### 2.6 Further Reading

Two surveys by Bianchini et al. [44] and by Langville and Meyer [232] provide in-depth discussions of many of the mathematical, algorithmic, and technological challenges involved in link-based search engines. The survey by Langville and Meyer has meanwhile been expanded into a comprehensive book [233].

Some of the technical infrastructure underlying successful search engines is described in [39, 157]. The focus in this chapter, on the other hand, is on link-based inferences about relevance and content, and more specifically spectral methods in this context.

A good, though somewhat outdated, overview of spectral techniques in link-based web search is given by Borodin et al. [53, 54]. Among the many suggested changes or improvements in HITS and PageRank are the following: SALSA [237] performs a random walk on the co-citation graph, but normalizes probabilities. Thus, high-degree hubs divide their hub weight among the authorities they link to, rather than giving full hub weight to each.

An idea similar to PageRank was around the same time proposed by Marchiori [247], building on previous work by Frisse [152] for singly authored hypertext corpora. The relevance of $i$ is its own content, plus a weighted sum of the contents of pages reachable from $i$. The weights decrease exponentially, and if page $i$ has multiple outgoing edges, the assumption is that they will be visited in the order maximizing the relevance of $i$. A similar idea is used by Boyan et al. [55], who also aggregate relevance scores of pointed-to pages with geometrically decreasing weights. In addition, [55] proposes taking into account user feedback to train better retrieval functions and learn the weights for different features. Using clickthrough data (information on which pages were or were not clicked by the searcher) was pursued much more extensively subsequently by Joachims et al., for example [196, 294].

The idea of using eigenvalues as measures of “importance”, “authority”, or “status” long precedes HITS and PageRank. Katz [204] proposes, and Bonacich [52] elaborates on, an approach similar to PageRank to measure social status in a social network: the status of $v$ is the number of paths from other nodes to $v$, with path weights exponentially decreasing in the length of the path. This is similar to giving a node a weight derived from the (discounted) sum of weights of its neighbors. Hubbell [186] uses an essentially mathematically equivalent idea to derive a matrix measuring how strongly pairs of nodes are connected. This matrix is then used to identify “cliques” as densely connected subsets of nodes (differing from the mathematical definition of a clique). Pinski and Narin [289] use a technique very similar to PageRank to measure the importance or influence of scientific publications.
Finding authoritative pages on particular topics is also called topic distillation. Many heuristics have been proposed to avoid topic drift (which happens when a majority of the pages in the focused subgraph are not exactly matching the query) and undue influence of individual nodes. In addition to the two papers discussed above [41, 73], many others propose different heuristics. For instance, Chakrabarti et al. [75] suggest analyzing pages at a more fine-grained level of individual document objects of the page.

In Sections 2.2 and 2.3, as part of our analysis of HITS and PageRank, we also discussed convergence speeds, and remarked that it depends crucially on the spectral gap, the difference between the two largest eigenvalues. By introducing the random jump with probability $\epsilon$, PageRank ensures that the spectral gap is small, and the Markov Chain converges quickly. An experimental study of convergence speeds of PageRank computations is performed by Arasu et al. [16]. They find that PageRank converges faster on the graph of hosts than of pages, and that the Gauss-Seidel method for computing eigenvectors converges faster than the Power Iteration. The convergence speed of PageRank-like computations has received a lot of attention in the numerical analysis community. However, as observed by McSherry [250, 249], performing computations on the right block structure of matrices and reusing values appropriately are sufficient to run a full PageRank computation on an ordinary laptop. At this point, the technological challenges of computing PageRank values (offline) can therefore be considered solved.

A large spectral gap has a second advantage beyond fast convergence: the principal eigenvector does not change much when the matrix is perturbed slightly, e.g., when links are added or deleted, or when rounding errors occur during the computation. In this sense, PageRank is more stable than HITS, a fact pointed out by Ng et al. [280]. Based on this observation, in subsequent work [281], they suggest augmenting HITS by a random jump with probability $\epsilon$ as well. Alternatively, one can consider the top $k$ eigenvectors, up to a point where the eigenvalues drop significantly in magnitude. Such an eigenspace would be more stable. For an in-depth discussion of the effects of perturbations on matrix computations, see [165, 317].

In order to analyze formally the empirically observed performance of LSI and related methods, one needs to describe a model of the underlying “ground truth”, which is usually a randomized model for generating page content. In this context, Papadimitriou et al. [284] were the first to propose a model under which LSI provably performs well: basically, their model posited that different topics define distributions over terms, and pages are drawn from a distribution of topics. If there were few topics with sufficiently distinct signatures, LSI can recover the latent linear structure.

Similar approaches were applied in the context of web search, where generational models also include the links. For instance, a model of Cohn and Chang [94], based on a similar idea of Hoffman [183] for term-document scenarios, assumes that links are generated probabilistically with topic-dependent probabilities. From the sites that a page links to, a maximum likelihood estimation then tries to infer the actual topics. Cohn and Hofmann [95] extend this model to include both terms within the page and links. Achlioptas et al. [3] take this approach yet one step further: for each page, an (unknown) low-dimensional vector describes to what extent the page is a hub and an authority on each of $k$ topics. Text on a page is generated according to distributions parametrized by the topics, and linking probability depend on the authority weights of the target and the hub weights of the source of the link. The algorithm then uses an appropriate SVD to infer the low-dimensional vectors for pages.

The linearity constraint implicit in the use of SVD techniques has been recognized by the community. As a result, standard PCA has been extended to Kernel PCA (see, e.g., [309]). The idea is that the standard inner product (by which one implicitly assumes entries of the matrix to be computed) can be replaced by different notions of inner products. The Kernel of Kernel PCA techniques provides such a notion. Using these techniques requires knowing what inner product is appropriate for the problem at hand.

Implicit in essentially all of the work described in this chapter is the assumption that pages which are “close” in terms of number of hops on the web should be for the most part “similar” in terms of content or topic. Davison [106] and Menczer [252] explicitly test this hypothesis and mostly find it to be true.

To ensure the assumption that links express some notion of endorsement, it is important to deal with nepotistic links and link spamming. Nepotistic links include navigational aids as well as link spam. Several papers attempt to automatically identify such links. Davison [105] uses machine learning techniques to identify decision tree rules for recognizing nepotistic links from examples. The work of Bharat and Henzinger [41] also implicitly discounts nepotistic links. Zhang et al. [346] argue that the PageRanks of nodes with
large numbers of nepotistic links are particularly sensitive to the reset parameter $\epsilon$ in PageRank, and can heuristically be identified by varying $\epsilon$.

Despite the intuitive mathematical justification of the search algorithms presented here, the ultimate judgment is the relevance of the results returned. Naturally, in this context, the success of search engines like Google or Yahoo! constitutes some amount of verification. A principled experiment was attempted by Amento et al. [14], who asked 40 human users for the most authoritative documents, and then compared with the results of search algorithms, as well as various features. They found that link-based analysis yields good results overall, but surprisingly, merely the number of pages co-hosted with the given page is a very good predictor as well.
Chapter 3

Community Structure and Clusters in Graphs

In this chapter, we turn to the problem of identifying communities from link information. A community is defined informally as a cohesive set of nodes, in the sense that the node set has “many” links within, and “few” links connecting it to other nodes. Of course, “many” and “few” are still not very clearly defined terms, and we will see below how different instantiations of these concepts will lead to different objectives.

But first, let us turn to the question of why we are interested in identifying community structure in networks. The most important reason is that community structure may reveal useful information about the network. In a network of friendships, we may identify a large group of tightly knit friends. In a terrorist network, with edges representing phone calls or other interactions, we may identify a cell. In a graph such as the WWW, we might discover communities of pages pertaining to the same specialized subject. (Here, we would be using the implicit assumption of homophily: that pages on similar topics are more likely to link to each other.) For another example, many biologists believe that if a set of proteins contains many pairwise interactions, then these proteins are often involved in some specific common function. In summary, the most important benefit of identifying communities based on links is the functional or content predictions one can obtain from the structure.

Two side benefits may not be as readily apparent. Identifying high-level community structure may permit studying a network at a more appropriate level of detail. This could include replacing communities with one or few meta-nodes, and studying the network at the level of those meta-nodes. Or one could zoom in on one particular community, disregarding other parts of the network. For instance, Dill et al. [117] found empirically that the WWW graph exhibits a lot of self-similarity: within communities defined based on common topics or other sufficiently homophilous criteria, degree distributions and other statistics closely resembled those of the entire web graph. A fairly detailed discussion of these and other motivations can be found, for instance, in Newman’s paper [275].

A further application of community identification is the following: often, social scientists are particularly interested in links that cross between communities. Granovetter [169] calls such links weak ties. The particular importance of weak ties is that they provide individuals with much more “fresh” information than strong ties to their good friends, presumably because the close friends share the same social circles, and therefore do not have as much novel information. Granovetter [169] reports on data showing that weak ties were much more instrumental than strong ties in helping individuals find jobs. Further corroborating the importance of weak ties in connecting social networks was an experiment analyzing a high-school friendship network. The graph representing only the strong ties, where each person is linked only to their two best friends, contains many small, isolated communities. On the other hand, including links to the top eight friends (which will include weaker ties) results in a graph with a giant connected component, and only a few smaller clusters.

The importance of links crossing between different communities has meanwhile been reiterated and stressed in many other papers. Most notably, Burt (e.g., [67, 68]) advances a theory of structural holes,
arguing that the existence of distinct communities tends to imply entrepreneurial opportunities for individuals or companies bridging those communities.

In order to pose questions about weak links and structural holes mathematically, and analyze models for their formation or function, it is first necessary to identify communities as such. This would also allow us to formalize the problem of identifying weak links or structural holes.

For the remainder of this chapter, we will use the following notation.

**Definition 3.1** Let $G = (V, E)$ be a graph.

1. $e(S, T) = E \cap (S \times T)$ denotes the set of edges with exactly one endpoint in $S$ and one in $T$. $e(S) = e(S, S)$ is the set of edges with both endpoints in $S$.

2. We use $d_S(v) = |e(\{v\}, S)|$ to denote the degree of $v$ within the set $S$, i.e., number of edges between $v$ and nodes in $S$.

3. The edge density of a node-set $S$ is defined as $\frac{|e(S)|}{|S|}$.

In exploring the spectrum of different definitions of the term “community”, there are several parameters along which we can distinguish definitions:

1. Can nodes be in multiple communities? Does each node have to be in at least one community? If these questions are answered “No” and “Yes”, respectively, then we are seeking a partition of the nodes. As such, the problem of identifying community structure shares a lot of overlap with the problem of clustering, and we will explore several objective functions in Sections 3.4 and 3.5.

2. If the answers are “Yes” and “No” instead, then we are looking to identify individual, possibly overlapping communities. These will be “unusually dense” sets of nodes. Here, an interesting distinction is whether we want to include nodes which have overall high degree, or focus on nodes which have “most” of their edges within a set. This distinction will give rise to different definitions in Sections 3.1 and 3.2.

3. If we follow the reasoning about hubs and authorities of Section 2.2, we would suspect that in the web graph, communities could be identified by dense bipartite cores of hubs and authorities. Heuristic approaches based on these ideas are discussed in Section 3.3.

### 3.1 Dense Subgraphs

Perhaps the simplest and most “obvious” definition of a community is a dense subgraph, i.e., a subgraph with large edge density in the sense of Definition 3.1. If we want to find the best community, that would be the subgraph with largest density.

**Problem 1** *In an arbitrary graph $G$, find the densest subgraph $S$, i.e., the set $S$ maximizing $\frac{|e(S)|}{|S|}$.*

We begin with the decision version of the problem: Given a constant $\alpha$, is there a subgraph $S$ with $\frac{|e(S)|}{|S|} \geq \alpha$? This constraint can be re-written as:

$$|e(S)| - \alpha|S| \geq 0. \quad (3.1)$$

Each internal edge of $S$ contributes 2 to the total degree in $S$, and each edge leaving $S$ contributes 1, so the number of internal edges $|e(S)|$ in $S$ is $|e(S)| = \frac{1}{2}(\sum_{v \in S} d(v) - |e(S, S)|)$. Substituting this in Equation (3.1) gives us

$$\sum_{v \in S} d(v) - |e(S, S)| - 2\alpha|S| \geq 0, \quad (3.2)$$
or, equivalently,

\[
\sum_{v \in V} d(v) - \left( \sum_{v \in \mathcal{S}} d(v) + |e(S, \mathcal{S})| + 2\alpha |\mathcal{S}| \right) \geq 0. \tag{3.3}
\]

As \(2|E|\) is a constant (independent of \(S\),) this constraint is satisfiable if it is satisfied by the set \(S\) with minimum \(\beta(S)\). To find such a set, we formulate a mincut problem as follows. Consider the graph \(G'\) with \(V' = V \cup \{s, t\}\), where \(s\) is connected to all vertices \(v\) with an edge of capacity \(d(v)\), and \(t\) is connected to all vertices in \(V\) with an edge of capacity \(2\alpha\). The cost of the cut \((S + s, \mathcal{S} + t)\) in \(G'\) is exactly \(\beta(S)\). Thus, if the minimum \(s\)-\(t\) cut \((S + s, \mathcal{S} + t)\) of \(G'\) satisfies the constraint in Equation (3.3), then \(S\) is a set of density at least \(\alpha\). If not, then no such set exists. The minimum \(s\)-\(t\) cut can of course be computed with any of the standard Mincut algorithms (see, e.g., [222]).

Of course, our real goal was to find the maximum value of \(\alpha\) for which a set \(S\) of density \(\alpha\) exists. We could accomplish this with a binary search over values of \(\alpha\). A more efficient way uses the parametric max-flow algorithm of Gallo, Grigoriadis and Tarjan [155] to compute minimum \(s\)-\(t\) cuts for all values of \(\alpha\) in one computation. Since this is a fairly general and useful technique, we state the main result here as a theorem:

**Theorem 3.2 (Parametric Maximum Flow [155])** Let \(\alpha\) be a real-valued parameter, and \(G = (V, E)\) a graph with non-negative edge capacities \(c_e\) which can depend on \(\alpha\) in the following way:

1. Capacities of edges out of \(s\) are non-decreasing in \(\alpha\).
2. Capacities of edges into \(t\) are non-increasing in \(\alpha\).
3. Capacities of edges not incident with \(s\) or \(t\) are constant.

Then, the maximum flow and minimum cut for all values of \(\alpha\) can be computed simultaneously with one invocation of a modified Goldberg-Tarjan algorithm [161], in time \(O(n^2m)\). Furthermore, the minimum cuts \((S_\alpha, \overline{S}_\alpha)\) are nested, in the sense that \(S_\alpha \subseteq \overline{S}_{\alpha'}\) whenever \(\alpha \leq \alpha'\).

**Remark 3.3** While we gave an algorithm for the problem of finding the overall densest subgraph, the problem becomes significantly more difficult once we restrict the size of the subgraph. For instance, finding the densest subgraph of exactly (or at most) \(k\) nodes is an NP-hard problem, as can be seen easily by reducing from the \(k\)-clique problem, setting the density to be \(k - 1\).

The densest \(k\)-subgraph problem has been studied in several papers [20, 43, 139, 142]. The current best approximation ratio is \(O(n^{1/4+\epsilon})\) (in time \(O(n^{1/4})\) for any \(\epsilon > 0\) [43].

Only recently did Khot [214] rule out the existence of a PTAS for the densest \(k\)-subgraph problem: unless the Unique Games Conjecture [213] is violated, there is some \(\alpha > 1\) such that no \(\alpha\)-approximation to the densest subgraph can be found in polynomial time. Obviously, the upper and lower bounds are far from matching, and closing this gap is an interesting open question.

A straightforward, but quite useful, generalization of the Densest Subgraph problem specifies a set \(X\) of vertices that must be included in the set \(S\).

**Problem 2** Given a graph \(G\), find the densest subgraph \(S\) containing a specific vertex set \(X\) (i.e., \(S\) maximizes \(|\mathcal{S}| \subseteq |\mathcal{S}|\) over all sets \(S \supseteq X\)).

This allows us to find out what the nodes in \(X\) “have in common”, by identifying a dense subgraph \(S\) containing them, and then inspecting the nodes in the dense subgraph for their attributes.

This generalization can be solved quite easily, simply by giving all edges from \(s\) to vertices \(v \in X\) a capacity of \(\infty\). This will ensure that all nodes in \(X\) are on the \(s\)-side of the cut, and the rest of the analysis stays the same.
3.1.1 A 1/2-Approximation

In some real-world graphs, such as the WWW or the graph of all friendships among people in the US, the running time of $O(mn^2)$ from Theorem 3.2 is still prohibitively large. Thus, we are interested in faster algorithms, preferably running in linear or near-linear time. As it is not known how to compute minimum $s$-$t$ cuts much faster, we will have to settle for an approximation algorithm. The following greedy algorithm was first analyzed by Charikar [78]. We present it immediately for the generalized version (Problem 2) of the Densest Subgraph Problem.

Algorithm 1 A Greedy $\frac{1}{2}$-Approximation Algorithm for finding dense subgraphs

Let $G_n \leftarrow G$.

for $k = n$ downto $|X| + 1$ do

Let $v \notin X$ be the lowest degree node in $G_k \setminus X$.

Let $G_{k-1} \leftarrow G_k \setminus \{v\}$.

Output the densest subgraph among $G_n, \ldots, G_{|X|}$.

Claim 3.4 Algorithm 1 is a $\frac{1}{2}$-approximation.

Proof. Let $S \supseteq X$ be the densest subgraph. If our algorithm outputs $S$, then it is clearly optimal. If not, then at some point, we must have deleted a node $v \in S$. Let $G_k$ be the graph right before the first $v \in S$ was removed. Because $S$ is optimal, removing $v$ from it would only make it worse, so

$$|e(S)| \geq |e(S-v)| \geq |e(S)| - d_S(v) \geq |e(S)| - |S| + 1.$$

Multiplying through with $|S|(|S| - 1)$ and rearranging gives us $d(v) \geq \frac{|e(S)|}{|S|}$.

Because $G_k$ is a supergraph of $S$, the degree of $v$ in $G_k$ must be at least as large as in $S$, so $d_{G_k}(v) \geq d_S(v) \geq \frac{|e(S)|}{|S|}$. The algorithm chose $v$ because it had minimum degree, so we know that for each $u \in G_k \setminus X$, we have $d_{G_k}(u) \geq d_{G_k}(v) \geq \frac{|e(S)|}{|S|}$. We thus obtain the following bound on the density of the graph $G_k$:

$$\frac{|e(G_k)|}{|G_k|} \geq \frac{\sum_{u \in S} d_S(u) + \sum_{u \in G_k \setminus S} \frac{|e(S)|}{|S|}}{2|G_k|} \geq \frac{|e(S)| |S| + |G_k \setminus S|}{2|G_k|} \geq \frac{|e(S)| |S| + |G_k \setminus S|}{2|G_k|} \frac{1}{2|S|}.$$

The graph that the algorithm outputs is certainly no worse than $G_k$, as $G_k$ was available as a potential solution. Hence, the algorithm is a $\frac{1}{2}$-approximation.

3.2 Strong Communities

In the previous section, we defined a community as a dense subgraph. That means that the group of nodes overall has many edges. One feature of this definition is that it will “usually” favor the inclusion of nodes with overall high degree. For instance, in terms of communities in the WWW, we would expect that for any sufficiently large set $S$ of nodes, the most commonly linked to other nodes would be yahoo.com, google.com, or cnn.com, even if the nodes in $S$ do share a much more narrow feature. In terms of our introductory
discussion, we were requiring many edges within the community, but not few edges connecting it to other
nodes. For some applications, one could argue that communities should also be well separated from other
nodes, in the sense that each node has most of its links inside the set \( S \). This leads us to the following
definition, slightly generalizing one by Flake et al. [144].

**Definition 3.5** Let \( G \) be a graph, and \( \alpha \in [0,1] \). A set \( S \subseteq G \) is called a strong \( \alpha \)-community in \( G \) iff
\[
d_S(v) \geq \alpha d(v) \text{ for all nodes } v \in S.
\]

This definition captures that each individual should “belong” to the community. The value \( \alpha = \frac{1}{2} \) has
a particularly natural interpretation, in that it requires that each node in \( S \) have at least as many edges inside
\( S \) as to nodes outside \( S \). It is also the case considered in [144]. If we omit the value of \( \alpha \), we assume \( \alpha = \frac{1}{2} \).
The larger \( \alpha \), the more “tightly knit” the community is. So the best communities are the ones with large
\( \alpha \). However, this leads to the problem that the whole graph is the “best” community, because all of its edges
are within the selected set (itself), making it a 1-community.

The entire graph is not a very interesting community. We are more interested in discovering significantly
smaller communities which nevertheless have many edges inside. As an extension, we may also wish to force
the inclusion of one or more nodes, and find the best community containing them.

Unfortunately, finding the smallest (or even approximately smallest to within a factor of \( c \log(n) \) for some
constant \( c \)) community is NP-hard, with or without the inclusion of specific sets.

**Theorem 3.6** 1. It is NP-complete to decide whether a given graph \( G \) has an \( \alpha \)-community of size at
most \( k \).

2. Unless \( P=NP \), the size of the smallest \( \alpha \)-community cannot be approximated within a factor of \( O(c \log n) \)
in polynomial time, for some \( c \).

3. The same results hold if we require to find a community of a specified vertex \( v \).

**Proof.** Membership in NP is obvious. The NP-hardness will be implied by the approximation hardness,
so we will only prove the second part of the theorem (and show the few simple extensions necessary to prove
the third). To avoid notational clutter, we assume here that \( \alpha = \frac{1}{2} \).

We reduce from the **Set Cover** problem. Let \( X = \{x_1, \ldots, x_n\} \) be the set of elements, and \( S_1, \ldots, S_m \subseteq X \) the sets. The goal is to find a minimum number of sets \( S_i \) whose union covers \( X \). To do so, we construct
an undirected and unweighted graph \( G \). Let \( x_i \) for an element \( x_i \) denote the set of all indices \( j \) such that
\( x_i \in S_j \).

We start with a large clique \( B \) of \( 4mn^2 + 1 \) bulk nodes. These will be used solely to increase the degree
of other nodes and force them to include many of their more “legitimate” neighbors — we don’t even need
to give them individual names.

We also have a set \( F \) of \( n \) forcer nodes, each of which is connected to \( n \) distinct bulk nodes. Forcer nodes
are not connected among each other.

The other nodes are more important, and fall into two classes: first, there is an element node \( v_i \) for each
element \( x_i \), and the \( v_i \) form a clique. In addition, each \( v_i \) is connected to \( 2n + 1 - |X_i| \) distinct bulk nodes,
and to each of the \( n \) forcer nodes.

For each set \( S_j \), there are \( 2n + 1 \) set nodes \( s_{j,a} \), which are connected to form a clique (for fixed \( j \)). Each
\( s_{j,a} \) for \( 2 \leq a \leq 2n + 1 \) is connected to \( 2n \) distinct bulk nodes, whereas \( s_{j,1} \) is connected to \( 2n + |S_j| \) distinct
bulk nodes. Finally, there is an edge between \( s_{j,1} \) and \( v_i \) whenever \( x_i \in S_j \) in the **Set Cover** instance.

In the resulting graph, each element node \( v_i \) has degree exactly \( 4n \), each set node \( s_{j,a} \) for \( a \geq 2 \) has degree
\( 4n \), the nodes \( s_{j,1} \) have degree \( 4n + 2|S_j| \), forcer nodes have degree \( 2n \), and the bulk nodes have degree at
least \( 4mn^2 \).

Given a set cover of size \( k \), we can obtain a community as follows: Let \( C \) consist of all element nodes \( v_i \),
all forcer nodes, and all set nodes \( s_{j,a} \) for which \( S_j \) is in the set cover. The size of \( C \) is \( k(2n + 1) + 2n \).
To verify that it is a community, notice that in the induced subgraph of \( C \), each node \( s_{j,a} \) for \( a \geq 2 \) has degree
\( 2n \), and node \( s_{j,1} \) has degree \( 2n + |S_j| \). Each node \( v_i \) has degree at least \( 2n \), because \( C \) contains all other
\( v'_j \), all forcer nodes, and a node \( s_{j,1} \) for some set \( S_j \) containing \( x_i \).
Conversely, consider a community $C$ of size at most $k(2n+1)+2n$ for some $k < n$ (solutions with $k \geq n$ are of course uninteresting for set cover approximations). Because it is quite small, $C$ cannot contain any bulk nodes — if it contained any one of them, it would have to contain at least $2mn^2 > m(2n+1)+2n$ of them to be a community.

Now, if $C$ contains any set node $s_{j,a}$, it must contain all of the $s_{j,a'}$ for that particular $j$ in order to be a community (otherwise, the degree requirement for $s_{j,a}$ would be violated, as none of its adjacent bulk nodes are included). By the size constraint, $C$ contains the set node clique for at most $k$ sets $S_j$, and we will prove that these sets $S_j$ form a set cover.

In order to satisfy the degree requirement, all of the $v_i$ adjacent to $s_{j,1} \in C$ must be included — in particular, any community $C$ (of size at most $k(2n+1)+2n$) contains at least one element node $v_i$.

In order for the element node to have its required degree of $2n$ without bulk nodes, $C$ must contain a forcer node, for $C$ can contain at most $n$ element nodes and $k \leq n-1$ nodes $s_{j,1}$. But in a community $C$ with a forcer node and without bulk nodes, all element nodes $v_i$ must be included to meet the forcer node’s degree requirement. Thus, any community $C$ without bulk nodes contains all element nodes.

Now, to meet an element node $v_i$’s degree requirement of $2n$ without including bulk nodes, at least one set node adjacent to $v_i$ must be included in $C$, for all forcer and element nodes would only yield degree $2n-1$. Hence, the sets $S_j$ for which $s_{j,1} \in C$ indeed form a set cover of size at most $k$.

This completes the proof that the above reduction is approximation preserving, thus showing the inapproximability of the community problem, because Set Cover is known to be inapproximable within $O(n \log n)$ for some $c$ (unless $P=\text{NP}$) [138, 297].

When we ask about communities including a particular node $v$, we can use exactly the same reduction, and fix $v$ to be a forcer node — the result is going to be the same, since forcer nodes were just shown to be included in all sufficiently small communities.

Given that even approximating the smallest $\alpha$-community within $O(n \log n)$ is NP-hard, we will look at heuristics: approaches that may often work in practice, even if they give us no guarantees. Let $(S, \overline{S})$ be a minimum $s$-$t$ cut in $G$. Then, $S$ is almost a $\frac{1}{2}$-community because each $v \in S \setminus \{s,t\}$ has at least as many edges inside $S$ as crossing $(S, \overline{S})$: otherwise moving $v$ to the other side of the cut would make the cut cheaper. If this also held for $s$ and $t$, then $(S, \overline{S})$ would be a $\frac{1}{2}$-community.

If we are looking for communities including a given node $s$, we can use the above heuristic approach to compute the minimum $s$-$t$ cut for all $t$ in $n$ min-cut computations. Then, we simply take the best cut found this way.

If we are looking for just communities, without specifying a node $s$, then we can try all $(s, t)$ pairs ($\Theta(n^2)$ min-cut computations). We can reduce that number of computations to $O(n)$ using Gomory/Hu trees [167]. The idea is that all min-cuts can be “encoded” in a tree.

**Theorem 3.7 (Gomory-Hu Trees [167])** Let $G = (V, E)$ be a graph. For all node pairs $i, j \in V$, let $f_{ij}$ be the maximum flow (min-cut) between $i$ and $j$. Let $G'$ be the complete graph on $V$ with edge costs $f_{ij}$. Let $T'$ be a maximum spanning tree of $G'$.

For each $i, j \in V$, let $P_{ij}$ denote the (unique) $i$-$j$ path in $T$. Then, the tree $T$ has the property that $f_{ij} = \min_{e \in P_{ij}} f_e$ for all $i, j$. If $e$ is the edge attaining the minimum, then the two connected components of $T \setminus \{e\}$ define a minimum $i$-$j$ cut in the original graph $G$.

**Proof.** We will prove that $f_{ij} = \min_{e \in P_{ij}} f_e$ by contradiction, ruling out inequality in both directions.

If $f_{ij} > \min_{e \in P_{ij}} f_e$, then inserting $(i, j)$ into $T$ and removing the edge $e \in P_{ij}$ minimizing $f_e$ would create a more expensive tree $T'$. This contradicts the assumption that $T$ was a maximum spanning tree of $G'$. So $f_{ij} \leq \min_{e \in P_{ij}} f_e$.

For the other direction, let $(S, \overline{S})$ be an $i$-$j$ cut of capacity $f_{ij}$. Because $P_{ij}$ is an $i$-$j$ path, it must cross this cut, i.e., there is an edge $e = (u, v) \in P_{ij}$ with $u \in S, v \in \overline{S}$. So $(S, \overline{S})$ is also a $u$-$v$ cut, and thus, $f_{uv} \leq f_{ij}$. But then, $\min_{e \in P_{ij}} f_e \leq f_{uv} \leq f_{ij}$, completing the proof.

Therefore, the $s$-$t$ cuts for all $s, t \in V$ can be compactly represented. An interesting additional consequence is that there are only $n-1$ different min-cuts (and associated min-cut values) for the $n(n-1)$
different pairs of nodes. Gomory and Hu [167] also show how to compute $T$ from $G$ using only $n - 1$ min-cut computations. Using the approach of first computing $T$, we can thus find the communities for all nodes with only $O(n)$ min-cut computations.

### 3.2.1 Other values of $\alpha$

If we want to find $\alpha$-communities for $\alpha \neq \frac{1}{2}$, we need to adapt the approach. The idea is to adapt some degrees so that $d_S(v) \geq \alpha d_G(v)$ holds in the original graph if and only if $d'_S(v) \geq \frac{1}{2}d'_G(v)$ holds in the new graph. We do this by adding a source and sink, and connecting them to each node $v$ in an effort to “balance” the sides of the inequality.

If $\alpha < \frac{1}{2}$, each node $v$ has an edge to the source with capacity $(1 - 2\alpha)d(v)$; otherwise, each node $v$ is connected to the sink with capacity $(2\alpha - 1)d(v)$. Nodes $x$ whose inclusion (or exclusion) is required are connected to $s$ (resp., $t$) with infinite capacity. (Notice that if no nodes are connected with infinite capacity, then the minimum $s$-$t$ cut will always just separate $s$ or $t$ from the rest of the graph, corresponding to our above intuition that the entire graph is the best community.) The algorithm then just looks for the minimum $s$-$t$ cut in the resulting graph $G'$.

**Claim 3.8** This approach will produce “almost $\alpha$-communities,” in the sense that all nodes except those whose inclusion was forced will meet the community constraint.

**Proof.** Here, we give the proof for the case that $\alpha \geq \frac{1}{2}$ — the proof for $\alpha < \frac{1}{2}$ is symmetric.

As in the argument for plain min-cuts in a graph, the fact that a node $u$ (without infinite-capacity edges) is on the $s$-side as opposed to the $t$-side implies that $d_S(u) \geq d_G(u) + (2\alpha - 1)d(u)$, or $d_S(u) + d(u) - d_G(u) \geq 2\alpha d(u)$. Now, because $d(u) = d_S(u) + d_T(u)$, this becomes $2d_S(u) \geq 2\alpha d(u)$, or $\frac{d_S(u)}{d(u)} \geq \alpha$.

The sets found by this approach can violate the community constraint at the nodes whose inclusion was forced. Notice that this may happen even when there are communities including/excluding specific nodes, i.e., the fact that the algorithm did not find a community does not mean that none exists.

Flake at al. [144] propose another min-cut based heuristic, which makes more of an effort to avoid the “entire graph” problem described above. For a parameter $\delta$ that will be varied, all nodes other than a specified node $s$ to be included are connected to a new sink $t$ with capacity $\delta$. We then vary the parameter $\delta$, and look for the minimum $s$-$t$ cut. The cuts found this way can then be inspected manually, and interesting ones extracted. (Notice that they can be found with one parametric max-flow computation [155], as described in Theorem 3.2.)

For $\delta = 0$, the minimum cut is $(V, \{t\})$. On the other hand, for very large $\delta$, the minimum cut is $((s) \cup (t) \setminus \{s\})$. If along the way, some $\delta$ yields a non-trivial solution, that is an “almost community,” in that it violates the constraint only at the node $s$.

### 3.3 Bipartite Subgraphs

The previous views of communities, both of which essentially looked for dense subgraphs, were based on the implicit assumption that nodes within a community are “likely” to link to each other. This homophily often applies in social settings. However, as we argued in Section 2.2, it is a more questionable assumption in competitive settings such as the WWW. For instance, within the “community” of car manufacturers, we would expect relatively few or no links between the most prominent members. At the time, the solution was to look (implicitly) for links in the co-citation graph, or look for bipartite graph structures of hubs and authorities.

Kumar et al. [229] propose the same approach for community identification, and argue that large and dense bipartite graphs are the “signature” of communities in the WWW.

Ideally, we would like to enumerate “all” such signatures, and expand them to communities. However, the complexity of doing so would be prohibitive. If nothing else, deciding the presence of large complete bipartite graphs is NP-hard, and as hard to approximate as the CLIQUE problem itself. However, in this
context, it is interesting to note that if the given graph is dense enough (i.e., contains enough edges), it always has a large complete bipartite subgraph.

**Lemma 3.9** If a bipartite graph has $\Omega(b^{1/3}n^{5/3})$ edges, it contains a $K_{3,b}$, i.e., a complete bipartite subgraph with $3$ nodes on one side, and $b$ nodes on the other.

**Proof.** For each node $v$ on the right side, we write $\delta(v)$ for the set of its neighbors, and $d(v) = |\delta(v)|$ for its degree. Each node $v$ is labeled with each 3-element subset $T$ of $\delta(v)$ (i.e., with all $T \subseteq \delta(v), |T| = 3$). Notice that each node thus has many labels, namely $\left(\frac{d(v)}{3}\right)$. Taken over all nodes on the right side, the total number of labels is $\sum_v \left(\frac{d(v)}{3}\right)$. We assumed in the statement of the lemma that $\sum_v d(v) = \Omega(b^{1/3}n^{5/3})$. The total number of labels is minimized when all $d(v)$ are equal, i.e., $d(v) = \Omega(b^{1/3}n^{2/3})$. Even then, the number of labels is

$$\sum_v \left(\frac{d(v)}{3}\right) = n\left(\frac{\Omega(b^{1/3}n^{2/3})}{3}\right) = n\Omega(bn^2) = \Omega(bn^3).$$

But the total number of distinct labels is only $\binom{n}{3} = O(n^3)$. Hence, by the Pigeonhole Principle, some label must appear at least $b$ times. The $b$ nodes on the right side sharing the label, and the three nodes on the left side who form the parts of the label, together form a $K_{3,b}$. 

By considering $a$-tuples instead of triples for labels, we can obtain the generalization that any bipartite graph with $\Omega(b^{1/a}n^{2-1/a})$ edges contains a $K_{a,b}$.

While it is interesting to know that sufficiently dense graphs will contain a $K_{a,b}$, it does not necessarily help us in finding one, in particular if the graph is not dense. For large $a$ and $b$, the problem is NP-hard in general, but we may still be interested in speeding up the search for smaller, and practically important, values, such as finding $K_{3,6}$ graphs. By brute force (trying all 9-tuples of nodes), this would take $\Theta(n^9)$ steps. A first and simple improvement is given by realizing that we only need to look at triples of nodes on one side. Given nodes $v_1, v_2, v_3$, we can take the intersection of their neighborhoods $\bigcap_i \delta(v_i)$. If the intersection contains at least $b$ elements, then a $K_{3,b}$ has been found, else those three nodes cannot be part of a $K_{3,b}$. This reduces the running time to $\Theta(n^4)$.

The ideas underlying this improvement can be extended further. Obviously, any node of indegree less than 3 can be pruned, and similarly for outdegrees less than $b$. Once nodes have been pruned, we can iterate, as the degree of other nodes may have been reduced. In addition, if a node reaches indegree exactly 3 (or outdegree exactly $b$), it can be verified easily if it and all its neighbors form a $K_{3,b}$, after which they can either be reported (and pruned), or just pruned. These heuristics, while not useful in a worst-case scenario, help a lot in practice. They were reported, along with other heuristics, by Kumar et al. [229], and used to identify several 10,000 communities in the web graph.

**Remark 3.10** Finding large complete bipartite graphs can be likened to finding dense areas of a 0-1 matrix, a task known as association rule mining in the data mining community [5]. A common approach there (see, e.g., [6]) is to take simple rules, and combine them into larger rules. The idea is that any subgraph of a larger complete (or dense) graph must itself be complete (or dense). Hence, looking only at combinations of small dense graphs rules out a lot of unnecessary attempts. By starting from $K_{1,1}$ graphs, extending them to $K_{1,2}$ and $K_{2,1}$, then to $K_{2,2}$, $K_{3,1}$, and $K_{1,3}$, etc., we make sure to only look at relevant data, which leads to a lot of speedup in practice (though again no theoretical guarantees in the case of dense graphs).

### 3.3.1 Spectral techniques

If we follow the intuition of hubs and authorities a little further, we arrive at an alternate approach to identify community structure. Rather than looking merely for dense bipartite subgraphs (or large complete bipartite graphs), we could instead look at other eigenvectors of the adjacency matrix (or cocitation matrix). Nodes with large positive or negative entries in those eigenvectors (hub or authority weights) could then be considered as the cores of communities [160].
To gain some intuition for why this might make sense, consider the following adjacency (or cocitation) matrix $B$.

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

The largest eigenvalues of $B$ are $\lambda_1 = 3$ and $\lambda_2 = 2$, and the corresponding eigenvectors are $e_1 = (0, 0, 1, 1, 1)$, and $e_2 = (1, 1, 0, 0, 0)$, respectively. The sub-communities found are the bottom three nodes (a $K_{3,3}$), and the top two nodes (a $K_{2,2}$).

Here is the intuitive reason why this works: consider a vector $e$ which is an eigenvector for a large eigenvalue $\lambda$. Let $i$ be some index such that $e_i$ is “large”. The $i^{th}$ row of $B \cdot e$ is $\sum b_{ij} e_j$, and must be equal to $\lambda e_i$ because of the eigenvector property. Because both $\lambda$ and $e_i$ were assumed to be “large”, the sum is large, and most of the contribution to the sum must come from the entries $j$ for which $e_j$ is large. But then, the corresponding $b_{ij}$ must also be large. In other words, the $b_{ij}$ values are “fairly large” among the pairs $(i, j)$ such that the $e_j$ are large. But this just means that the nodes $i$ with large $e_i$ values are relatively densely connected in $B$, which is exactly what we are looking for in a community.

While the above provides some intuition, there are also quite a few things that can be said formally about this approach. One is the following lemma, which is a well-known basic fact about random walks on undirected graphs.

**Lemma 3.11** Let $G$ be an undirected graph, and define the matrix $A$ as $a_{ij} = \frac{1}{d(j)}$ if $(i, j)$ is an edge of $G$, and $a_{ij} = 0$ otherwise. (A corresponds exactly to the transition probabilities for a random walk on $G$.) Then, the largest eigenvalue of $A$ is 1, and its multiplicity $k$ is the number of connected components of $G$. Furthermore, the $k$ (right) eigenvectors corresponding to the eigenvalue 1 identify the connected components of $G$.

**Remark 3.12** We are deliberately vague about “identifying the connected components”. If there is more than one component, then the eigenvectors are not unique: any set of $k$ orthonormal vectors spanning that space is possible. We can say the following, however: there is a set of $k$ orthonormal eigenvectors such that the $i^{th}$ eigenvector has non-zero entries exactly for the nodes in the $i^{th}$ connected component. Furthermore, given any set of $k$ eigenvectors, it is fairly easy to obtain instead a set with this “identifying” property.

Of course, there are easier ways to compute connected components, and furthermore, connected components are not necessarily the most exciting type of communities. However, Lemma 3.11 at least motivates the approach of Gibson et al. [160] described above.

Much more can be said about eigenvectors of different matrices derived from the adjacency matrix of a graph, and the sense in which they identify communities. In particular, eigenvectors are useful in identifying sets with small “surface-to-volume” ratio, i.e., many edges inside the set, and few edges outside. For a detailed introduction to spectral graph theory, see for instance [87, 45, 259]. In particular, the Cheeger Inequality (e.g., [13, 87]) gives a relationship between the expansion of a graph and its eigenvalues; the proof also shows why eigenvectors of the Laplacian of a graph permit finding a low-expansion node set.

### 3.4 Modularity

We now move away from the problem of identifying just one community, and instead try to partition the graph into a collection of **disjoint** communities. In this sense, the problem of finding communities is synonymous to the problem of partitioning a graph, or of clustering a graph. Naturally, there are many different objectives describing the quality of such a partitioning. Here, we will focus on a measure called **modularity** proposed by Newman et al. [278, 274].
To motivate the idea behind the modularity measure, we return to an observation made earlier in the context of α-communities: high-degree nodes, by definition, tend to contribute more edges to communities. Thus, if we include many high-degree nodes in a community together, we expect to capture a lot of the edges. On the other hand, if we isolate low-degree nodes, we expect not to lose many edges. This really does not tell us much about communities, but only recovers known information about the degree distribution. The idea of the modularity measure is to capture how many more edges a partitioning explains beyond what could be predicted merely from the degree distribution. To capture what we mean by “from the degree distribution”, we look at how many more edges are inside communities than would be in a random graph with the same degree distribution.

Formally, consider a graph \( G = (V, E) \), and a partition \( P = \{S_1, \ldots, S_k\} \) of that graph. If the partition “explains” the communities in the graph, then just the partition, without any information about the edges, should allow us to “reconstruct” the edges relatively well. How well is captured by modularity. The actual number of edges inside \( S_i \) is exactly \(|e(S_i)|\). If all we knew about \( G \) was its degree distribution, then the expected number of such edges can be calculated as follows: we assume that the graph is a uniformly random multi-graph of the same degree distribution, i.e., we allow self-loops and parallel edges.

Consider a vertex \( v \in S_i \) of degree \( d(v) \). For any vertex \( u \), the probability that a given edge \( e \) of \( v \) has \( u \) as its other endpoint is \( \frac{d(u)}{2m} \), where \( m \) is the total number of edges. The reason is that with \( m \) edges, there are \( 2m \) edge endpoints, and for any one edge, it might choose any one of these endpoints (except the ones that’s already used up). By linearity of expectation, the expected number of edges between \( u \) and \( v \neq u \) is therefore \( \frac{d(u)d(v)}{2m-1} \). The expected number of self-loops of \( v \) is \( \frac{d(v)(d(v)-1)}{2m-1} \), because one of the edge endpoints \( v \) is already used up for the other end of the edge. Summing up over all pairs \( u, v \) in \( S_i \), now gives us that the expected number of edges with both endpoints inside \( S_i \) is \( \frac{1}{2} \frac{d(S_i)(d(S_i)-1)}{2m-1} \), where we write \( d(S_i) = \sum_{v \in S_i} d(v) \). The factor of \( \frac{1}{2} \) here arises because each edge is considered twice, once for each of its endpoints. Ignoring the two \(-1\) terms, and summing up over all of the partitions gives us that the expected number of edges inside partitions is (roughly) \( \frac{1}{4m} \sum_i d(S_i)^2 \).

We define the modularity \( q(P) \) of a partition \( P \) as the difference between the actual number of edges inside clusters, and the expected such number under a uniformly random multigraph, normalized to a scale of \([-1, 1]\) by dividing out the total number of edges. Thus, the definition is

\[
q(P) = \frac{1}{m} \sum_i |e(S_i)| - \frac{1}{4m} d(S_i)^2.
\]  

(3.5)

Thus, we can now formally state the algorithm question as finding a partitioning \( P \) maximizing \( q(P) \). Notice that the number \( k \) of clusters is not specified; avoiding the need to pre-specify the number of clusters is one of the advantages of the modularity objective. Of course, if desired, one can also consider a version in which the number of clusters is actually given.

Little is known about algorithms with provable guarantees for the problem of finding a modularity-maximizing partitioning. Only recently did Brandes et al. prove that the problem is NP-hard [61]. So far, no algorithms with provable guarantees are known. Many heuristics have been proposed, based on bottom-up greedy aggregation [274, 90], identifying edges which lie on many shortest paths [278], extremal optimization [123], simulated annealing [172], and rounding of linear and semi-definite programs [4].

One of the more interesting heuristics by Newman [275, 276] also uses the top eigenvector of a matrix termed the modularity matrix. The key observation is the following: if we define \( a_{u,v} = 1 - \frac{d(u)d(v)}{2m} \) if there is an edge between \( u \) and \( v \), and \( a_{u,v} = -\frac{d(u)d(v)}{2m} \) otherwise, then Equation 3.5, the definition of modularity, can be rewritten as \( \frac{1}{m} \sum_i \sum_{u,v \in S_i} a_{u,v} \). Now imagine that we only want a partition into two communities, and use an indicator variable \( y_v = \pm 1 \) for each node \( v \), depending on whether \( v \) is in one or the other partition. Then, the objective can be written as \( \frac{1}{m} \sum_{u,v} a_{u,v}(1 + y_v) \). Because \( \sum_{u,v} a_{u,v} = 0 \), we can rewrite the modularity as

\[
\frac{1}{4m} y^T Ay,
\]  

(3.6)

where \( A \) is the matrix of all \( a_{u,v} \) entries, and \( y \) is a vector of \( \pm 1 \) entries describing the partition.
While finding such a vector $y$ is NP-hard, it ceases to be hard if we remove the restriction that all entries be $\pm 1$. In fact, the vector maximizing $y^T A y$ is exactly the top eigenvector of $A$. Thus, it appears to be a good heuristic to first find the top eigenvector of the modularity matrix $A$, and then round the entries of $y$. Different ways suggest themselves: the most obvious is to put all nodes with positive $y$ entries in one partition, and all nodes with negative $y$ entries in the other. But instead, one can choose any threshold $\tau$, and put all nodes with $y_v \geq \tau$ in one partition, and $y_v \leq \tau$ in the other. In fact, it probably makes sense to try all possible $\tau$ values for this.

Akin to considering other eigenvectors in Section 3.3.1, Newman [275, 276] also suggests looking at the next eigenvectors, and the partitionings defined by them. Again, heuristically, this may lead to the discovery of interesting community structure.

Since its introduction by Newman, modularity has become a very popular measure of community structure, in particular in the physics and biology communities. However, there are also important questions to ask. For instance, at what modularity measure is an observed community structure really meaningful?

Clauset et al. [90] suggest that the community structure identified by an algorithm is significant in the context of a graph when the modularity is above about 0.3. On the other hand, Guimerà et al. [173] argue that those values will be obtained merely by random fluctuations in $G(n, p)$ random graphs. While for any given community partitioning, the expected number of edges inside it under a random graph will be (essentially) equal to the baseline used in the definition of modularity, this does not hold if the community structure is chosen after the random graph is generated. A more in-depth examination of the idea of “conditioning on the degree sequence” is performed by Gaertler et al. [153]. They generalize the notion to express the “significance” of a clustering as the amount of information (edges) it contains compared to any posited baseline model. A uniform random multi-graph subject to a degree distribution is one particular example, but far from the only one. Similarly, the objective functions falling into the broad class of “modularity like” are explored further by Reichardt and Bornholdt as well [298].

3.5 Correlation Clustering

In the previous section, we began considering the problem of community identification as partitioning a graph. The assumption there was that each edge expresses that two nodes are “more likely” to be in the same community. In many cases, this is a reasonable assumption. However, in particular in contexts such as the web (including blogs), we may often find edges that explicitly suggest the two endpoints might not be in the same community. For instance, in politics or sports, pages will link to other pages with the explicit goal of deriding the content. This can be frequently identified from anchor text and similar clues. Thus, it makes sense to consider a clustering problem where edges are annotated with ‘+’ or ‘-’, expressing whether the endpoints appear to be similar or dissimilar.

This is a clustering problem called correlation clustering [25]. In the setting described above — each edge of the graph is annotated with a label of ‘+’ or ‘-‘ — the goal is to find a clustering that puts many ‘+’ edges inside clusters, and ‘-‘ edges between clusters. However, these goals may be conflicting, as can be seen for a triangle with two edges labeled ‘+’ and one labeled ‘-‘. As with the problem of maximizing modularity, we do not pre-specify the number of clusters, though of course one can also consider the variant where the target number of clusters is given.

More formally, given the graph $G = (V, E)$ on $n$ vertices, we write $+ (i, j)$ if the edge between $i$ and $j$ is labeled ‘+’ and similarly for $- (i, j)$. The optimization problem can now be expressed in two ways:

- Maximize the number of agreements, i.e., the number of ‘+’ edges inside clusters plus the number of ‘-‘ edges across clusters. This problem is called MAXAGREE.
- Minimize the number of disagreements, i.e., the number of ‘+’ edges across clusters plus the number of ‘-‘ edges inside clusters. This problem is called MINDISAGREE.

Clearly, the solution optimizing the first criterion is the same as the one optimizing the second. However, we will see that the two objective functions differ with respect with how well they can be approximated.
The correlation clustering problem is NP-hard, even for the complete graph (where each possible edge exists, and is labeled either ‘+’ or ‘-’). Hence, we are interested here in approximation algorithms.

### 3.5.1 A Simple Algorithm for Maximizing Agreements

For the maximization version, we may choose to go just after one of the two types of edges. By putting all nodes in one cluster, we get all the ‘+’ edges right — by putting each node in its own cluster, we get all the ‘-’ edges right. This suggests the following algorithm [25]:

If the number of ‘+’ edges in the graph is larger than the number of ‘-’ edges, then put everything into one cluster, else put every node in its own cluster.

**Claim 3.13** This is a $\frac{1}{2}$-approximation.

**Proof.** If the graph has $m$ edges, then the optimal solution can have at most $m$ agreements. Our algorithm produces a clustering that has at least $\frac{m}{2}$ agreements. Hence, it is a $\frac{1}{2}$-approximation. □

One consequence is that in a complete graph with assignments of edges, there must exist a clustering with at least $\frac{n(n-1)}{2}$ agreements (half of the total number of edges in the complete graph).

### 3.5.2 Improving the Approximation Factor

While the algorithm is a $\frac{1}{2}$-approximation, it is hardly satisfactory in a practical sense. We don’t need a new clustering model or algorithm if all it does is lump everything together, or put each node in its own cluster. So naturally, we want to know if the approximation guarantee can be improved.

In [25], the authors develop a PTAS (Polynomial Time Approximation Scheme) for the MAXAGREE problem in a complete graph. That is, they present a class of algorithms, parametrized with some $\epsilon > 0$, such that the algorithm with parameter $\epsilon$ is a $(1 - \epsilon)$ approximation with running time $O(n^2e^{O(1/\epsilon)})$. While this grows exponentially in $\epsilon$, for any fixed $\epsilon$, the algorithm takes polynomial time.

### 3.5.3 A 4-Approximation for MINDISAGREE in complete graphs

Given that MAXAGREE is (at least theoretically) settled for complete graphs, we next look at the minimization version. [79] gives a 4-approximation based on rounding the solution to a Linear Program.

The idea of the Linear Program is to start with an Integer Program. For each pair $i, j$ of nodes, we have a variable $x_{ij}$ which is 1 if they are in different clusters, and 0 otherwise. To make these variables consistent, we have to require that if $i$ and $j$ are in the same cluster, and $j$ and $k$ are in the same cluster, so are $i$ and $k$. This can be expressed by saying that $x_{ik} \leq x_{ij} + x_{jk}$. Subject to this consistency requirement, we want to minimize the number of ‘+’ edges between clusters, plus the number of ‘-’ edges within clusters. Thus, we obtain the following IP.

**Minimize** $\sum_{+ij} x_{ij} + \sum_{-ij} (1 - x_{ij})$

**subject to** $x_{ik} \leq x_{ij} + x_{jk}$ for all $i, j, k$

$x_{ij} \in \{0, 1\}$ for all $i, j$ 

As we know, solving IPs is itself NP-hard, so as usual, we want to relax the IP to a linear program, by replacing the last constraint with the constraint that $0 \leq x_{ij} \leq 1$ for all $i, j$. This LP can now be solved; however, the output will be fractional values $x_{ij}$. Our goal is to somehow convert these fractional values into integer ones (“round” the solution) in such a way that the objective function value of the rounded solution is not much more than that of the fractional solution. As the latter is a lower bound on the optimum integer solution, we will derive an approximation guarantee.

After obtaining the fractional $x_{ij}$ values (in polynomial) time, we need some intuition for dealing with them. We notice that the main constraint is just the triangle inequality. Hence, if we define $x_{ii} := 0$ for all

36
i, the $x_{ij}$ exactly form a metric space. In other words, the optimum solution is the metric minimizing the objective function. Nodes that are close in the metric space can be considered as being “almost in the same cluster”, while distant nodes are “quite separated”. So we likely want our clusters to consist of nodes that are mutually close. At the same time, we need to choose the boundaries carefully, so as not to cut too many edges. It turns out that the following Algorithm 2 from [79] carefully trades off between the different types of costs:

**Algorithm 2** LP-Rounding  

1. Start with a set $S$ containing all the nodes.  
2. while $S \neq \emptyset$ do  
3. Select an arbitrary node $u \in S$.  
4. Let $T := \{i \mid x_{ui} \leq \frac{1}{2}\} \setminus \{u\}$  
5. if the average distance from $u$ to the vertices in $T$ is at least $\frac{1}{4}$ then  
6. Let $C := \{u\}$ (a singleton cluster).  
7. else  
8. Let $C := T \cup \{u\}$.  
9. Remove all of $C$ from $S$.

**Theorem 3.14** Algorithm 2 is a 4-approximation.

**Proof.** The idea of the proof is to compare the mistakes that are incurred with the above algorithm against the LP cost, by showing that whenever a cluster is formed, the number of mistakes (‘+’ edges across clusters plus ‘-’ edges inside clusters) is at most four times the corresponding LP cost of the corresponding edges.

First, by simple applications of the triangle inequality and reverse triangle inequality, we obtain the following relationships between the distances (which will be used later).

**Lemma 3.15** For any nodes $i, j, u$:

- The cost of a ‘+’ edge $(i, j)$ incurred by the LP is at least $x_{ij} \geq x_{uj} - x_{ui}$.
- The cost of a ‘-’ edge $(i, j)$ incurred by the LP is at least $\max(0, 1 - x_{uj} - x_{ui})$.

We will show the claim for each cluster separately. So let $u$ be a cluster center of the cluster $C$. We can then sort the other nodes by increasing distance from $u$, i.e., whenever $x_{ui} < x_{uj}$, we say that $i < j$ (ties are thus broken arbitrarily).

We distinguish between the two cases when a singleton cluster is created, or when the cluster includes $T$.

1. Singleton cluster $\{u\}$:
   - This case occurs when the average distance from $u$ to the vertices in $T$ is at least $\frac{1}{4}$. Our algorithm then pays for all the ‘+’ edges incident with $u$. For each such ‘+’ edge whose other endpoint is not in $T$, the LP pays at least $\frac{1}{2}$. And for the edges whose other endpoint is in $T$, the LP pays at least $\frac{1}{2}$. And for the edges whose other endpoint is in $T$, the LP pays at least
     \[ \sum_{i \in T, +(ij)} x_{ui} + \sum_{i \in T, -(ij)} (1 - x_{ui}) \geq \sum_{i \in T} x_{ui} \geq \frac{|T|}{4}. \]

   Therefore, for a singleton cluster, our algorithm pays at most four times the LP cost.

2. Non-singleton clusters
   - In this case, the algorithm could be making two kinds of mistakes:
     - Allowing ‘-’ edges inside $C$.
     - Cutting ‘+’ edges between $C$ and $S \setminus C$.  

37
We first analyze the costs of negative edge mistakes. If we have a '-' edge \((i, j)\) such that both \(i\) and \(j\) are close to \(u\), then they must be close to each other, so the fractional solution must also pay a lot for this edge. Specifically, if \(x_{uj}\) and \(x_{ui}\) are both at most \(\frac{3}{8}\), then from the second part of Lemma 3.15, \(x_{ij}\) is at least \(1 - x_{uj} - x_{ui} \geq \frac{1}{4}\). So the cost for any such edge incurred by our algorithm is at most 4 times that of the LP solution.

For the remaining '-' edges, we will not be able to come up with a bound on an edge-by-edge basis. Indeed, if two nodes \(i\) and \(j\) are at distance \(\frac{1}{2}\) each from \(u\), they may be at distance 1 from each other, so the fractional solution incurs no cost. Instead, we will compare the costs node by node: specifically, we will show that for any node \(j\), the number of '-' edges between it and nodes \(i\) that are closer to \(u\) than itself is at most four times the corresponding cost of the LP solution.

So we fix a node \(j\) with \(x_{uj} \in (\frac{3}{8}, \frac{1}{2}]\). The total cost of edges \((i, j)\) with \(i < j\) (both '+' and '-' edges) that the LP solution incurs is

\[
\sum_{i < j, + (ij)} x_{ij} + \sum_{i < j, - (ij)} (1 - x_{ij}) \geq \sum_{i < j, + (ij)} (x_{uj} - x_{ui}) + \sum_{i < j, - (ij)} (1 - x_{uj} - x_{ui}).
\]

By writing \(p_j\) for the number of '+' edges \((i, j)\) with \(i < j\), and \(n_j\) for the number of '-' edges \((i, j)\) with \(i < j\), we can rewrite this as

\[
p_j x_{uj} + n_j (1 - x_{uj}) - \sum_{i < j} x_{ui}.
\]

Because the algorithm chose not to have a singleton cluster, the average distance of all nodes to \(u\) is at most \(\frac{1}{4}\). The last sum \(\sum_{i < j} x_{ui}\) only leaves out some subset of nodes furthest away from \(u\), so the average distance of those nodes is also at most \(\frac{1}{4}\). Hence, the last sum is at most \(\frac{1}{4}(p_j + n_j)\), and the entire LP cost is at least

\[
p_j x_{uj} + n_j (1 - x_{uj}) - \frac{p_j + n_j}{4} = p_j (x_{uj} - \frac{1}{4}) + n_j (1 - x_{uj} - \frac{1}{4}).
\]

The number of negative edge mistakes the algorithm will make is at most \(n_j\). On the other hand, because \(\frac{3}{8} \leq x_{uj} \leq \frac{1}{2}\), the LP cost is at least \(\frac{3}{4} + \frac{n_j}{4}\). Thus, the sum of edge costs incurred by the algorithm is at most four times that of the LP. Since this holds for any cluster, and we only account for the LP cost of any edge once, we have shown that the algorithm is a 4-approximation for '-' edges.

Next, we turn our attention to '+' edges. A '+' edge \((i, j)\) only contributes to the cost of our solution if one endpoint, say, \(i\), is included in the cluster that is formed, while the other, \(j\), is not. So we are dealing with the case that \(x_{ui} \leq \frac{1}{2}\) and \(x_{uj} \geq \frac{1}{2}\). If in fact, \(x_{uj} \geq \frac{3}{4}\), then \(x_{ij} \geq \frac{1}{4}\) by the triangle inequality, so the cost paid by our algorithm is within a factor of 4 of the LP-cost for any such edge. Hence, we now focus on the case of nodes \(j\) with \(x_{uj} \in (\frac{3}{8}, \frac{3}{4})\). We compare the number of '+' edges cut by our algorithm to the total LP cost of all edges ("+' and '-'') between \(j\) and nodes in the cluster \(C\).

By the triangle inequality (captured in Lemma 3.15), writing \(p_j\) and \(n_j\) for the number of '+' resp. '-' edges between \(j\) and nodes in \(C\), we obtain that

\[
\text{LP}_j = \sum_{i \in C, + (ij)} x_{ij} + \sum_{i \in C, - (ij)} (1 - x_{ij})
\]

\[
\geq \sum_{i \in C, + (ij)} (x_{uj} - x_{ui}) + \sum_{i \in C, - (ij)} (1 - x_{ui} - x_{uj})
\]

\[
= p_j x_{uj} + n_j (1 - x_{ui}) - \sum_{i \in C} x_{ui}.
\]
Because the algorithm didn’t form a singleton cluster, the average distance of nodes in $C$ from $u$ is at most $\frac{1}{4}$, so $LP_j$ is bounded below by $p_j x_{uj} + n_j (1 - x_{uj}) - \frac{p_j + n_j}{4}$. But $\frac{3}{4} \geq x_{uj} \geq \frac{1}{2}$, so

$$LP_j \geq \frac{p_j}{2} + \frac{n_j}{4} - \frac{p_j + n_j}{4} = \frac{p_j}{4}.$$ 

As the algorithm cuts $p_j$ ‘+’ edges, the total cost of edges cut by the algorithm is at most four times the LP-cost. By summing this over all nodes $j$ and all clusters formed, we obtain that the algorithm is a 4-approximation.

We mentioned above that, while the optimal solution for the minimization and maximization version is the same, the approximation guarantees differ. For the minimization version on complete graphs, the algorithm from [79] we just analyzed gives a 4-approximation. On the other hand, [79] also shows that the problem is APX-hard, i.e., there is some constant such that the minimization version on complete graphs cannot be approximated to within better than that constant unless P=NP. On arbitrary graphs, the best known approximation is $O(\log n)$; however, it is open whether the problem can be approximated to within a constant.

Even though there is a constant-factor approximation for the minimization version on complete graphs, together with an APX-hardness result, we may wonder what is the best possible constant. We will show that the LP used above has an integrality gap of 2, so no rounding approach solely based on that LP can yield an algorithm with a better guarantee. The example is the “wheel” graph, in which all nodes of a $k$-cycle are connected to one additional center node with a ‘+’ edge (while all edges of the cycle are labeled ‘-‘). Then, the integral optimal solution puts all of the cycle nodes in different clusters, paying a total of $k-1$, while the fractional optimum assigns $x_{ij} = \frac{1}{2}$ to all edges between the center and the cycle nodes, paying $\frac{k}{2}$. The ratio approaches 2 as $k \to \infty$.

The authors of [79] show that no rounding algorithm based on the same type of “region growing” will lead to an approximation guarantee of better than 3, and conjecture that in fact, no similar approach will give a better approximation than the factor of 4 obtained.

For the maximization version, there is a PTAS on complete graphs [25]. For graphs that are not complete, the problem is APX-hard. However, in this case, it is known how to approximate it to within a constant factor. The factor of 0.7664 from [25] was improved to an 0.7666 approximation via semi-definite programming by Swamy [321].

### 3.6 Further Reading

Since the topic of community discovery in graphs has attracted so much attention (in particular in physics and biology), there is a very large literature by now. Several good surveys exists, including ones by Fortunato and Castellano [145, 147], Danon et al. [103], and Porter et al. [292]. They share a focus on physics-style heuristics, and are short on surveying important techniques (such as approximation algorithms or rounding of linear programs) from the computer science literature. Good overviews of some of the basic techniques common in computer science are given in Part II of the book “Network Analysis” by Brandes and Erlebach [62] and a survey by Schaeffer [307].

Naturally, the identification of cohesive groups of individuals has long been an important task for sociologists as well. Some of the traditional approaches can be found in classical texts on social networks by Wasserman and Faust [333] and Scott [311]. An overview of some of the more recent work on group cohesiveness is given in the article by White and Harary [338].

The definition of Flake et al. [144] is called a strong community because of the stringent condition on each node. One can instead consider weak communities. In a weak $\alpha$-community $S$, we only require that the total number of edges inside $S$ be at least an $\alpha$ fraction of the sum of degrees of all nodes in $S$. This definition was proposed by Radicchi et al. [293], simultaneously with a paper by Hayrapetyan et al. considering the same motivation and definition [178]. As with the definition by Flake et al., the entire graph is always a weak 1-community. So again, one could like to find the smallest weak $\alpha$-communities. Currently, no approximation
algorithm with provable guarantees is known for this problem. On the hardness side, only NP-hardness is known. Lemma 3.9 is an example of the topics studied in the area of extremal graph theory [51]. Broadly, extremal graph theory studies the relationships between different graph properties. For instance, in the case of the lemma, we see that a certain number of edges implies that the graph must have bipartite cores.

As discussed in Section 3.4, maximizing modularity has become very popular, in particular among physicists. The paper [275] contains a fairly detailed survey of much of the work that has been done on the topic. Among the more critical papers is one by Fortunato and Barthélémy [146], which proves that modularity has an intrinsic resolution limit, in the sense that the modularity of a proposed clustering (in a large graph) can always be improved by merging smaller communities into larger ones. In fact, Fortunato and Barthélémy show that this is the case for every clustering objective that summarizes the quality of a partitioning in a single number.

There is a direct analogy between the modularity maximization problem and correlation clustering, as observed in [4]. Consider again Equation (3.6). There is really no reason why \( A \) would have to be restricted to be the modularity matrix. For instance, if instead, we made \( A \) the matrix with +1 for ‘+’ edges, and -1 for ‘-’ edges, then we measure the difference between correctly and incorrectly classified edges. By adding \( |E| \) and dividing by two, we thus get the number of correctly classified edges, i.e., the MaxAgree objective. Since \( |E| \) is a constant, this does not change optimality of solutions. Indeed, the formulation is the basis of the approximation algorithm by Swamy [321], who adds the constraint that \( y_i^2 = 1 \) for all \( i \), and then uses semi-definite programming with an appropriate rounding scheme.

One can also apply the same reasoning to the LP (3.7). Instead of having terms \( x_{ij} \) and \( 1 - x_{ij} \) in the objective function for edges labeled ‘+’ resp. ‘-’, we can write \( \sum_{ij} a_{ij} x_{ij} \) for the modularity matrix entries \( a_{ij} \), under the same constraints. While the objective function now takes a different form, we can still aim to apply the same kind of rounding techniques. Unfortunately, they do not come with approximation guarantees any more.

The previous discussion suggests a more general question: given a matrix \( A \), find a \( \pm 1 \) vector \( y \) maximizing \( y^T A y \). This would subsume both modularity maximization and correlation clustering. Some recent work has identified conditions on \( A \) under which provable approximation guarantees can be obtained. Nesterov [269] shows that if \( A \) is positive semi-definite, then solving the semi-definite program alluded to above and rounding appropriately gives a \( 2/\pi \) approximation. Charikar and Wirth [80] give an \( \Omega(1/\log n) \) approximation algorithm based on semi-definite programming and rounding, under the assumption that the diagonal elements of \( A \) are all zero. While this is the case for correlation clustering, it does not apply to modularity maximization.
Chapter 4

Link-Based Classification

When looking at correlation clustering, we started from the motivation that nodes that have ‘+’ edges between them are more likely to belong to the same cluster. In a sense, we could then consider the clusters we formed as communities, sharing perhaps some similarity in topic. If we have an a priori estimate of the topics that nodes (web pages) are about, we can use a similar approach to correct and refine our estimates. For instance, if a page about coffee points to a lot of pages about programming, then perhaps we misinterpreted the meaning of “Java” in one case, and should revise our initial estimate. Hence, we will try to optimize two conflicting goals: agreements with an a priori estimate, and agreements between nodes with edges between them.

Essentially the same problem arises frequently in the context of computer vision [59]. The goal there is to label pixels into classes (classification problem) and assign labels representing these classes as foreground or background objects (or different colors). Again, the competing constraints are that we don’t want many disagreements with the a priori labels, but also not between adjacent (physically close) pixels.

4.1 Markov Random Fields and the Optimization Problem

To formalize this intuition, we can think about the following model, first proposed in the context of link-based text classification by Chakrabarti, Dom, and Indyk [74], and later also studied by Broder, Krauthgamer, and Mitzenmacher [64]. Our goal is to assign some label $f(v)$ to each node $v$ in the graph. Focus on one node $v$, and consider its neighborhood $\delta(v)$. (For now, we will ignore the direction of edges.) Suppose that we knew the (correct) labels $f(u)$ of each neighbor $u \in \delta(v)$. With this knowledge, a classifier could give us a conditional probability distribution over the possible labels of $v$, e.g., by taking text and other content at node $v$ into account. This is exactly the definition of a Markov Random Field (MRF) (see, e.g., [38, 288]).

Notice that the specification of the distribution at each node $v$ is still exponential in the size of $\delta(v)$, but if nodes have low degree, this allows us to specify probability distributions concisely.

Of course, we do not know the correct labels in neighborhoods of $v$, either, so what we are really looking for is a complete labeling maximizing the overall probability, the product over all nodes $v$ of their conditional probability. As expected, this is still a hard problem; among others, it contains Graph Coloring or Independent Set as special cases.

We can transform this optimization problem into a more familiar looking one, as described by Kleinberg and Tardos [221]. First, by the Hammersley-Clifford Theorem [38], the probability can be written as

$$\frac{1}{Z} \exp(- \sum_{C \in \mathcal{C}} \Gamma_C(f|_C)).$$

Here, $Z$ is a normalizing constant, $\mathcal{C}$ is the collection of all cliques in the graph, and $\Gamma_C$ is a clique potential. The important thing here is that the clique potentials depend only on the restriction $f|_C$ of the labeling to $C$, i.e., only on the labels of the nodes in the clique.
The real observation will not necessarily be drawn from the given Markov Random Field. Instead, we assume that it was obtained by first drawing a labeling from the Markov Random Field, and then introducing independent random noise at each node, according to known distributions.

If we further assume that only pairwise interactions matter, then instead of all cliques \( C \), we only need to look at penalty terms \( \Gamma_e \) for edges. In principle, these terms could depend in arbitrary ways on the labels of the edge endpoints. However, it is natural to assume homogeneity: each \( \Gamma_e \) is of the form \( w_e \cdot \Gamma \). That is, the type of penalty is the same for all edges; only the weights are different. Combining all of these assumptions, and using Bayes’ Rules, we obtain the following form for the probability that the true labeling is \( f \), given that the observed labeling is \( f' \).

\[
\frac{1}{\gamma Z} \prod_v \text{Prob}[f'(v) \mid f(v)] \cdot \exp \left( - \sum_{e=(u,v) \in E} w_e \Gamma_e(f(u), f(v)) \right).
\]

Here, \( \gamma \) is a constant depending only on \( f' \), but not on \( f \). The terms in the first product are the probabilities of observing labels \( f'(v) \) if the true underlying labels are \( f(v) \), and are thus given by the description of the random noise.

By taking logarithms, we see that maximizing the probability is the same as minimizing the objective

\[
\sum_v \log \frac{1}{\text{Prob}[f'(v) \mid f(v)]} + \sum_{e=(u,v) \in E} w_e \Gamma_e(f(u), f(v)).
\]

This motivates the following, somewhat more general, optimization problem: Given a graph \( G \), where each edge \( e = (u, v) \) has weight \( w_e \) representing the cost or strength of the relationship between vertices \( u \) and \( v \), we want to assign labels \( f(v) \in L \) for all vertices \( v \) such that we minimize the assignment cost

\[
\sum_v c(v, f(v)) + \sum_{e=(u,v)} w_e \cdot d(f(u), f(v)). \tag{4.1}
\]

\( c(v, f(v)) \) is the cost of choosing label \( f(v) \) for node \( v \), which will be a result of deviating from the observed or otherwise determined label. \( d \) represents how “different” the two labels are, and how much penalty therefore should be assigned to adjacent nodes with these labels. Since \( d \) is a measure of dissimilarity, it makes sense to assume that it is a metric. From now on, we will therefore focus on the metric Markov Random Field labeling problem.

In fact, for the remainder of this chapter, we will focus on the special case of the uniform labeling problem. In that case, for two labels \( a, a' \), we have that \( d(a, a') = 0 \) if \( a = a' \), and \( d(a, a') = 1 \) otherwise. Thus, labels are either identical, or simply different — there are no gradations of “different”. We will discuss the general problem briefly at the end of the chapter.

Remark 4.1 In the physics literature, this type of Markov Random Field is also called the Ising Model. It is used to describe the distributions of states of spin systems such as magnets. In those systems, adjacent atoms tend to have the same orientation of spin in ferro-magnetic states. Depending on the temperature, the state will be more or less aligned, and physicists try to answer how likely different spin configurations will be at different temperatures.

4.2 A 2-approximation using LP-rounding

In this section, we will present and analyze a polynomial-time 2-approximation algorithm using LP-rounding. The algorithm and analysis were given by Kleinberg and Tardos [221]. Recall that we focus on uniform Markov Random Fields. That is, our goal is to find a labeling \( f(v) \) of vertices \( v \) minimizing

\[
\sum_v c(v, f(v)) + \sum_{e=(u,v), f(u) \neq f(v)} w_e. \tag{4.2}
\]
To formulate this problem as an (integer) linear program, we define variables $x_{v,a}$, which will be equal to 1 if node $v$ is assigned label $a$, and equal to 0 otherwise.

Using the $x_{v,a}$ values, we would like to define a variable $y_e$ that is equal to 1 if the two endpoints of $e$ have different labels, and equal to 0 otherwise. To do this, we notice that if the endpoints $u$ and $v$ of $e$ have the same label, then $x_{v,a} - x_{u,a} = 0$ for all $a$. Otherwise, it is 0 for all but two of the $a$, and equal to 1 and -1 for the other two. Thus, for $e = (u,v)$,

$$y_e = \frac{1}{2} \sum_a |x_{v,a} - x_{u,a}|.$$

However, this is still no linear constraint, as $|x_{v,a} - x_{u,a}|$ is not a linear function. We will see in a moment how to deal with this problem. For now, notice that our goal is to minimize

$$\sum_{v,a} x_{v,a} \cdot c(v,a) + \sum_e w_e \cdot y_e. \tag{4.3}$$

The only constraint is that each node should obtain exactly one label, so $\sum_a x_{v,a} = 1$ for all $v$, and each $x_{v,a} \in \{0, 1\}$.

To return to the issue of how to express $y_e$ via linear constraints, we first write $y_e = \frac{1}{2} \sum_a y_{e,a}$, where $y_{e,a} = |x_{v,a} - x_{u,a}|$. While we cannot directly express this as a linear constraint, we can require that $y_{e,a} \geq |x_{v,a} - x_{u,a}|$, by writing $y_{e,a} \geq x_{v,a} - x_{u,a}$ and $y_{e,a} \geq x_{u,a} - x_{v,a}$ for all $e = (u,v)$. But since the objective is to minimize the objective function (4.3), the optimum solution will never choose $y_{e,a}$ larger than necessary, so we will indeed have $y_{e,a} = |x_{v,a} - x_{u,a}|$. In summary, we have derived the following LP for the metric labeling problem:

Minimize

$$\sum_{v,a} x_{v,a} c(v,a) + \sum_e w_e \cdot y_e$$

subject to

$$\sum_a x_{v,a} = 1 \quad \text{for all } v$$

$$y_e = \frac{1}{2} \sum_a y_{e,a} \quad \text{for all } e$$

$$y_{e,a} \geq x_{v,a} - x_{u,a} \quad \text{for all } e = (u,v), a$$

$$y_{e,a} \geq x_{u,a} - x_{v,a} \quad \text{for all } e = (u,v), a$$

$$x_{v,a} \in \{0, 1\} \quad \text{for all } v, a.$$

As usual with LP rounding algorithms, we relax the constraint $x_{v,a} \in \{0, 1\}$ to $x_{v,a} \in [0, 1]$, i.e., we allow the solutions to take on fractional values. Then, we can solve the LP in polynomial time. We will investigate how to round the solution to obtain a labeling which does not perform much worse than the fractional optimum. In particular, that will show that it is not much worse than the integral optimum as well.

For the underlying idea, notice that the $x_{v,a}$ summed over all $a$, add up to 1 for each $v$. Hence, we can view them as fractions to which we would label each node $v$ with $a$ with probability $x_{v,a}$. Notice that we would do very well on the labeling cost $\sum_{v,a} x_{v,a} \cdot c(v,a)$ this way, as

$$E[\text{labeling cost of node } v] = \sum_a \text{Prob}[v \leftarrow a] \cdot c(v,a) = \sum_a x_{v,a} \cdot c(v,a),$$

Here (and later), we denote by $v \leftarrow a$ the fact that node $v$ is assigned label $a$, and use the fact that $\text{Prob}[v \leftarrow a] = x_{v,a}$.

However, this does not ensure that we do well on separation costs, too. In fact, if we make the choices for different nodes $v$ independently, it can happen that this rounding procedure does very badly. Suppose that the graph consists of just two vertices, $v_1$ and $v_2$, with a non-zero cost edge between them. There are two labels $a$ and $b$, and all vertex labeling costs are $c \equiv 0$. Then, the fractional solution can choose any fractional assignment, so long as $v_1$ and $v_2$ have the same fractional share of $a$. One of the optimal solutions of the LP, which we may have to round, would thus be assigning all probabilities equal to $\frac{1}{2}$. The LP cost is then 0. But if $v_1$ and $v_2$ receive different labels, we pay a non-zero separation cost, and thus have infinitely bad approximation.
Thus, we need to make more coordinated choices between labels of nodes and their neighbors. Suppose that we label a node \( v \) with the label \( a \). A first stab would be to label all neighbors \( u \) of \( a \) with \( x_{u,a} \geq x_{v,a} \) with \( a \) as well. But then, \( u \)'s label may conflict with that of one of its neighbors when it is labeled later. So we may also want to require that \( a \) be the label with the largest fractional value among \( u \)'s labels. But that does not really solve the problem either. It seems that we may also have to label the neighbors of \( u \), and so on recursively.

As a result, one thing we can do is that once we pick a node \( v \) and assign \( a \) to it, we pick \( x_{v,a} \) as a threshold, and assign \( a \) to all nodes \( u \) having \( x_{u,a} \geq x_{v,a} \). Then, when a conflict occurs between some such \( u \) and another node \( w \) which gets a different label \( b \), at least we know that they didn’t have exactly identical fractional values.

This motivates the following rounding algorithm, which picks a label \( a \) uniformly at random, as well as a random threshold, and labels all nodes with that particular label if their probability for the chosen label is above the threshold. This process is repeated until there are no more nodes left to be labeled.

### Algorithm 3 LP Rounding algorithm

1: repeat
2:   Pick a uniformly random label \( a \in L \), and a uniformly random \( \alpha \in [0,1] \).
3:   Label all \( u \) with \( x_{v,a} \geq \alpha \) with \( a \), and remove them.
4: until no more nodes are left

#### 4.2.1 Approximation Guarantee

**Theorem 4.2** Algorithm 3 is a 2-approximation.

**Proof.** We prove the theorem by analyzing the assignment costs and separation costs separately. We show that for both cases, the cost incurred using the approximation algorithm is within a factor of two of the optimum fractional solution of the LP, and hence also within the same factor of the integral solution (which can be no better than the best fractional one). Let \( k = |L| \) denote the number of labels.

1. To analyze the assignment cost, we first notice that in any particular iteration, \( v \) is labeled \( a \) with probability \( \frac{x_{v,a}}{k} \). That is because label \( a \) is picked with probability \( 1/k \), and conditioned on picking label \( a \) we label \( v \) iff the random threshold is at most \( x_{v,a} \).

   Because in each iteration, label \( a \) is assigned to \( v \) with probability proportional to \( x_{v,a} \), the overall probability of assigning \( a \) to \( v \) is \( x_{v,a} \). We can thus use the same argument as above to obtain that the expected labeling cost is \( \sum_a \text{Prob}[v \leftarrow a] \cdot c(v,a) = \sum_a x_{v,a} \cdot c(v,a) \).

2. To analyze the separation cost, we first notice that we only incur cost \( w_{u,v} \) for the edge \( e = (u,v) \), when \( u \) and \( v \) get different labels. But that can only happen when \( u \), \( v \) are labeled in different iterations of the algorithm. Therefore, there must have been an iteration labeling exactly one of \( \{u,v\} \). We bound the probability of that happening.

   Given a label \( a \), exactly one of \( \{u,v\} \) is labeled \( a \) iff the threshold \( \alpha \in (x_{u,a}, x_{v,a}] \) (assuming \( x_{u,a} < x_{v,a} \)). So the probability of labeling exactly one of \( u \) and \( v \) with label \( a \) is \( |x_{v,a} - x_{u,a}| \). Summing over all labels \( a \), the probability of labeling exactly one of \( u \) and \( v \) with any label is \( \sum_a \frac{|x_{v,a} - x_{u,a}|}{k} \).

   For any time \( t \), we let \( F_t \) be the event that at least one of \( \{u,v\} \) is labeled in iteration \( t \) and \( E_t \) the event that exactly one is labeled. We are thus interested in \( \text{Prob}[E_t \mid F_t] \). First off, notice that \( E_t \subseteq F_t \), so \( \text{Prob}[E_t] = \text{Prob}[E_t \cap F_t] = \text{Prob}[E_t \mid F_t] \cdot \text{Prob}[F_t] \). Rearranging yields that \( \text{Prob}[E_t \mid F_t] = \frac{\text{Prob}[E_t]}{\text{Prob}[F_t]} \). But we already calculated \( \text{Prob}[E_t] \) above. And the probability of \( F_t \), the event that at least one of \( u, v \) is assigned a label, is at least \( 1/k \), the probability that \( u \) is assigned a label. Hence, we have that

\[
\text{Prob}[E_t \mid F_t] \leq \frac{\sum_a \frac{1}{k} |x_{u,a} - x_{v,a}|}{k} = \sum_a |x_{u,a} - x_{v,a}| = 2w_{u,v}.
\]
Applying this at the time $t$ where the first of $u, v$ is actually assigned a label gives us that they are separated with probability at most $2y_{u,v}$. (Notice that $t$ is a random variable, but this does not really hurt us here.)

Hence, the expected total separation cost is

\[
E[\text{separation cost}] = \sum_{e=(u,v)} \text{Prob}[e \text{ separated}] \cdot w_e \leq \sum_{e=(u,v)} 2y_{u,v}w_e.
\]

Summing up over the two cases, the total cost is at most

\[
\sum_{v,a} x_{v,a}c(v, a) + 2\sum_{e=(u,v)} w_ey_{u,v} \leq 2 \cdot \text{cost(LP-OPT)}.
\]

Thus, the algorithm is a 2-approximation.

Having obtained a 2-approximation, we naturally want to know if we can do better. While this question has not been entirely resolved, we will show here that we cannot do better than a 2-approximation if the bound on the optimum uses only the LP solution. We exhibit an integrality gap of 2, i.e., we show that there are instances where the best integral solution is worse than the best fractional one by a factor arbitrarily close to 2.

The example consists of a complete graph $K_n$ with edge weights of 1 for each edge. There are $n$ labels, one corresponding to each vertex. The labeling cost is $c(v, v) = \infty$, and $c(v, a) = 0$ for all nodes $u \neq v$. That is, no node can be labeled with its own name, but any other labeling is free.

The best integral solution assigns each node label “1”, except for node 1, which is given label “2”. Then, exactly the edges incident with node 1 are cut, so the total cost is $n - 1$. On the other hand, a fractional solution can assign

\[
x_{v,a} = \begin{cases} 0, & \text{if } a = v \\ \frac{1}{n-1}, & \text{if } a \neq v. \end{cases}
\]

Then, the cost incurred by each edge $(u, v)$ is $y_{u,v} = \frac{1}{n-1}$, so the total cost is \( \frac{1}{n-1} \binom{n}{2} = \frac{n}{2} \). Hence, the integrality gap is \( \frac{n-1}{2} \). As $n \to \infty$, this integrality gap approaches 2.

### 4.3 Uniform Labeling via Local Search

We can avoid having to solve an LP (which, while polynomial, tends to be fairly slow), by using a different 2-approximation algorithm based on Local Search where search moves use Min-Cut computations [60, 330]. In local search algorithms, we start with a solution, and repeatedly apply one of several simple moves, as long as it leads to an improvement. When no more improvement is possible with the simple moves, the algorithm terminates. It is rare that provable guarantees can be shown for efficient local search algorithms, but this is one of the cases.

Here, a local search move works as follows: We pick a label $a$ and try if converting other vertices to that label will reduce the total cost. So we may label an arbitrary additional set of vertices with $a$, but no vertex gets its label changed to anything except $a$ in one step. Among all such new labelings, we choose the best one, and then iterate over all labels. We thus obtain the following algorithm:

Notice that in this algorithm, the labeling $f_a$ can be found using a single Min-Cut computation. Indeed, the algorithm is based on an insight by Greig et al. [171], who showed that for just two labels, the optimum solution can be found in polynomial time using a single Min-Cut computation.

\[\text{Remark 4.3} \] Perhaps the most natural local search algorithm would only consider relabeling one vertex at a time, or a constant number. It is easy to see that any such algorithm can get stuck in highly suboptimal local minima, and no approximation guarantee can be proved. Allowing wholesale relabeling of many vertices is crucial.

45
4.3.1 Approximation Guarantee

**Theorem 4.4** The algorithm produces a 2-approximation.

**Proof.** The idea behind the proof is to show that the termination condition — the fact that no single-label relabeling yielded any more improvement — is enough to be close to the optimum solution. Roughly, we will do this by “decomposing” the optimum solution into the algorithm’s solution.

Specifically, let \( f \) be the algorithm’s solution and \( f^* \) the optimal solution. For each \( a \), we let \( S_a = \{ v \mid f^*(v) = a \} \) be the set of all nodes that the optimum labeled \( a \). We define a labeling \( f_a \) that “interpolates” between \( f \) and \( f^* \), by saying that \( f_a(v) = a \) if \( f^*(v) = a \), and \( f_a(v) = f(v) \) otherwise. That is, the interpolating labeling agrees with \( f^* \) whenever \( f^* \) chose \( a \), and with the algorithm’s labeling otherwise.

Because \( f_a \) was a candidate for relabeling from \( f \), the termination of the algorithm implies that \( \gamma(f) \leq \gamma(f_a) \) for all labels \( a \). We now divide the cost of \( f \) into that incurred inside \( S_a \), outside \( S_a \) and across the boundary. So we define

\[
\gamma_S(f) := \sum_{v \in S} c(v, f(v)) + \sum_{(u, v) \in S \times f(u) \neq f(v)} w_{(u, v)}
\]

for any set \( S \). As a result, we can rewrite \( \gamma(f) \) and \( \gamma(f_a) \) as

\[
\gamma(f) = \gamma_{S_a}(f) + \gamma_{\overline{S}_a}(f) + \sum_{(u, v) \in S_a \times \overline{S}_a, f(u) \neq f(v)} w_{(u, v)};
\]

\[
\gamma(f_a) = \gamma_{S_a}(f_a) + \gamma_{\overline{S}_a}(f_a) + \sum_{(u, v) \in S_a \times \overline{S}_a, f_a(u) \neq f_a(v)} w_{(u, v)}.
\]

Above, we argued that \( \gamma(f) \leq \gamma(f_a) \), and because the two labelings agree on the set \( \overline{S}_a \), we have that \( \gamma_{\overline{S}_a}(f) = \gamma_{\overline{S}_a}(f_a) \). On the other hand, \( f_a \) and \( f^* \) agree on \( S_a \), so \( \gamma_{S_a}(f_a) = \gamma_{S_a}(f^*) \). Taken together, this implies that

\[
\gamma_{S_a}(f) + \sum_{(u, v) \in S_a \times \overline{S}_a, f(u) \neq f(v)} w_{(u, v)} \leq \gamma_{S_a}(f^*) + \sum_{(u, v) \in S_a \times \overline{S}_a, f_a(u) \neq f_a(v)} w_{(u, v)}.
\]

But under the last sum, whenever \( f_a(u) \neq f_a(v) \), then also \( f^*(u) = a \neq f^*(v) \), so we can upper bound the sum by replacing \( f_a(u) \) and \( f_a(v) \) with \( f^*(u) \) and \( f^*(v) \). Thus, we have derived

\[
\gamma_{S_a}(f) + \sum_{(u, v) \in S_a \times \overline{S}_a, f(u) \neq f(v)} w_{(u, v)} \leq \gamma_{S_a}(f^*) + \sum_{(u, v) \in S_a \times \overline{S}_a, f^*(u) \neq f^*(v)} w_{(u, v)}.
\]

Because the \( S_a \) for all \( a \) form a disjoint cover of all nodes, summing up over all \( a \) now gives us that
\[
\gamma(f) \leq \sum_a \gamma_S(a) + \sum_{(u,v) \in S \times S, f(u) \neq f(v)} w(u,v) \\
\leq \sum_a \gamma_S(f^*) + \sum_{(u,v) \in S \times S, f^*(u) \neq f^*(v)} w(u,v) \\
\leq 2\gamma(f^*),
\]

because in the last sum, each edge between differently labeled nodes is counted exactly twice, while the assignment cost for each node is counted exactly once. So the algorithm is a 2-approximation.

### 4.3.2 Running Time

For a given labeling \( f \), we let \( \gamma(f) = \sum_v c(v,f(v)) + \sum_{v=(u,v),f(v)\neq f(u)} w_e \) denote its total labeling cost. The algorithm only chooses \( f_a \) over \( f \) if it is better, so whenever changes happen, \( \gamma \) is strictly decreasing. As a result, the algorithm cannot loop. Further, if all \( w_e \) and \( c(v,a) \) are integers, then the decrease will be at least 1 in each iteration, so there can be at most \( W = \sum w_e \) iterations. Therefore, the algorithm runs in pseudopolynomial time.

In order to make the running time actually polynomial, we can apply a standard trick (see, e.g., [19]): only perform an update step if it decreases the cost by a factor of at least \( (1 - \frac{\epsilon}{p(n)}) \) for some polynomial \( p \) and some small constant \( \epsilon \). At the cost of an additional factor depending on \( \epsilon \) in the approximation guarantee, this ensures that the number of steps is at most \( \log W / \log (1 - \epsilon / p(n)) \). (The analysis for this becomes a bit more messy, though.) At this point, it appears to be open whether this or a similar local search algorithm runs in strongly polynomial time, i.e., time not depending on \( W \) at all.

### 4.4 Further Reading

In the presentation of approximation algorithms, we focused on the case of uniform metric labeling. The paper by Kleinberg and Tardos [221] also gives an approximation algorithm for arbitrary metrics \( d(a, a') \) on labels \( a, a' \). The idea is to first probabilistically embed the metric into a hierarchically well-separated tree metric with low distortion. A hierarchically well-separated tree metric is defined by assigning distance labels on the edges of a tree. These distance labels must decrease exponentially from the root to the leaves. For each node pair, their distance is the length of the unique path in the tree. The distortion of such an embedding here is the largest factor by which the distance of any pair \( a, a' \) increases when going from the metric \( d \) to the tree metric. (Here, the tree metric has to be such that no distance shrinks.)

[221] shows how a modification of the LP-rounding based approach presented in Section 4.2 gives a constant-factor approximation for hierarchically well-separated tree metrics. Combining this result with the fact that each metric on \( k \) points can be probabilistically embedded into hierarchically separated tree metric with expected distortion \( O(\log k \log \log k) \) [32, 31] now provides an \( O(\log k \log \log k) \) approximation algorithm if there are \( k \) different labels to assign. Subsequently, Fakcharoenphol et al. [136] improved the approximation guarantee for probabilistic embeddings into hierarchically well-separated trees to \( O(\log k) \), immediately implying the same approximation guarantee for the algorithm of Kleinberg and Tardos.

Several papers focus on algorithms for particular types of metrics. For example, Boykov, Veksler and Zabih [59] show that if the metric is the linear metric \( d(a,a') = |a - a'| \), which applies for instance to pixel intensities, then an optimum solution can be found using minimum-cut computations. Gupta and Tardos [174] consider the variation of the truncated linear metric \( d(a,a') = \min(M, |a - a'|) \), which assigns all pairs beyond a certain maximum distance the same penalty. For this version, they show that a variation of the local search algorithm presented in Section 4.3 gives a 4-approximation in polynomial time.

Chekuri et al. [81] give an interesting novel LP-formulation for the general metric problem. As special cases, they recover a 2-approximation for the uniform metric, an \( O(\log k) \) approximation for arbitrary metrics, and a \( 2 + \sqrt{2} \) approximation for truncated linear metrics.

47
The objective functions discussed in this chapter are all concave. This corresponds to penalties increasing more slowly as the assigned labels grow more and more different. If the increase in penalties instead grows more steep, the “metric” \( d \) will instead be convex. In fact, this simplifies the problem significantly. Hochbaum [180] shows that if the function \( d(a, a') \) is convex (in fact, even if there are different such convex functions on each edge), and each node penalty function \( c(v, a) \) is convex in \( a \), then the problem can be solved optimally in polynomial time.

The idea of using cut-based algorithms for classification has proved useful in other applications as well. In particular, Blum and Chawla [48] propose using it in the context of semi-supervised learning [77], where a small set of labeled examples for learning is augmented by many more unlabeled training examples. These unlabeled examples can be labeled by considering a graph representation of similarities. If the labels are only binary (“Yes” or “No”), then the best labeling of the additional instances can be obtained using a Min-Cut algorithm.

Classification problems of this type have been studied extensively in the literature on machine learning. For example, Taskar et al. [323] study the problem for learning the weights of edges (or, more generally, cliques) from example data. They also propose heuristics based on the popular belief-propagation algorithm [288] for the classification problem. However, these heuristics come with no provable guarantees. Subsequent to the paper by Taskar et al., several papers have proposed heuristics for dealing with sparse labels in a networked classification task. For some examples and overviews, see [156, 270, 242, 154].
Chapter 5

Rank Aggregation and Meta-Search

In previous chapters, we had discussed the problem of searching for relevant results on the WWW by exploiting the link structure. Nowadays, there are already multiple search engines giving quite good results. However, given that there are so many search engines already, employing different techniques, we may be interested in combining their positive features, and constructing a meta-search engine that uses all of their results [312, 253]. This leads very naturally to the problem of rank aggregation: given several orders on items (such as search results), determine a consensus ordering, which somehow reflects all orderings together.

This problem can be studied from different perspectives. We could treat it as a machine learning problem, by treating each search engine as an “expert”, and trying to learn which engine’s advice to trust, based on past performance. The goal would then be to perform no worse than the best expert, but without knowing ahead of time which one is actually the best expert.

Another approach is to consider it as an optimization problem. By defining an appropriate notion of distance between rankings, we can then look for a ranking that is close to all given rankings in a certain sense.

A third approach, and the one we begin with, exploits the close connection between rank aggregation and voting. In both settings, we have items (candidates or web pages), and orderings on them (voter preferences or search engine rankings). The goal is to find a ranking that is “agreeable” to all voters/search engines. By using this analogy, we can leverage several centuries of thought on the issue of voting and social choice.

5.1 Rank Aggregation as Social Choice

To get a feel for the difficulties in determining a consensus ordering, let us look at a simple example where a first choice is to be determined. Suppose there are three candidates, GWB, AG, and RN, and the voters’ preferences are as follows:

- 49% of voters prefer the order GWB–AG–RN
- 48% of voters prefer AG–GWB–RN
- 3% of voters prefer RN–AG–GWB

Who should be considered the winner legitimately? There are different plausible answers, depending on the view we take of these preferences. We could argue that GWB was the candidate desired as winner by the largest number of voters, and hence should be the winner. This is the outcome of plurality voting. On the other hand, we could argue that a majority of voters prefers AG over GWB, and also over RN. Hence, AG wins all pairwise comparisons, and should be the winner.

Over the years, different rules have been proposed for deciding on a winner, or even a complete ranking. One of the earliest is due to Borda [107]. His method essentially uses the average position of a candidate, averaged over all voters’ preferences, and sorts candidates in this order. More formally, for each voter, each
candidate obtains a score equal to the number of other candidates he beat in this voter’s ranking. Candidates are then ranked by the sum of their scores, summed over all voters. In our above example, this would give us AG as the winner, as the score would be $49 \cdot 1 + 48 \cdot 2 + 3 \cdot 1 = 148$, whereas the score of GWB is $49 \cdot 2 + 48 \cdot 1 = 146$ (and the score of RN is $3 \cdot 2 = 6$).

Although Borda’s rule looks useful at first, a major problem with it is that even the candidate with the largest number of pairwise wins can lose. For instance, if 51% of the voters prefer the order ABC, and 49% prefer BCA, A’s score is $51 \cdot 2 = 110$, while B scores $49 \cdot 2 + 51 = 149$, thus winning the election though clearly more people prefer A to B than vice versa.

Taking the latter idea further, we may wish to attain the following property, called Condorcet Criterion (after N. de Condorcet): If candidate A beats candidate B in pairwise comparison (i.e., more voters prefer A to B than vice versa), then A should precede B. We quickly run into a problem, already observed by Condorcet [108]: If three voters have the preference orders ABC, BCA, and CAB on three candidates A, B, C, then A should precede B, B should precede C, and C should precede A. Obviously, not all three can be accomplished simultaneously.

A relaxed version of the property, called Extended Condorcet Criterion (XCC) was proposed by Truchon [325], and requires the following: If $X, Y$ are a partition of the set $V$ of all candidates (i.e., $X \cap Y = \emptyset$ and $X \cup Y = V$), and for all $x \in X, y \in Y$, $x$ beats $y$ in direct comparison, then all of $X$ should precede all of $Y$.

This version is much less restrictive; in particular, for the example given above, it allows us to choose an arbitrary ranking of candidates. In fact, we can show that for any input orderings, there is always an ordering satisfying the XCC.

**Proposition 5.1** For any input rankings, there is a consensus ordering satisfying the XCC.

**Proof.** Consider the directed graph $G$ on $V$ with an edge from $x$ to $y$ if $x$ beats $y$ in pairwise comparison (we also write $x > y$ for this). Notice that this is a tournament graph, i.e., a graph in which for each pair of nodes, there is an edge one way or the other. We prove below that every tournament graph contains a Hamiltonian Path. Consider the candidate ordering determined by the sequence of vertices along such a Hamiltonian Path. Assume that this ordering fails the XCC. Then, there is a partition $(X, Y)$ of $V$ such that each $x \in X$ beats each $y \in Y$ in direct comparison, yet some $y \in Y$ precedes some $x \in X$ in this ordering. Then, there must also be an adjacent such pair in the ordering, i.e., one where $y$ immediately precedes $x$ on the Hamiltonian Path. But this is a contradiction, as there must have been an edge from $y$ to $x$ in the Hamiltonian Path, so $y$ must actually beat $x$ in direct comparison.

**Lemma 5.2** Each tournament graph $G$ contains a Hamiltonian Path.

**Proof.** We use induction on the number of vertices, $n$. For $n = 1$, the claim is trivial. For $n \geq 2$, let $x \in V$ be an arbitrary vertex. By induction hypothesis, $G[V - x]$ has a Hamiltonian path $x_1, x_2, \ldots, x_{n-1}$. If there is an edge from $x$ to $x_1$, or from $x_{n-1}$ to $x$, then insert $x$ at the beginning or end of the ordering, respectively. Otherwise, there must be a $k$ with $1 \leq k < n - 1$, such that there is an edge from $x_k$ to $x$, and from $x$ to $x_{k+1}$ (start from $x_1$, and follow the path until a node has an edge from $x$). By inserting $x$ between $x_k$ and $x_{k+1}$, we obtain a Hamiltonian path for $G$.

Truchon [325] proves Proposition 5.1 in a different way, by showing that a Kemeny order also satisfies the XCC. A Kemeny order [205] is defined as follows: for each ordered pair $(x, y)$ of alternatives, let $w_{x,y}$ be the number of voters who prefer $x$ over $y$. A Kemeny order then minimizes the total sum of weights $w_{x,y}$ going “backwards”, i.e., going from $x$ to $y$ with $x > y$ in the ordering. We will prove this fact in Lemma 5.7 below.

### 5.1.1 Formalization of Social Choice Properties

So far, we have investigated several concrete approaches for determining a consensus ordering from several given orderings. All approaches suffered from “unnatural” outcomes in some cases or others. Perhaps, it
would thus make sense to axiomatize which properties a voting scheme should satisfy, and then look for voting schemes meeting these axioms. We describe here the axioms proposed by Arrow [17].

Formally, we will associate with each voter \( i \) a preference order \( \preceq_i \). A social choice function \( f \) takes all of these \( k \) orders, and outputs a consensus order

\[
\prec = f(\preceq_1, \ldots, \preceq_k).
\]

Such a function \( f \) should intuitively satisfy certain properties.

**Monotonicity** : If \( a \prec b \), and \( b \prec_i a \), then swapping \( a \) and \( b \) in \( \prec_i \) does not result in \( b \prec a \). Intuitively, this means that if \( a \) ranks above \( b \) overall, then changing another vote in \( a \)'s favor does not affect the relative ranking between \( a \) and \( b \).

**Non-triviality** : For each pair \( a \) and \( b \) of candidates, there is some choice of orderings \( \prec_i \) such that \( a \prec b \). This ensures that the relative ordering of \( a \) and \( b \) actually depends on the votes, and is not predetermined.

**Independence of Irrelevant Alternatives (IIA)** : Let \( \preceq_1, \ldots, \preceq_k \) and \( \prec'_1, \ldots, \prec'_k \) be two different preference orders for each voter, and \( \preceq = f(\preceq_1, \ldots, \preceq_k) \) and \( \prec' = f(\prec'_1, \ldots, \prec'_k) \) the corresponding consensus orderings. If \( B \subseteq V \) is a subset of candidates such that \( a \prec b \) if and only if \( a \prec'_i b \) for all \( a, b \in B \) (i.e., all \( \prec_i \) and the corresponding \( \prec'_i \) agree on the orderings of \( B \)), then \( a \prec b \) if and only if \( a \prec' b \) for all \( a, b \in B \). What this expresses is that if no voter changes his relative preference between any two candidates in \( B \), then the final ordering among candidates of \( B \) does not change. In other words, changing preferences merely with regards to third candidates does not affect the order of any two other candidates.

Monotonicity and non-triviality together imply the property of *unanimity*: if \( a \prec_i b \) for all \( i \), then \( a \prec b \). That is, if every voter prefers \( a \) over \( b \), then \( a \) ends up ahead of \( b \). To prove this fact, start with some set of orders \( \prec_i \) such that \( a \prec b \) (such a set exists by the non-triviality property). Then, we keep swapping the positions of \( a \) and \( b \) in all \( \prec_i \) that previously had \( b \prec a \). By monotonicity, the outcome will still be that \( a \) is ranked ahead of \( b \), and eventually, \( a \prec_i b \) for all \( i \).

In trying to find social choice functions satisfying all these axioms, one quickly notices that this is not so easy. In particular, the IIA property is not satisfied by many schemes. However, one class of social choice functions meeting these requirements is *dictatorship*: the dictator function \( f_i \) is defined as \( f_i(\preceq_1, \ldots, \preceq_k) = \preceq_i \). That is, the output is simply the preference of just one voter. Dictatorship is an undesirable quality for a voting scheme. That gives us our last property:

**Non-Dictatorship** : \( f \neq f_i \) for all \( i \). That is, the aggregation will not disregard the opinions of all but one voter.

Unfortunately, with this additional requirement, we have ruled out all remaining social choice functions:

**Theorem 5.3 (Arrow, 1951 [17])** There is no function \( f \) satisfying all the above four properties. Hence, the only functions satisfying the first three properties are the dictatorship functions.

**Proof.** To prove Arrow’s theorem, we show that any social choice function satisfying monotonicity, non-triviality, and IIA is in fact a dictatorship function. We will do this by first proving the existence of a single voter who can decide the order between two candidates; then, we prove that this voter is in fact a dictator.

First, we define sets that decide the outcome between two candidates. We call a set \( J \) of voters \( (A, B) \)-decisive if the fact that all of \( J \) ranks \( A \) ahead of \( B \) is enough to guarantee that \( A \) will be ranked ahead of \( B \) in the output. Formally, we write \( A \prec_J B \) to denote that \( A \prec_i B \) for all \( i \in J \). We then say that \( J \) is \( (A, B) \)-decisive if \( A \prec_J B \) implies \( A \prec B \). We call a set \( J \) of voters *decisive* iff \( J \) is \( (A, B) \)-decisive for some pair \( (A, B) \).

One useful characterization of decisiveness can be derived from the monotonicity property: \( J \) is \( (A, B) \)-decisive if and only if \( A \prec_J B \) and \( B \prec_J A \) imply that \( A \prec B \). We will use this characterization later. Notice
also that by the unanimity property, the set $V$ of all voters is always $(A, B)$-decisive for each pair $(A, B)$. Hence, there must be some smallest decisive set $J^*$. We will prove that $J^*$ is actually a singleton.

Let $\{A, B\}$ be the alternatives such that $J^*$ is $(A, B)$-decisive. Consider an arbitrary voter $i \in J^*$ from this set, and define $J^* = J^* - \{i\}$. Assume that the voter $i$ has preference order CAB, each voter in $J^*$ has order ABC, and each voter in $J^*$ has order BCA. (Other candidates can be ignored by the IIA property). In the final order, $A$ must precede $B$, because $J^*$ is $(A, B)$-decisive, and all of $J^*$ ranks $A$ ahead of $B$. In addition, we know that $C$ must precede $A$. The reason is that only $J'$ would prefer $A$ to precede $C$, so if it actually did, then $J'$ would be $(A, C)$-decisive, contradicting the assumption that $J^*$ is a smallest decisive set. Now, by transitivity, we conclude that the final output must be CAB, so $C$ precedes $B$. But $i$ is the only voter preferring $C$ to precede $B$, so by our characterization above, $\{i\}$ must be $(C, B)$-decisive. Because $\{i\}$ is thus decisive, it cannot be smaller than $J^*$, so we know that $J^* = \{i\}$, and $\{i\}$ is also $(A, B)$-decisive.

Next, we show that $\{i\}$ is also $(A, D)$-decisive for all $D \neq A$ and $(C, D)$-decisive for all $D \neq C$. We consider the scenario where voter $i$ has the order ABD, and everyone else has order BDA. Then, by unanimity, $B \prec D$, and because $\{i\}$ is $(A, B)$-decisive, we have $A \prec B$. By transitivity, the ordering is ABD. In particular, this means that $A$ precedes $D$, which only voter $i$ prefers. Thus, $\{i\}$ must be $(A, D)$-decisive. Replacing $A$ by $C$ in the previous proof shows that $\{i\}$ is also $(C, D)$-decisive. Also, by substituting $C$ resp. $A$ for $D$, we further conclude that $\{i\}$ must be $(C, A)$-decisive and $(A, C)$-decisive.

Next, we show that $\{i\}$ is also $(D, A)$-decisive for all $D \neq A$, as well as $(D, C)$-decisive for all $D \neq C$. Here, we assume that voter $i$ prefers the order DCA, and everyone else prefers ADC. By unanimity, we have $D \prec C$ in the outcome, and because $\{i\}$ is $(C, A)$-decisive, we have $C \prec A$. So the final order will be DCA. Again, $i$ is the only voter preferring $D$ over $A$, so $\{i\}$ must be $(D, A)$-decisive. Since $\{i\}$ is both $(A, C)$-decisive and $(C, A)$-decisive, we can simply switch the order $A$ and $C$ in the previous construction, and prove that $\{i\}$ is $(D, C)$-decisive as well.

As a final step, we show that $\{i\}$ is in fact $(D, E)$-decisive for all $D$, $E$, and hence a dictator. We assume that $i$ votes DAE, and everyone else votes EAD. Because $\{i\}$ is $(D, A)$-decisive and $(A, E)$-decisive, the final ordering must be DAE. But by the same argument as before, this means that $\{i\}$ is $(D, E)$-decisive.

In summary, we have proved that voter $i$ is a dictator, completing our proof of Arrow’s Theorem.

This result is quite disappointing. It suggests that for a very reasonable definition of what democratic decision-making is, the process is impossible. Yet, we observe frequently in practice that voting does work (reasonably) well. So one direction to pursue further is to ask: what kind of restrictions do voter preferences in practice seem to satisfy? Our construction in the proof was based on some very carefully crafted scenarios — perhaps, those don’t appear in practice.

This line of thought is pursued further in several papers. Black [46, 47] suggests single-peaked preferences on the line, i.e., each voter is located on a one-dimensional space (such as the political spectrum from “left” to “right”), and ranks preferences by distance from his own location. Barberà et al. [30] extend this to higher-dimensional spaces with the $L_1$-norm, and Richards et al. [300] extend it further to graph structures shared by all agents, in which agents have single-peaked preferences at nodes. In these cases, consensus voting is in fact possible, i.e., Arrow’s Axioms can be achieved.

### 5.2 Aggregation as Optimization

A different approach from the axiomatic one (which, in a sense, we just saw fail) would be to treat the determination of a consensus ordering as an optimization problem: to find a ranking that is as close as possible to the given set of rankings.

In order to phrase the problem this way, we first need to define a notion of distance between two rankings. As rankings are really nothing but permutations (so long as they are complete rankings — more about partial rankings later), we can look at known distance measures on permutations. Two well-known such measures are Spearman’s Footrule and Kendall’s $\tau$.

**Definition 5.4 (Spearman’s Footrule)** Let $\prec_1$, $\prec_2$ be two orderings, and $\prec_1 (a)$ the position of element
a in the ordering \( \preceq_1 \). Then, the Footrule distance is defined as

\[
F(\preceq_1, \preceq_2) = \sum_a |\preceq_1(a) - \preceq_2(a)|
\]

Notice that this is identical to the \( L_1 \) distance between the vectors of positions.

**Definition 5.5 (Kendall’s \( \tau \))** Kendall’s \( \tau \) counts the number of inversions (or Bubble Sort swaps) between the two permutations. Writing

\[
K_{a,b}(\preceq_1, \preceq_2) = \begin{cases} 
1 & \text{if } a \preceq_1 b \text{ and } b \preceq_2 a \\
0 & \text{otherwise},
\end{cases}
\]

we define

\[
\tau(\preceq_1, \preceq_2) = \sum_{a,b} K_{a,b}(\preceq_1, \preceq_2).
\]

The maximum value that can be attained by these metrics is \( n(n-1) \) for Spearman’s Footrule, and \( \frac{n(n-1)}{2} \) for Kendall’s \( \tau \). In both cases, the maximum is attained for two orderings that are the reverse of each other.

It is easy to see that the two metrics are not the same. One obvious example is when \( \preceq_1 = (1 \ 2) \), and \( \preceq_2 = (2 \ 1) \). Then, the Spearman Footrule distance is \( F = 2 \), while the Kendall distance is \( \tau = 1 \). While the two can be different, we can prove that they are not *very* different.

**Theorem 5.6 (Diaconis and Graham [114])** For any orderings \( \preceq_1 \) and \( \preceq_2 \), we have

\[
\tau(\preceq_1, \preceq_2) \leq F(\preceq_1, \preceq_2) \leq 2\tau(\preceq_1, \preceq_2).
\]

**Proof.** We first show that \( F(\preceq_1, \preceq_2) \leq 2\tau(\preceq_1, \preceq_2) \). We do this by induction on the value of \( \tau(\preceq_1, \preceq_2) \). In the base case, when \( \tau(\preceq_1, \preceq_2) = 0 \), both the orderings are the same, and hence \( F(\preceq_1, \preceq_2) = 0 \).

In the induction step, we look at \( \preceq_1 \) and \( \preceq_2 \) such that \( \tau(\preceq_1, \preceq_2) > 0 \). Let \( \preceq' \) be obtained from \( \preceq_2 \) by one switch towards \( \preceq_1 \). Then, \( \tau(\preceq_1, \preceq') = \tau(\preceq_1, \preceq_2) - 1 \), and \( \tau(\preceq', \preceq_2) = 1, F(\preceq', \preceq_2) = 2 \). By the Triangle Inequality, applied to the metric \( F \), we have that

\[
F(\preceq_1, \preceq_2) \leq F(\preceq_1, \preceq') + F(\preceq', \preceq_2) = F(\preceq_1, \preceq') + 2\tau(\preceq', \preceq_2).
\]

We apply the Induction Hypothesis to \( \preceq_1 \) and \( \preceq' \), obtaining that \( F(\preceq_1, \preceq') \leq 2\tau(\preceq_1, \preceq') \). Thus,

\[
F(\preceq_1, \preceq_2) \leq F(\preceq_1, \preceq') + 2\tau(\preceq', \preceq_2) \leq 2\tau(\preceq_1, \preceq') + 2\tau(\preceq', \preceq_2) = 2\tau(\preceq_1, \preceq_2),
\]

completing the inductive proof.

For the other inequality, \( \tau(\preceq_1, \preceq_2) \leq F(\preceq_1, \preceq_2) \), we use induction on \( F(\preceq_1, \preceq_2) \). In the base case \( F(\preceq_1, \preceq_2) = 0 \), both the orderings are the same, so \( \tau(\preceq_1, \preceq_2) = 0 \).

For the inductive step, we have the problem that simple switches will not necessarily improve the value of \( F \). As an example, we can look at \( \preceq_1 = (1 \ 2 \ 3 \ 4) \) and \( \preceq_2 = (4 \ 3 \ 2 \ 1) \). In this case, a switch towards \( \preceq_1 \) will result in the ordering \( \preceq' = (4 \ 3 \ 1 \ 2) \), which has \( F(\preceq_1, \preceq') = F(\preceq_1, \preceq_2) \). Thus, here we try to “meta-swap” two elements such that one is too far to the left of its position in \( \preceq_1 \) and the other is too far right of its position in \( \preceq_1 \).

Without loss of generality, we may assume that \( \preceq_1 = (12 \ldots n) \) (i.e., the elements are sorted), and \( \preceq_2 \) has element \( i \) in position \( a_i \). Thus \( F(\preceq_1, \preceq_2) = \sum_i |a_i - i| \).

Let \( i \) be maximal such that \( a_i \neq i \), i.e., the rightmost element that is out of position. Notice that this implies that \( a_i < i \), for otherwise, the element \( a_i \) (the one that is in position \( a_i \) in the ordering \( \preceq_1 \)) would be a larger index out of place. Let \( j \leq a_i \) be the largest index with \( a_j > a_i \) and \( a_i > j \), i.e., the rightmost element to the left of \( i \) in the order \( \preceq_2 \) which is too far right. (See Figure 5.1 for an illustration.) Notice
that such an element $j$ must exist by the Pigeonhole Principle: at most $a_i - 1$ of the elements $1, \ldots, a_i$ can be in positions $1, \ldots, a_i - 1$. Also notice that $a_j \leq i$, as otherwise, the element $a_j$, which is out of position, would show that $i$ is not maximal.

Let $\prec'$ be the ordering obtained by swapping $i$ and $j$ in the ordering $\prec_2$. Notice that in $\prec'$, neither $i$ nor $j$ “overshoot” their actual positions, because $j \leq a_i$ and $a_j \leq i$. Thus, both $i$ and $j$ move $|a_i - a_j|$ positions closer to their destination, and all other elements stay in position. We obtain that $F(\prec_1, \prec_2) = 2|a_i - a_j| + F(\prec_1, \prec')$. Applying the Induction Hypothesis, we thus have that $F(\prec_1, \prec_2) \geq 2|a_i - a_j| + \tau(\prec_1, \prec').$

We can also calculate the Kendall $\tau$ distance between $\prec_2$ and $\prec'$ quite easily. By making $|a_i - a_j|$ switches to the right of element $i$, we get it into position. Then, another $|a_i - a_j| - 1$ switches to the left of element $j$ move it to position $a_i$. All elements in between are moved once to the left and once to the right, and thus end up in the same position. Hence, we just proved that $\tau(\prec', \prec_2) \leq 2|a_i - a_j| - 1$. Now, by the Triangle Inequality for the metric $\tau$, we obtain that

$$\tau(\prec_1, \prec_2) \leq \tau(\prec_1, \prec') + \tau(\prec', \prec_2) \leq \tau(\prec_1, \prec') + 2|a_i - a_j| - 1 < F(\prec_1, \prec_2).$$

This completes the inductive proof. \hfill \qed

Now that we have metrics to measure the difference between two orderings, given a list of several orderings $\prec_1, \prec_2, \ldots, \prec_k$, we want to find an ordering $\prec$ “close to all” of them. There are multiple concrete optimization criteria we could be trying to minimize, for example:

1. The average $\tau$ distance $\frac{1}{k} \sum_i \tau(\prec, \prec_i)$.
2. The average Footrule distance $\frac{1}{k} \sum_i F(\prec, \prec_i)$.
3. The maximum $\tau$ distance $\max_i \tau(\prec, \prec_i)$.
4. The maximum Footrule distance $\max_i F(\prec, \prec_i)$.

Notice that the first definition exactly defines Kemeny orderings, as discussed in the previous section. For the sum of all Kendall $\tau$ distances is exactly the total number of disagreements between the ordering $\prec$ and all other orderings. Next, we show that a Kemeny ordering satisfies the extended Condorcet property.

**Lemma 5.7** If $\prec$ minimizes $\sum_i \tau(\prec, \prec_i)$, then it satisfies the extended Condorcet property.

**Proof.** We prove the lemma by contradiction. Assume that we have an ordering $\prec$ which minimizes $\sum_i \tau(\prec, \prec_i)$, but does not satisfy the extended Condorcet property. Then, there is a partition $(S, \overline{S})$ of the alternatives $\{1, \ldots, n\}$ such that every alternative in $S$ beats all those in $\overline{S}$, yet $\prec$ ranks some $j \in \overline{S}$ ahead of some $i \in S$. Therefore, there must be such a pair such that $i$ and $j$ are adjacent in the ordering $\prec$. Now, swapping $i$ and $j$ improves $\sum_i \tau(\prec, \prec_i)$ by $i$’s margin of victory. Since the relative orderings of no other elements are affected, the objective function actually improves, contradicting the optimality of $\prec$. \hfill \Box

\footnote{Recall that $i$ is said to beat $j$ if there are more orderings in which $i$ precedes $j$ than vice versa.}
Unluckily, minimizing the average $\tau$ distance is NP-hard, as proved by Dwork et al. [124] (we don’t give the proof here).

**Fact 5.8** If $k \geq 4$, then minimizing $\sum_i \tau(\prec, \prec_i)$ is NP-hard.

Instead, we may focus on the second objective function: minimizing $\sum_i F(\prec, \prec_i)$. This one can actually be minimized in polynomial time. The intuition is to look at an ordering/permutation as a matching between elements and positions. Then, we can express the total objective function value as a sum of “penalties” incurred by each element for the position it is in. Specifically, we assign a penalty $\phi_{j,p} = \sum_i |p - \prec_i(j)|$ for putting element $j$ in position $p$. (Recall that $\prec_i(j)$ denotes the position in which element $j$ appears in the order $\prec_i$). For each $j$, we thus need to assign a (distinct) $p(j)$ to minimize $\sum_j \phi_{j,p(j)}$. Since feasible assignments are exactly perfect matchings between elements and positions, our goal is to find the cheapest perfect matching in the complete bipartite graph where the edge $(j, p)$ has cost $\phi_{j,p}$. Finding minimum-cost perfect matchings is known to be solvable in polynomial time, for instance by first computing any perfect matching, and then repeatedly eliminating negative cycles in the residual graph (see [8, 222] for discussions of polynomial algorithms for minimum cost matching).

Notice that the ordering minimizing the average Footrule distance is also a 2-approximation to the problem of minimizing the average Kendall’s $\tau$ distance. For if $\text{OPT}$ denotes the best Kendall ordering, and $\prec$ the best Footrule ordering, then

$$\sum_i \tau(\prec, \prec_i) \leq \sum_i F(\prec, \prec_i) \leq \sum_i F(\text{OPT}, \prec_i) \leq 2 \sum_i \tau(\text{OPT}, \prec_i).$$

### 5.2.1 Partial Orderings

In the previous discussion, we assumed that the orderings were actually complete. What happens when we are not given full orderings? Obviously, this is relevant in combining search engine results, since different search engines will have different web crawls, containing different sets of pages in the search results. What techniques can we then use to compare these orderings? Fagin et al. [135] study the issue of how to compare top $k$ lists. If the lists don’t all contain all of the elements, one needs to extend (or replace) the distance notions defined above.

Fagin et al. suggest several techniques to deal with this. One is to augment all the lists such that all the elements appear in all the lists, by appending the missing elements at the end of each list (since they were clearly not considered to be in the top $k$ by that list). This raises the question in which order the extra elements should be appended to the lists. An “optimistic” view would define the distance between lists based on the assumption that the missing elements appear in the same relative order as in the other list. Another solution is to append the elements in a random order, and define the distance as the average.

More generally, we can define an extension of Kendall’s $\tau$ with a penalty $p$ as follows. If both elements $i, j$ are in one list, and at least one element is in the other list, then the transposition penalty for $i, j$ is just the same as for Kendall’s $\tau$, i.e., 0 or 1 depending on whether their order is the same or different. (If only $i$ is in the second list, then the penalty is 0 if $i$ appears before $j$ in the first list, and 1 otherwise.) If one list contains both $i$ and $j$, and the other list contains neither, then the penalty is $p$. Then, $p = 0$ corresponds to the optimistic assumption, and $p = \frac{1}{2}$ to the random ordering of absent elements.

[135] shows that neither of these approaches for pairwise comparison of lists satisfies the triangle inequality, and hence, we do not obtain a metric. However, they show that the distance measures can be both upper and lower bounded by a metric.

Spearman’s Footrule can also be generalized by introducing a new location $\ell$, and assuming that all missing elements are at location $\ell$, and calculating the normal Footrule distance. Notice that this does not really define a permutation any more. The most obvious choice would be $\ell = k + 1$, but other choices are possible. The resulting Footrule distance actually does define a metric.
5.3 Further Reading

Even though this chapter is motivated mostly by rank aggregation for search engines and similar applications, our initial example focused on elections. There, the goal is frequently different: rather than finding an aggregate ordering of all alternatives, the goal is solely to determine one winner. While we could design any number of simple rules (e.g., counting votes, runoff voting, etc.), there is now the issue that voters may reason that incorrectly stating their preferences might help them achieve a more desirable outcome. For instance, in our introductory example, knowing the distribution of other votes, the third type of voters (with the order RN-AG-GWB) might now declare AG as their first choice to avoid their least favorite outcome.

Hence, it becomes important to study which types of voting systems are strategy-proof against voters misrepresenting their preferences. The Gibbard-Satterthwaite Theorem [158, 306] gives an equally pessimistic characterization to Arrow’s Theorem. It posits the following axioms:

1. **Fairness**: Each candidate can win if all voters unanimously prefer this candidate to all others.

2. **Strategy-Proofness**: No voter is ever better off misrepresenting his preferences.

The only voting rules satisfying these axioms (for three or more candidates) are again the dictator functions.

Again, an interesting question is what natural restrictions on the preferences would lead to non-trivial social choice functions. In this context, the survey paper of Barberà [29] gives a nice summary of results. Again, single-peaked preferences on the line and in higher dimensions with the $L_1$-norm (e.g., [267]) lead to non-trivial mechanisms that have voters truthfully revealing their preferences.

When we allow randomized social choice functions, the situation is not quite as dire as the one of the Gibbard-Satterthwaite Theorem. Gibbard [159] shows that a rule is strategy-proof only if it is a randomization over dictator functions and two-choice rules (i.e., rules that prune all but two candidates, and then decide on a winner among these two candidates by considering the votes). Conitzer [97] extends this result to the case when the voting rule also needs to be anonymity-proof, in the sense that no agent benefits by casting multiple votes. (This is necessary, for instance, for online polls, where the inherently anonymous environment cannot prevent a user from voting multiple times.) Conitzer shows that in this case, the remaining randomized social choice functions are significantly more restricted: the voting rule randomizes between picking a uniformly random alternative, or picking two alternatives, and then checking if all agents prefer one over the other. If so, that alternative is the winner; if not, then a fair coin is flipped. Thus, the outcome is — except for a few unanimous choices — essentially always random.

Dictator functions also play an important role in the Fourier Analysis of Boolean Functions [283]. Using the setup of Fourier Analysis, and a proof based on uniformly random orderings by all voters, Kalai [198] gives a very short proof of Arrow’s Theorem. Kalai’s survey [199] describes the connections between Fourier Analysis and Social Choice Theory in more depth, and relates them to the outcomes of elections with random choices.

Indeed, the impact of randomization already impacted Condorcet’s description. He posits a model where each voter chooses between two alternatives, and makes the societally most beneficial choice with some probability $p > \frac{1}{2}$. Condorcet [108] then shows that the majority vote has the highest probability of identifying the right choice. This is a direct precursor to analysis showing that the Kemeny order, under a natural model of randomization, is most likely to be the societally preferred one. This is discussed in more detail in papers by Young [343] and Truchon [325].

Finding a Kemeny ordering is equivalent to the well-studied problem of finding a Minimum Feedback Arc Set. A set $S$ of edges in a directed graph $G$ is a Feedback Arc Set if and only if its removal breaks all cycles in $G$. Clearly, once all cycles are broken, the nodes can be ordered by any topological sort order, respecting all preferences. Thus, finding the smallest number of edges to remove in a (multi-)graph is equivalent to finding the ordering minimizing the total number of inversions.

In general, the best known approximation for Feedback Arc Set is $O(\log n \log \log n)$ [131]. However, for the case of rank aggregation, we know that the graph is a tournament graph, where an edge is directed only from the candidate who wins the pairwise comparison to the loser of that comparison. The edge weight is the difference in votes between the two relative orders. For this special case, Kenyon-Mathieu and Schudy...
provide a PTAS [212], improving on a previous (simpler) 3-approximation algorithm due to Ailon, Charikar, and Newman [11]. Ailon [10] shows how to extend these algorithms to the case when the rankings are only partial (as discussed in Section 5.2.1), and gives a 2-approximation, and a more complex $\frac{3}{2}$-approximation algorithm.

An issue somewhat similar in spirit to aggregating top-$k$ rankings is the aggregation of rankings with ties. Fagin et al. [134] give extensions of the Kendall $\tau$ and Spearman Footrule measures for those cases, similar to the ones discussed for top-$k$ lists in Section 5.2.1. Again, they generalize the result of Diaconis and Graham to show equivalence of the measures to within constant factors, and discuss the issue of aggregation.

A different approach to ranking is proposed by Hochbaum [181]. She posits that not only are we given pairwise comparisons, but for each pair $(i, j)$ also a weight $w_{i,j}$ indicating how much one beat the other. These weights may be inconsistent, e.g., A beat B, B beat C, and C beat A. The goal is to “smooth” them to make sure that all numbers end up consistent, in the sense that for each pair of alternatives $(i, j)$, all paths from $i$ to $j$ have the same sum of weights in the end. At the same time, the smoothed weights $w'_{i,j}$ should be “close” to the original ones. [181] shows that so long as the penalty function for deviations is convex, an optimal solution can be found in polynomial time using the dual of a minimum cost network flow problem [7].

The issue of learning orderings is addressed in more depth by Cohen et al. [93]. The idea is to use an expert learning [240, 241, 150] approach to finding scores for pairwise comparisons, and then find the ordering that agrees best with these scores. The learning approach is also pursued by Joachims [196], who uses Support Vector Machines to learn rankings from user preferences that are expressed in clickthrough data.
Chapter 6

Power-Law Distributions

In Section 1.2.3, we briefly remarked on the fact that the WWW graph appears to have a power-law degree distribution, i.e., the fact that the frequency of nodes of degree $d$ decays polynomially as $d^{-\gamma}$ for some constant $\gamma$ (as opposed to — say — exponentially). This observation clearly invalidated the Erdős-Rényi $G(n, p)$ model as a model for the WWW. In this chapter, we will investigate power-law distributions in much more detail, with a focus on generational models which would predict power laws. Most of the material in this chapter, and much more, is covered in the excellent survey of the topic by Mitzenmacher [257].

The existence of power law distributions has been observed empirically over the years in many natural and man-made scenarios. Both the observation of such distributions and possible generational models have received quite a lot of attention. A major surge in interest within the computer science (and physics) community resulted from the paper by Faloutsos, Faloutsos, and Faloutsos [137], which demonstrated power laws in a crawl of the Internet and WWW graphs. Some of the other contexts in which power laws had been observed previously include:

- Financial Models: price changes of securities in financial markets [245].
- Biology: lengths of protein sequences in genomes [193] or variation among plant species [314].
- Linguistics: word length frequencies [130, 347] and degrees of words in association networks [318].
- City populations or distribution of income among people [285].
- Number of papers published by scientists [271].

Formally, we define power law distributions as follows:

**Definition 6.1** A non-negative random variable $X$ is said to have a power law distribution if $\text{Prob}[X \geq x] = cx^{-\alpha}$, for some constants $c, \alpha > 0$.

If the variable is actually discrete, then the definition implies that $\text{Prob}[X = x] = c'x^{-\alpha'}$, for different values $c', \alpha' > 0$. We will use the definitions interchangeably.

Power law distributions fall into the class of heavy tailed distributions: the probability that $X$ assumes a large value is only polynomially small, compared to the exponentially small probabilities for Gaussian, Binomial, or other common distributions.

When given an actual set of data, we can recognize it as a power law most easily in a log-log plot, i.e., in a plot where both axes scale logarithmically. For if $f(x) = c \cdot x^{-\alpha}$ denotes the frequency with which value $x$ was observed, then $\log f(x) = \log(c) - \alpha \cdot \log(x)$. Hence, in a log-log plot, we will observe a straight line with a slope of $-\alpha$ and $y$-intercept of $\log(c)$. However, two important caveats should be noted:

1. In order to infer with confidence that a power law is indeed present, the straight line should extend over several orders of magnitude at least. A very short straight-line segment in a log-log plot does not indicate a power law very conclusively.
2. To determine the coefficient $\alpha$, it is usually not right to fit the best straight line to the data, e.g., via regression, since the log-log scaling of the plot invalidates the stochastic assumptions behind the regression. A more in-depth discussion of this issue is given by Clauset et al. [91].

6.0.1 Power Laws in the WWW

For the WWW, the paper of Broder et al. [65] was the first to show that the distribution of in (resp. out) degree $d$ is proportional to $d^{-\alpha}$ for some value $\alpha \in [2.1, 2.2]$ (resp. $\alpha \in [2.7, 2.8]$). Given the values of $\alpha, c$, we can calculate the mean of the power law distribution as

$$\mu = \sum_d d \cdot c \cdot d^{-\alpha} = c \sum_d d^{1-\alpha}.$$ 

Notice that the mean is finite iff $\alpha > 2$. In particular, this implies that if the power law indegree distribution were to continue to hold as the WWW grows, the average indegree of pages would remain finite. The fact that the mean indegree of the WWW is finite is not too surprising, given that the average indegree equals the average outdegree, and we don’t expect the average web page to have more than a constant number of outlinks, as each requires actual work.

However, there are natural examples of networks that exhibit power laws with $\alpha < 2$. For such networks, we would predict that if the power law persists as the size of the network grows, the mean would diverge to infinity:

• The WWW at a site level ($\alpha \approx 1.6$) [40]. Here, single sites can include large numbers of web pages, and we may expect the number of pages within a site to grow along with the web.

• Co-authorship in high-energy physics ([271], $\alpha \approx 1.2$). High-energy physics papers often have very large numbers of co-authors (in excess of 1000). It is not clear how to predict the scaling of this graph in the future, but it is not inconceivable that future papers may have even larger author lists. (Notice that by comparison, the mathematics co-authorship graph has a rather large value of $\alpha \approx 2.5$. Many mathematics papers still have a single author.)

6.1 Preferential Attachment

In trying to explain observed power laws, many models posit a “rich get richer” phenomenon: entities (such as nodes or web pages, dollars, plant genera, cities, individuals, etc.) that already have a large value (degree, number of species, populations, wealth, etc.) have a tendency to attract more of the same. If this attraction is linear in the current value, then power laws emerge.

In the case of the WWW or other graphs, such a behavior is called Preferential Attachment. It is posited that newly arriving nodes link to an existing node $v$ with probability proportional to $v$’s current degree. Thus, high-degree nodes are more likely to attract new links. The underlying assumption here is not necessarily that nodes make this choice deliberately, but rather that high-degree nodes (e.g., web pages with high indegree) are more likely to be discovered in the first place, and thus linked to. Kumar et al. [228] make this explicit by investigating a copying model, wherein newly arriving nodes randomly select another node, and copy some of that node’s outlinks. A mathematically rigorous analysis of such models tends to be quite involved (and is carried out, for instance, in [228]). However, by making some non-rigorous simplifications, we can obtain qualitatively identical results.

Here, we study the following simple model: The graph starts with zero nodes. At each time $t \in \mathbb{N}$, a new node arrives and generates $k$ outlinks; we will label the node with its arrival time $t$. Each outgoing link is either uniformly random, which happens with probability $1 - \alpha$, or preferential with probability $\alpha$. In the latter case, the edge links to existing nodes with probabilities proportional to their degrees.

Our analysis will follow the outline by Barabási et al. [27, 28]. Let $d_i(t)$ denote the in-degree of node $i$ at time $t$. Then, because node $i$ arrives at time $i$ with no links yet, we have that $d_i(i) = 0$. At any time $t$, the total number of inlinks (into all nodes) is $kt$, and the total number of nodes is $t$. Hence, the probability
Hence, after \( t \) at most \( d \) each node higher degree, as they are around for more rounds, and can thus receive more incoming edges. Notice that nodes never attach preferentially, but always choose uniformly at random. Older nodes will still end up with respect to distributions. While “rich-get-richer” models may be considered close to the truth for WWW-like graphs, in the last section, we investigated the preferential attachment model for the generation of power law degree structure (such as triangles or other unusually dense subgraphs). The copying model does slightly better in this respect.

To obtain the density function from this, we take the derivative with respect to \( d \), which gives us density \( \frac{1}{\alpha} \ln(\alpha d_i(t) + \beta) = \ln(t) + c \). Solving for \( d_i(t) \) now shows that \( d_i(t) = \frac{\alpha d_i(t) + \beta}{\alpha} \). We still have to find out the value of the constant \( c \). To determine it, we can use the initial condition that \( d_i(i) = 0 \). This gives us that \( \frac{c e^{\alpha d_i(t)}}{\alpha} = 0 \), which by rearranging yields that \( e^{\alpha c} = \beta \cdot i^{-\alpha} \). Substituting this back into the expression for \( d_i(t) \) now shows that

\[
\frac{\beta}{\alpha} \left( \frac{(t/i)^\alpha}{1} - 1 \right).
\]

To find the cumulative function for the number of nodes with degree greater than or equal to \( d \), we first solve the inequality \( d_i(t) \leq d \). This yields that the expected degree is at most \( d \) whenever \( t \geq (d \cdot \frac{2}{\beta} + 1)^{-\frac{1}{\alpha}} \). Thus, the fraction of nodes with degree greater than \( d \) at time \( t \) is

\[
\frac{t - i}{t} = \frac{t - t \cdot (d \cdot \frac{2}{\beta} + 1)^{-\frac{1}{\alpha}}}{t} = 1 - (d \cdot \frac{2}{\beta} + 1)^{-\frac{1}{\alpha}}.
\]

To obtain the density function from this, we take the derivative with respect to \( d \), which gives us density

\[
\frac{1}{\beta} \cdot \left( d \cdot \frac{2}{\beta} + 1 \right)^{-\left(1 + \frac{1}{\alpha}\right)}.
\]

Thus, we observe a power law with exponent \( 1 + \frac{1}{\alpha} > 2 \). In particular, for \( \alpha \approx 0.9 \), we obtain the same degree distribution as for the WWW.

Notice that for \( \alpha = 0 \), the above analysis breaks down (and the result is meaningless). Indeed, for \( \alpha = 0 \), nodes never attach preferentially, but always choose uniformly at random. Older nodes will still end up with higher degree, as they are around for more rounds, and can thus receive more incoming edges. Notice that each node \( i \), in each round \( t \), receives an expected \( k/t \) new incoming links (as there are \( t \) competing nodes). Hence, after \( t \) rounds, node \( i \) will have expected indegree \( \sum_{j=i+1}^{t} \frac{k}{j} \approx k \cdot \log(t/i) \). The indices \( i \) having degree at most \( d \) are thus \( i \geq t \cdot e^{-d/k} \), and the fraction with degree at least \( d \) is \( 1 - e^{-d/k} \). Taking a derivative with respect to \( d \) gives us a density function of \( \frac{1}{\beta} \cdot e^{-d/k} \), which is sharply concentrated. In particular, we will not obtain a power law. So the age advantage alone does not explain the power law distribution derived above; rather, the preferential attachment was a crucial part of the model.

Notice that the preferential attachment model can be easily extended to include deletion of nodes and edges and rewiring etc., and the same type of analysis based on differential equations can be carried out. Usually, the power law degree distribution is preserved under these modifications, although they usually result in a large number of parameters, whose relative sizes affect the exact value of the exponent.

While preferential attachment gives us the same degree distribution as was observed for the WWW, it fails to obtain several other key properties (among others, all graphs generated are acyclic). For instance, as the links are generated independently at random, the graph will likely not exhibit much community structure (such as triangles or other unusually dense subgraphs). The copying model does slightly better in this respect.

### 6.2 Heavy-Tailed Distributions via Optimization

In the last section, we investigated the preferential attachment model for the generation of power law degree distributions. While “rich-get-richer” models may be considered close to the truth for WWW-like graphs,
they do not seem appropriate for graphs such as the Internet, as there appears no natural reason why nodes would choose to attach to high-degree nodes (or be more likely to find out about them). Hence, this class of models is not very useful for explaining the power laws observed in the Internet at the AS level [137].

Fabrikant et al. [133] argue that the reason power laws evolve in the Internet is a heuristic optimization performed by the owners of machines. Indeed, optimization as a cause for power laws has been investigated before. Mandelbrot [244] and Zipf [348] argue that power laws in word lengths are caused by the “Principle of Least Effort”: languages evolve so as to optimize the average information transmitted per character or phoneme. (Notice, however, that the same result about word lengths can be obtained by assuming that characters, including the space bar, are pressed completely randomly [256], as we see in Section 6.3 below.) Carlson and Doyle [71] extend the argument for file sizes and other parameters, and Fabrikant et al. [133] apply it to Internet-like graphs.

They posit the following graph growth model. A communication tree is built as nodes arrive uniformly at random in a unit square. Let $O$ be the first node that arrives, and assume that it arrives at the center of the square. Each node $i$ arriving subsequently connects to a node $j$ that had arrived earlier. The issue is which node $j$ should $i$ connect to. [133] argues that nodes want to be central in the network, i.e., few hops from $O$. At the same time, they want to have low “last mile” cost for their connection to $j$. The tradeoff is accomplished by considering the objective function $d_{ij} + \beta h_j$ where $d_{ij}$ is the Euclidean distance between nodes $i$ and $j$, $h_j$ is the number of hops from $j$ to $O$ in the communication tree, and $\beta$ is a given constant. That is, node $i$ connects to the node $j$ minimizing $d_{ij} + \beta h_j$, and consequently has $h_i = h_j + 1$.

Depending on the value of the parameter $\beta$, the graph evolves in very different ways.

- If $\beta$ is large (e.g., $\beta \geq \frac{1}{\sqrt{2}}$), then the hop count is more important than any possible distance (all distances to $O$ are at most $\frac{1}{\sqrt{2}}$), so the graph will evolve to be a star with node $O$ as the center. In particular, the degree distribution will not be heavy-tailed, since all nodes but one have degree 1.

- If $\beta = 0$, then the Euclidean distances become the only criterion for minimizing the objective function. Thus, nodes connect to their closest neighbors. We will analyze this process briefly.

If node $i$ has two neighbors $j, j'$, such that $d_{ij}, d_{ij'} \geq r$, then we also have that $d_{jj'} \geq r$ (else $j, j'$ would have been neighbors instead of both connecting to $i$). It follows that each neighbor $j$ of $i$ at distance $r$ has a circle of influence with area at least $\Omega(r^2)$ around it, which does not contain any other neighbors of $i$.

Now consider the $O(\log n)$ rings around $i$ of the form $R_k := \{ j \mid d_{ij} \leq 2^{-k} \}, k = 0, \ldots, 3\log n$. Note that with high probability, each node other than $i$ lies inside one of the $R_k$. Within each $R_k$, node $i$ can have at most a constant number of neighbors, because the area of each $R_k$ is at most $\Omega(2^{-2k})$, while each neighbor of $i$ in $R_k$ has an area of at least $\Omega(2^{-2(k+1)})$ in which no other neighbor can lie. Hence, there is at most a constant number of neighbors in each such ring, and a total of at most $O(\log n)$.

Thus, all degrees are bounded by $O(\log n)$, and the distribution of degrees is not power law, or even heavy-tailed.

- Thus, we next consider the case when $\beta > 0$ is small enough.

**Theorem 6.2** There is some $\varepsilon > 0$ (depending on $\beta$), such that there are at least $\Omega(n^{\varepsilon/6})$ nodes of degree at least $n^{1-\varepsilon}$.

Notice that while this proves that the distribution is heavy-tailed, it does not necessarily imply a power law. Indeed, Berger et al. [36] subsequently showed that the distribution is not power-law. It has a heavy tail, but does not exhibit all intermediate degrees to the corresponding extent.

**Proof.** We first show that each neighbor of the first node $O$ carves out a region of influence so that all nodes landing in the region must link to it. This gives those nodes high degrees. Then, we show that there will be enough such neighbors of $O$. Let $S$ denote the set of all neighbors of $O$. 

62
Consider a node $i$ that connects to $O$, and has $d_{i,O} \in [\beta, 2\beta]$. Let $r_i = d_{i,O} - \beta$. At the time node $i$ arrived, there are no other nodes from $S$ within distance $r_i$ of $i$, because $i$ linked to the root (otherwise, it would connect to such a node). No node $j$ arriving after $i$ within distance $d_{j,i} \leq r_i/2$ connects to node $O$, because it would be cheaper to connect to $i$. So there are no other nodes of $S$ within distance $r_i/2$ of $i$. Every node $j$ that lands within $r_i/4$ of $i$ links to $i$, because the cost to connect to node $i$ is at most $\beta + r_i/4$, while the cost to connect to another node $i'$ within distance $d_{i',i} \leq r_i/4$ of $i$ is at least $2\beta$, and the cost to connect to a node $i'$ with distance $d_{i',i} > r_i/4$ of $i$ is at least $\beta + d_{i',i} \geq \beta + r_i/4$ by the triangle inequality.

So node $i$ has a region of influence with area $\Omega(r_i^2)$ around it in which all nodes link to it, and thus receives at least an $r_i^2$ fraction of all edges.\footnote{Note: to show that $i$ will have degree $\Omega(r_i^2 n)$, we need to be sure that enough other nodes will arrive after $i$ to give it high degree. We can solve this problem by only considering the degrees for nodes that are among the first $n/2$ to arrive. Then, at least $n/2$ nodes arrive subsequently, giving us the desired high degree for those nodes. We only lose a factor of 2.}

So we know that any such node $i$ has high degree (as long as $r_i$ is large enough). We still need to show that there are enough such nodes $i$. Among others, we want to ensure that not too many nodes will link to $i$. To address this concern, consider a node $j$ arriving after $i$. The line of indifference between connecting to $i$ and connecting to $O$ is given by $d_{i,O} = \beta + d_{j,i}$. This curve is a hyperbola (see Figure 6.1(a)), i.e., the nodes linking to $i$ are a subset of the interior of that hyperbola.

Consider the following rings (Figure 6.1(b)) around $O$: Ring I has inner radius $\beta$ and outer radius $\beta + n^{-\varepsilon/2}$, Ring II has inner radius $\beta + n^{-\varepsilon/2}$ and outer radius $\beta + \frac{1}{2}n^{-\varepsilon/3}$, and Ring III has inner radius $\beta + \frac{1}{2}n^{-\varepsilon/3}$ and outer radius $\beta + n^{-\varepsilon/3}$ ($\varepsilon$ is chosen small enough that $\beta + n^{-\varepsilon/3} \leq 2\beta$.) Note that the area of Ring II is a constant fraction of the total area of Rings I, II, and III.

No node will ever link to a node $j \in S$ with $d_{j,O} < \beta$, because by triangle inequality, it would be cheaper to link to $O$ directly. Nor will any node that falls in these rings ever connect to any node more than one hop from node $O$, because the cost would be at least $2\beta$, so it would be cheaper to connect to node $O$ directly. Similarly, nodes in Ring II will not connect to a node $j$ outside the outer circle, because the cost is at least $\beta + \frac{1}{2}n^{-\varepsilon/3}$, while a direct connection to node $O$ would cost at most $\beta + \frac{1}{2}n^{-\varepsilon/3}$.

In summary, each node $i$ with $\beta + n^{-\varepsilon/2} < d_{i,O} < \beta + \frac{1}{2}n^{-\varepsilon/3}$ links either to $O$ or to some node $j \in S$ with $\beta \leq d_{j,O} \leq \beta + n^{-\varepsilon/3}$. Each such node that links to $O$ carves out a hyperbola, so that subsequent nodes that fall in the hyperbola do not link to node $O$.\footnote{Note: to show that $i$ will have degree $\Omega(r_i^2 n)$, we need to be sure that enough other nodes will arrive after $i$ to give it high degree. We can solve this problem by only considering the degrees for nodes that are among the first $n/2$ to arrive. Then, at least $n/2$ nodes arrive subsequently, giving us the desired high degree for those nodes. We only lose a factor of 2.}
If a new node \( i \) from Ring II arrives outside all existing hyperbolas of nodes \( j \) with \( \beta \leq d_j, O \leq 2\beta \), then \( i \) links to node \( O \). We want to find a lower bound on the number of hyperbolas defined by such \( i \). Some analytic geometry shows that the largest angle \( \theta \) of each hyperbola is \( O \left( \sqrt{r_i / \beta} \right) \). Because \( r_i < n^{-\varepsilon / 3} \) in Ring II, and \( \beta \) is a constant, \( \theta \) is \( O \left( \sqrt{n^{-\varepsilon / 3}} \right) = O(n^{-\varepsilon / 6}) \). So there is room for \( \Omega(n^{\varepsilon / 6}) \) disjoint hyperbolas.

Each node \( i \) that carves out a hyperbola has degree at least \( \Omega(n^{r_i^2}) \). Because \( r_i = \Omega(n^{-\varepsilon / 2}) \) in Ring II, the degree of these nodes is at least \( \Omega(n \cdot (n^{-\varepsilon / 2})^2) = \Omega(n^{1-\varepsilon}) \). Among all the nodes arriving in any of the rings (and thus claiming hyperbolas), at most a constant fraction lie outside of Ring II, so at least \( \Omega(n^{\varepsilon / 6}) \) nodes will be in Ring II, claim hyperbolas, and thus have degree at least \( \Omega(n^{1-\varepsilon}) \).

Notice that while we only talk about the expected number of such nodes (and their expected degree), standard occupancy bounds can be used to show concentration, i.e., that the actual outcome will be close to the expectation.

This completes the proof.

6.3 A Simple Model for Power Laws in Word Frequencies

While Mandelbrot [244] and Zipf [348] argue that power laws in word frequencies can be explained by optimization of transmitted content, Miller [256] proposes a much simpler model also resulting in power law distributions.

Miller posits a completely random typewriter (e.g., a monkey at a typewriter). The typewriter has \( b \geq 1 \) letter keys, each of which is hit with probability \( p/b \), for some constant \( p \in (0, 1) \). Word separators (such as newline or space) are hit with the remaining probability \( q = 1 - p \).

In this model, the probability that any particular word of \( i \) letters is typed is \( (p/b)^i \cdot q \), and all length-\( i \) words have the same probability. Thus, in this model, each length-\( i \) word is more likely than any word of length \((i + 1)\). Since there are \( b^i \) words of length \( i \), they occupy the next \( b^i \) positions (in word frequency order) starting from position \( k = \sum_{j<i} b^j = \frac{b^i - 1}{b - 1} = \Theta(b^{i-1}) \). Thus, the words at positions \( k \approx b^{i-1}, \ldots, b^{i} \) (roughly) each have probability \( (p/b)^i \cdot q \approx p^i \cdot q \cdot k^{-1} \). Taking logs with base \( b \), we have that the log of the probability is roughly \(-\log_b k + i \log_b p + \log_b q\), which, since \( i \approx \log_b k \), is \( \log_b k \cdot (\log_b p - 1) + \log_b q \). We thus obtain a power law with exponent \( \log_b p - 1 \).

The model of “monkeys at typewriters” has been extended by Conrad and Mitzenmacher [98] to the case of non-equal probabilities for hitting different keys. They show — using a significantly more complex proof — that power laws in the word frequency distribution persist even in this case.

We can also interpret the same mathematical model differently. Suppose that we have a hierarchy of topics, organized as an (infinite) \( b \)-ary tree. Topics higher up in the tree are more general, while topics further down are more specific. We traverse this tree starting from the root; at each step, we stop with probability \( q \), and otherwise choose a uniformly random child. Thus, the distribution of links to a particular topic is power law with exponent \( \log_b p - 1 \).

We can extend this to a simple model for topic-based linking in the WWW. Let the tree be truncated at a finite level now. For each of the topics (i.e., nodes of the tree), there are \( c \) pages covering that topic. When a page is created, it generates \( k \) outgoing links. Each outgoing link goes to a page on the same topic, or to a page on a more general topic, i.e., an ancestor in the tree. Say that a page links to another page \( h \) levels up with probability proportional to \( \gamma^h \), for some constant \( \gamma \). We calculate the in-degree distribution of a page \( p \). Suppose that \( p \) is located \( h \) levels above the leaves. Thus, there are \( b^h \cdot c \) pages in the descendants \( \ell \) levels below \( p \). Each gives rise to an expected \( k/c \cdot \gamma^\ell \) links to \( p \), so the expected number of links into \( p \) is \( \sum_{\ell=0}^{h} b^\ell \cdot c \cdot \gamma^\ell \cdot k/c = k \cdot \sum_{\ell=0}^{h} (b\gamma)^\ell = \frac{k}{b\gamma - 1} ((b\gamma)^{h+1} - 1) \).

Thus, if \( b\gamma < 1 \), the nodes’ degrees are nearly uniform, because \( \frac{(b\gamma)^{h+1} - 1}{b\gamma} \) is bounded on both sides by some constant. On the other hand, when \( b\gamma > 1 \), most of the contribution to any node comes from the leaves (since there are many of them, and they have high enough probability to link up). The nodes of degree at least \( \frac{k}{b\gamma - 1} ((b\gamma)^{h+1} - 1) \) are those at level \( h \) or above, which is a \( b^{-h} \) fraction of all nodes. Now, ignoring
constant factors and additive terms, we can solve for $h$ in $d = (b \gamma)^{h+1}$, giving us $h = \log_{b \gamma} d - 1$. As a result, roughly $d \approx \frac{\log b}{\log \gamma}$ nodes will have degree $d$, i.e., we again observe a power law.

### 6.4 Further Reading

In considering merely a log-log plot, power laws are easily confused with lognormal distributions. A random variable $X$ is lognormal if $\ln X$ has as normal (Gaussian) distribution. See the discussion in Mitzenmacher’s survey [257] for more details. Lognormal distributions can for instance be obtained by processes where in each step $t$, each quantity $X_j$ grows by a constant factor $F_{j+t}$ (rather than a constant additive term). If the $F_{j+t}$ are i.i.d., then in the limit, the $X_j$ will be lognormal. Huberman and Adamic [187] show that combining such a multiplicative growth process with different arrival times (as we saw in Section 6.1 for the special case $\alpha = 0$) also gives rise to power-law distributions. Similarly, Champernowne [76] shows that combining a multiplicative update with a hard lower bound on all sizes $X_j$ leads to a power law.

The power-law nature of the Internet topology reported in [137] has been the subject of quite a lot of discussion. As noted by Chen et al. [83], the data set used for analysis missed a large fraction of physical connections, comprising essentially a union of BFS trees of the topology. Lakhina et al. [231] showed via experiments that for sampling based on shortest paths like the one in the data set analyzed by Faloutsos, Faloutsos, and Faloutsos, even $G(n, p)$ random graphs would exhibit a power law (though with exponent $-1$). This was proved by Clauset and Moore [89], and in a more general setting by Achlioptas et al. [2]. The upshot here is that even random regular graphs will exhibit a power law if the edges are sampled in a manner akin to the traceroute tool used for the data set of [137].

Instead of asking for “natural” models which predict the structure of known graphs with power law degree distributions accurately, we could instead start out from the assumption that the graph is uniformly random subject to its degree distribution, and investigate which properties follow merely from this degree distribution. Along these lines, Molloy and Reed [260, 261] analyze the point at which a giant component emerges in a random graph with given distribution, and how the component’s size depends on the distribution. Using mean-field approximations and generating functions [340], Newman et al. [279] give an alternate derivation of the results of Molloy and Reed, and prove several other properties of random graphs with a given degree distribution.

While hard-wiring a desired degree sequence is a first step toward matching various properties of real-world observed networks, it is really only that. There are many other properties not nearly as well explained by such uniformly random graph models. One of the most prominent ones is “clustering”: the fact that nodes sharing a neighbor are more likely to be connected themselves. Several recent papers have proposed simple random graph models resulting in clustering. The simplest one is that of Random Intersection Graphs [315, 202]. Here, we are given two types of nodes: individuals and groups. We generate a random bipartite graph between those two types of nodes. Then, we connect two individuals if they share at least one group. Thus, all individuals in the same group form a clique, and we have built in significant clustering. This model has been extended to account for different degree distributions as well (see, e.g., [273, 112]).

While this class of models can match degree distributions while building in clustering, there may be other parameters of a real-world network we would also like to match. One class of random graphs lets us match basically any features of a graph we care about. The model is called ERG (Exponential Random Graphs) or $p^*$ graphs [149, 302, 334]. The high-level idea is to define several features (e.g., number of triangles, degrees, diameter, etc.), and posit that a graph $G$ is generated at random with probability proportional to \( \exp(\sum_{S \in G} \alpha_S) \), where $S$ ranges over all features of interest, and $[S \in G]$ denotes that feature $S$ is present in the graph $G$. (Notice that this approach is quite similar to the Markov Random Field analysis in Section 4.1.) If we have chosen the $\alpha_S$ parameters, we can sample random graphs from this distribution. A complementary problem is to find the parameter settings for $\alpha_S$ giving the maximum likelihood estimate of an observed graph $G$. Both problems appear to be computationally difficult in general, though there are heuristics for sampling which sometimes work well in practice.

One drawback of the basic preferential attachment model of Barabási et al. is that it produces only acyclic graphs. Hence, various papers have considered generalizations of this model. A fairly broad generalization
is due to Aiello et al. [9]. In their class of models, a random number of nodes and edges can be added in each step, including edges between previous nodes. The attachment is again preferential, and nodes can arrive with initial indegree and outdegree weights. The advantages of this model claimed in [9] include the fact that it can model arbitrary power-law degree distributions on both indegrees and outdegrees, while also allowing for arbitrary density, and subsuming most past models. A fairly detailed overview of models up to that point in time, together with results on structural properties and implications of these models (again using mean-field approximations) is given by Albert and Barabási [12].

Another generalization was proposed by Leskovec et al. [238] and termed the Forest Fire Model. It was based on the observation that real-world networks often appear to become denser over time (i.e., the number of edges increases super-linearly in the number of nodes), and that their diameter often shrinks over time. Standard random graph models do not predict such behavior. In the Forest Fire Model, nodes arrive one at a time. When a node arrives, it chooses one ambassador $w$ to link to. In addition, it selects $x \sim \text{Bin}(1/(1-p))$ of $w$’s neighbors, links to them, and recurses on these neighbors with the same value $x$. $p$ is a parameter of the model, and the model allows weighing forward and backward edges differently in determining a node’s neighbors. Using simulations, Leskovec et al. show that this model matches the evolution of density and diameter in several observed real-world graphs.

While the mean-field approximations in Section 6.1 and in several of the papers cited here are not mathematically rigorous, a beautiful result by Wormald [341, 342] gives sufficient conditions under which the results of the differential equations give a precise approximation (up to lower-order terms) of the correct outcomes, with high probability. A perhaps more intuitive description of this result is given by Achlioptas [1].
Chapter 7

Small-World Graphs and Decentralized Search

In the chapters so far, we have explored the structure of information networks, and how additional information can be extracted by looking at it in graph-theoretic terms. In the process, we found that it is often useful to think about generative models for the network structure: if we can explain how the network formed, then we might be in a better position to label or cluster its parts, or observe other useful high-level structures.

In this and the remaining chapters, we will think of the network not so much as being the artifact to analyze and extract information from, but as a substrate that itself disseminates information. Thus, we will analyze the dynamics of information spread within social and computer networks, and try to understand what network structure would enable quick and efficient dissemination.

We begin by discussing the small-world phenomenon, both in its historical context and algorithmic implications. The expression refers to the anecdotal observation that any two people appear to be connected by a small chain of intermediate social acquaintances. Restated in graph-theory terms, this means that social networks (graphs whose nodes are individuals, and whose edges capture friendships, acquaintances, or other pertinent interactions) appear to have small diameter.

Stanley Milgram, in 1967, was the first to design an experiment to explore this hypothesis more scientifically \[255\]. His experiment, described in more detail in Section 7.1, led him to the conclusion that the average path length between two individuals in the social network of acquaintances within the United States is about six hops. The phenomenon has made its way into pop culture, including movies and plays, via the phrase “Six Degrees of Separation”. Finding short paths in social networks is also the object of the game “Six Degrees of Kevin Bacon”: the goal is to find a short path from any actor to the actor Kevin Bacon in the graph consisting of actors and their co-appearances in movies\[2\]. Mathematicians have a similar concept: the Erdős number. It captures the shortest path distance between a scientist and the famous Hungarian mathematician Paul Erdős, where edges represent co-authorships of papers\[3\].

7.1 Milgram’s Experiment (1967)

In 1967, Stanley Milgram attempted to verify the small-world phenomenon quantitatively. His experimental setup was as follows: He selected random people from locations like Kansas or Nebraska, and asked them to forward a folder to a target in Cambridge, MA or Boston. The rules of forwarding were that the current holder of a letter could only mail it to someone they knew on a first-name basis. That person was to mail the

---

\[1\] Milgram’s other famous experiment explored obedience to authority, and involved making participants believe that they were administering very dangerous electro-shocks to other participants [254].

\[2\] See http://oracleofbacon.org.

\[3\] See http://www4.oakland.edu/enp/.
His first target person was the wife of a divinity student living in Cambridge, and starters were located in Wichita, Kansas. Milgram found that the very first folder reached her in just four days and took only two intermediate acquaintances. However, it appears that the completion rate of chains in this experiment was very low, so Milgram conducted a second study. Among other things, he made the folders much more impressive and official looking, which appears to have led to significantly higher participation. In the second study, the starters were located in Nebraska, and the target was a stockbroker from Sharon, MA, working in Boston. Milgram reported that “chains varied from two to ten intermediate acquaintances, with the median at five” [255, p. 65]. Subjects appeared to be able to reach the target with an average of six hops.

Graph-theoretically, Milgram was trying to find short paths in the social network of people in the US. Instead of using the “correct” approach of BFS, he had to rely on DFS, as the typical branching factor of several hundred social acquaintances would have resulted both in much work for all participating individuals, and an illegal chain letter using the US Postal System. So what Milgram found was really just an upper bound on the shortest paths.

As analyzed by Travers and Milgram [324], out of 296 chains, only 217 chains started, and 64 completed. The number of intermediate nodes varied between two and ten, with a median of 5 and a mean of 6. The frequencies that Milgram measured were monotonically increasing to chains of length 4, then dropped slightly for chains of length 5, and increased to the highest frequency for length 6, dropping again afterwards until length 10. Travers and Milgram argued that the drop in frequency at 5 is actually caused by the experimental results being a superposition of different populations.

The attrition rate (the frequency of not forwarding the letter) was about 25% at each step. Thus, it seems likely that the observed distribution is skewed to shorter chains: if a chain was naturally longer, then it has a higher probability of not finishing at all, and thus not being counted towards the average. White [339] gave a simple calculation correcting for this attrition. The argument is that if we observe a fraction of \( \frac{n}{j} \) chains of length \( j \), then the true fraction of such chains would have been roughly proportional to \( \frac{4}{3)^j \frac{n}{j}} \), as each chain of length \( j \) does not complete with probability \( (3/4)^j \). Correcting for attrition in this way, White suggests that Milgram’s data are more supportive of a median distance of 7–8 between individuals.

Additional interesting details about the Milgram experiments were that many chains had the same person as a last step: out of all completed chains, 25% went through the same last person, and 50% through one of three persons. This may be taken as an indication that society relies on “hub persons” for connectivity. Killworth and Bernard [215] analyzed the “rules” individuals used for deciding whom to forward letters to. In a separate experiment, they simply asked respondents (in West Virginia) to identify their “first hop” to a number of different targets. They found that geographic location and occupation were by far the dominant dimensions, with others used significantly less frequently.

Kleinfeld [223] revisits Milgram’s experiment and finds that there were a number of loopholes in Milgram’s experiment. Milgram’s sampling technique to determine starting points for the chains was biased since his target was a stockbroker, and a number of his starters were stockholders, who would be more likely to know stockbrokers. He recruited people through advertisements for people who were well connected, and the experiment was generally more likely to attract economically better-off people, who tend to have more long-distance connections. In addition, out of the Kansas and Nebraska studies, only the Nebraska one, which was much more successful, was published.

### 7.2 Formalization, and Expander Graphs

If we want to reason about small-world networks analytically, we first need to come up with a good definition of the concept. A simple (and popular) definition is to consider a graph a small world if it has small (e.g., logarithmic) diameter. However, this ignores another characteristic of small worlds, namely that the graph exhibits a lot of clustering and regular structure, similar to lattices. Watts and Strogatz [337] propose using the notion of clustering coefficient \( C \) (see, e.g., [333]), defined as the number of triangles in the graph divided by the number of adjacent edge pairs. Thus, it can be interpreted as the probability that two nodes are
connected, given that they share a mutual friend. For a complete graph, the clustering coefficient is 1, and for an empty graph, it is 0.

If we are interested merely in finding graphs with small diameter, our task is easy: the complete graph has diameter 1. However, the complete graph is certainly not a good model of social networks, as individuals tend to have a smaller number of contacts. So we are looking for graphs of small diameter and low degree.

One way in which we can achieve this goal is by virtue of expanders. A graph $G$ is an $\alpha$-expander (or $\alpha$ edge expander) if each vertex set $S \subseteq V$ with $|S| \leq n/2$ has at least $\alpha|S|$ edges leaving. Intuitively, this means that no two “large chunks” of the graph can be disconnected by removing few vertices or edges. Obviously, the larger $\alpha$, the better the expander the graph is. Normally, people are particularly interested in families of expanders where $\alpha$ is constant, i.e., does not become smaller as the number of nodes $n$ increases.

While the definition of expanders seems useful, we still have to find out whether they actually exist. It is pretty obvious that complete graphs are expanders (with $\alpha = n/2$), but again, we are looking for lower degrees. A binary tree is not a good expander, as the left half of the tree has $n/2$ nodes, but only one edge leaving. However, hypercubes are good expanders (with $\alpha = 1$), with degree $\log(n)$.

When we want to reduce the degree all the way down to constant, i.e., not increasing with $n$, the task becomes quite a bit more difficult. For a long time, no explicit constructions were known, and the first explicit constructions were based on Cayley graphs of rather difficult looking algebraic groups. More recently, Reingold et al. [299] presented a new way of constructing constant degree expanders via the zig-zag graph product. On the other hand, it is much easier to construct expanders randomly. For any $d \geq 3$, almost all $d$-regular graphs are actually expanders (notice that for $d = 2$, any $d$-regular graph is just a union of disjoint cycles, and thus certainly not an expander.)

It is not quite obvious how to generate a uniformly random $d$-regular graph. It is much easier to instead generate a $d$-regular multigraph, in which we also allow self-loops and parallel edges. For then, we can use the configuration model: each node has $d$ edge ends sticking out, and we construct a uniformly random matching among those edge ends, by randomly picking two of the remaining edge ends and pairing them up until all edge ends are matched. For this model, whenever $d \geq 6$, a fairly straightforward analysis using standard tail bounds (Martingale or Chernoff bounds) shows that with high probability, the resulting graph is an expander. For $3 \leq d < 6$, the analysis gets a little messier, as tail bounds are not quite strong enough: one has to reason a bit more carefully about binomial coefficients and such.

### 7.2.1 Properties of expanders

Expander graphs, and the notion of expansion, are important for several reasons. From a network design perspective, expander graphs are resilient to node or edge failures, as a small number of failures will not disconnect large parts of the network.

Expansion also plays an important role in the analysis of random walks and their mixing time (and thus in the context of MCMC sampling). The reason is that a low expansion means that few edges connect two large components. Thus, a random walk starting in one of the components has small probability of crossing the other component, and it will take a long time until the initial bias in the probability of being at any one node gets evened out.

For our purposes, the reason that expanders are interesting is that they have small diameter.

**Lemma 7.1** If $G$ is an expander of maximum degree $d$, then $G$ has diameter $O(\log(n))$.

**Proof.** Consider a BFS of $G$ starting from any node. For each layer $S$ of the BFS, at least $\alpha|S|$ edges are leaving the set of all nodes already reached, hence at least $\frac{\alpha}{2}|S|$ new nodes are hit, so long as $|S| \leq n/2$. Thus, at each layer, the total number of nodes in the BFS tree grows by a factor of at least $1 + \alpha/d$, and $\frac{\alpha}{2}$ nodes are reached in

$$\log_{1 + \alpha/d}(n/2) = O\left(\frac{\log n}{\log (1 + \alpha/d)}\right) = O\left(\frac{\log n}{\alpha/d}\right) = O\left(\frac{\log n}{\alpha^2}\right) = O\left(\frac{n}{\alpha^2}\right)$$

steps. Carrying out BFS from both the source $s$ and target $t$, they each reach $\frac{\alpha}{2}$ nodes. After one more step of BFS, they must therefore intersect at some node $v$, and the concatenation of the $s$-$v$ and $t$-$v$ paths gives an $s$-$t$ path of length $O\left(\frac{n}{\alpha^2}\right)$.
7.3 The Watts-Strogatz Model

In the past section, we began our quest for graphs that would model the empirically observed “small-world phenomenon”: the fact that most (or all) pairs of people tend to be connected through short paths of mutual acquaintances. We observed that expander graphs tend to have small diameter, i.e., short paths, and that random $d$-regular graphs for $d \geq 3$ are with high probability expander graphs.

While this in principle will give us “small worlds”, it is unsatisfactory in that real social networks do not at all look like random graphs. They exhibit a lot of “clustering” in the form of triangles or other short cycles. Indeed, two people sharing a common acquaintance are more likely to know each other as well. This observation has led to the definition of the clustering coefficient of node $v$ as

$$\frac{|\{(u, u') \mid (u, v) \in E, (u', v) \in E\}|}{\binom{d}{2}},$$

i.e., the fraction of actual triangles among pairs of neighbors of $v$. Random graphs have low clustering coefficient (roughly $d/n$ for $d$-regular graphs), while real social or collaboration networks have much higher clustering.

In order to form a more realistic model giving both small diameter and high clustering coefficient, Watts and Strogatz [337] propose a simple model of small worlds. Here, we look at a slight modification (mostly for analytical simplicity), but will still refer to it as the Watts-Strogatz model. Start with a (structured) graph $H$ (for instance, the 2-dimensional grid, or a ring), and add one or a few random edges from each node. (Alternately, “rewire” one or a few edges randomly.) The result will be that the edges of $H$ will lead to high clustering coefficient, while the random edges lead to small diameter, as per our analysis of expander graphs. The resulting graph will still “resemble” $H$ a lot. Intuitively, $H$ is supposed to model the “traditional” social network, based perhaps on physical proximity or related employment, while the random edges model “random” friends.

An interesting additional outcome of Milgram’s experiment [255] was not only the existence of short paths in social networks, but also the fact that individuals, who do not have a map of the graph, were able to actually find short paths. Their decisions were only based on local and qualitative information. Hence, an interesting question is what properties of a graph allow efficient decentralized routing based solely on local information. In addition to a better understanding of social networks, this question is also relevant for designing networks permitting simple routing.

To make this question more concrete, we can work with the Watts-Strogatz Model, and ask about the effect that different distributions of long-range links have on the ability to route with a decentralized algorithm. Watts and Strogatz considered the case of uniformly random links, i.e., each node links to one other uniformly randomly chosen one. We will show that in the case of the 2-dimensional grid of size $n \times n$, with this distribution, there is no local routing protocol reaching its destination in $o(n^{2/3})$ steps in expectation. Here, a local protocol is one in which a node only knows the history of the message in the network so far (including which nodes it has visited), and its own links.

For the proof, consider a ball $B$ of radius $n^{2/3}$ around the destination node $t$, and assume that the source $s$ of the message lies outside this ball. By the Principle of Deferred Decisions\footnote{The Principle of Deferred Decisions [266] states the intuitively obvious fact that we can defer coin flips or other random decisions until their outcomes are actually needed/referenced.}, we can assume that each node only generates its random outgoing link once it receives the message. Consider the first $\delta n^{2/3}$ steps of the protocol (for some constant $\delta \ll 1$). Because each node’s outgoing link hits $B$ with probability at most $O\left(\frac{n^{2/3}}{n^2}\right) = O(n^{-2/3})$, only $O(\delta)$ long-range links hit $B$ in expectation, and with constant probability (by Markov’s Inequality), none do. Thus, every step ending in $B$ must have been a step on the grid. As a result, with constant probability, $t$ was not reached in $\delta n^{2/3}$ steps, completing the proof. (Notice that the analysis is tight, in that in $O(n^{2/3})$ rounds, with constant probability, $B$ is reached, and within at most another $O(n^{2/3})$ grid steps, the message makes it to $t$).

To study the effect of the distribution of long-range links on the ability to route in a decentralized manner, let $\text{Prob}[v \to w]$ denote probability that $v$ connects to $w$ via a long-range link. Assume that this probability
is a function of $d_{v,w}$, the distance between $v$ and $w$. Intuitively, it seems that connections between distant individuals are less likely, so the probability should be monotone decreasing.

If $\text{Prob}[v \rightarrow w]$ is an inverse exponential function of $d_{v,w}$, then it decreases very rapidly, and long-range links are too unlikely. Hence, we use a polynomially decreasing function in $d_{v,w}$ [217], i.e., $\text{Prob}[v \rightarrow w] \sim (d_{v,w})^{-\alpha}$ for a constant $\alpha \geq 0$.

To understand this distribution well, we need to calculate its normalizing constant, i.e., the $\gamma$ such that $\text{Prob}[v \rightarrow w] = \frac{1}{\gamma} d^{-\alpha}$. By noticing that for each $d$, there are $\Theta(d)$ nodes at distance $d$ from any given node $v$, we can calculate

$$
\sum_{w} d_{v,w}^{-\alpha} \approx \sum_{d=1}^{\infty} d d^{-\alpha} = \sum_{d=1}^{\infty} d^{1-\alpha} = \begin{cases} 
\Theta(n^{2-\alpha}) & \text{for } \alpha < 2 \\
\Theta(\log n) & \text{for } \alpha = 2 \\
\Theta(1) & \text{for } \alpha > 2
\end{cases}
$$

We previously proved that for $\alpha = 0$, the case where the destination of each long-range link is a uniformly random node (independent of the distance between $v$ and $w$), nodes cannot route efficiently based solely on local information. For very large $\alpha$, we also expect local routing (or any routing, for that matter) to not work well, as long distance links will be exceedingly rare. Of the three cases (1) $\alpha = 2$ (2) $\alpha < 2$ and (3) $\alpha > 2$, it seems thus most promising that $\alpha = 2$ may work well for local routing. Indeed, we will show that it is the only exponent for which local routing can work in poly-logarithmic time.

**The case $\alpha = 2$**

**Claim 7.2** For $\alpha = 2$, we can route locally in $\Theta(\log^2(n))$ steps in expectation.

**Proof.** The algorithm is the simple greedy rule of always forwarding the message to the neighbor closest to the destination.

Let $s$ be the source node and $t$ the destination node, at distance $d$ from $s$. We will show that within an expected $O(\log n)$ steps, the distance of the message to $t$ is halved. (Notice that the time here is independent of $d$.) In order to prove this, we let $B_{d/2}(t) = \{v \mid d_{v,t} \leq d/2\}$ denote all nodes in the ball of radius $d/2$ around $t$.

Because there are $\Theta(d^2)$ nodes in $B_{d/2}(t)$, and each is at distance at most $\Theta(d)$ from $s$, we can lower-bound the probability for $s$’s long-range link to end in $B_{d/2}(t)$ as follows:

$$
\frac{1}{\Theta(\log n)} \sum_{v \in B_{d/2}(t)} d_{s,v}^{-\alpha} \geq \Theta \left( \frac{1}{\log n} \cdot d^2 \cdot d^{-\alpha} \right) = \Theta \left( \frac{1}{\log n} \right)
$$

Notice that this gives a notion of “uniformity of scales”: the probability of halving the distance is the same independently of what the distance itself is. Similarly, if we think of circles around $v$ of size $2^k$ for $k = 0, 1, \ldots$, then $v$ has (approximately) equal probability of having its long-range link in any of the annuli between circles of radius $2^k$ and $2^{k+1}$.

By our calculations above, any single long-distance edge reaches $B_{d/2}(t)$ with probability at least $\Theta \left( \frac{1}{\log n} \right)$. If not, the message is moved to another node, no further away, which again has at least the same probability of reaching $B_{d/2}(t)$, etc. Hence, the number of steps until a long-range edge will reach $B_{d/2}(t)$ is lower-bounded by a negative binomial random variable with parameter $\Theta \left( \frac{1}{\log n} \right)$. In particular, the expected number of steps to hit $B_{d/2}(t)$ is $\Theta(\log n)$, and the actual number is sharply concentrated around the expectation. As the distance of the message from $t$ is halved every $\Theta(\log n)$ steps (independently of the current distance), $t$ will be reached after $\Theta(\log n \log d) = O \left( \log^2 n \right)$ steps.

**The case $0 \leq \alpha < 2$**

We already saw that for $\alpha = 0$, local routing takes a polynomial number of steps. The problem was that links were too unstructured, and the chain was unlikely to encounter any node with long-range link into a (sufficiently small) polynomial-sized ball around the destination. Here, we will generalize the construction to show that the same problem occurs for all $\alpha < 2$. 


Claim 7.4 For $\alpha < 2$, local routing requires a polynomial number of steps in expectation.

Proof. Let $\delta = \frac{2-\alpha}{3}$, and let $B(t)$ be the ball of radius $n^\delta$ around $t$.

This time, we want to upper-bound the probability that a given node has a long-range link into $B(t)$. Even when two nodes are very close together, say $d = 1$, the probability of a long-range link between them is at most $\frac{1}{\gamma} = n^{\alpha-2} = n^{-3\delta}$. As there are only $\Theta(n^{2\delta})$ nodes in $B(t)$, the probability for a given node to have a long-range link into $B(t)$ is at most $\Theta(n^{-\delta})$. Thus, it takes $n^\delta$ steps in expectation to encounter a long-range link into $B(t)$. On the other hand, if no such long-range link is encountered, then the chain takes at least $n^\delta$ small steps (namely, all the last $n^\delta$ steps). In either case, the total number of steps is $\Omega(n^\delta)$.

The case $\alpha > 2$

For small $\alpha$, the problem is that, while short paths exist, they cannot be found with local information alone, as the random links are “too random”. For large $\alpha$, the problem is that there are not many long-range links — the argument below needs to be strengthened only slightly to show that the diameter of the graph is actually polynomial in $n$.

Claim 7.4 For $\alpha > 2$, the expected number of steps required to route from $s$ to $t$ is polynomial in $n$.

Proof. Following our intuition that long-distance links are unlikely, we first calculate the probability that a link is longer than some given number $m$. Using again that there are $\Theta(d)$ nodes at distance $d$ from a given node, this is at most

$$
\Theta(\sum_{d=m}^{\infty} dd^{-\alpha}) = \Theta(\sum_{d=m}^{\infty} d^{1-\alpha}) = \Theta(m^{2-\alpha}).
$$

Hence, the probability that a link’s length is at least, say, $n^{1/3}$, is at most $\Theta(n^{\frac{2}{3}-\alpha})$. So if we only take $n^{\frac{2}{3}-\alpha}$ steps, the probability of having any of them encounter a long-range link longer than $n^{1/3}$ is polynomially small ($O(n^{\frac{2}{3}-\alpha(1-1/\alpha)})$). But if none of them encounter any longer links, then the total distance covered is at most $n^{1/3} \cdot n^{\frac{2}{3}-\alpha} = n^{\frac{1}{3} - \alpha} = o(n)$, so the destination cannot be reached. (Notice that instead of $n^{1/3}$ and $n^{\frac{2}{3}-\alpha}$, we could have chosen other values $\beta > 0$ and $\gamma > 0$ such that the first $n^{\beta}$ steps don’t see a link of length greater than $n^{\gamma}$, so long as $\beta + \gamma < 1$.)

Based on simulations, it seems that the behavior is actually worse for $\alpha > 2$ than for $\alpha < 2$.

In the analysis, it is important to notice that while, for $n \to \infty$, any poly-logarithmic function is exponentially smaller than any polynomial one, this may not be so for finite, even fairly large, $n$. Indeed, it has been verified experimentally that the “correct” exponent $\alpha$ for finite $n$ is $\alpha = 2 - f(n)$, where $f(n) \to 0$ as $n \to \infty$. Determining the exact nature of $f(n)$ is still open. This may also help in explaining the apparent gap between the short paths observed by Milgram, and the apparently much longer paths guaranteed by our proof. In addition, our analysis only used one long-range link per node. If nodes have many more long-range links, the predicted path lengths will be smaller.

While our analysis was done for a 2-dimensional grid, we did not use many properties of it. In fact, the only property we used was that there were $\Theta(d)$ nodes at distance $d$ from a given one. The analysis extends easily to $r$-dimensional grids; the unique exponent for which local routing can be accomplished is then $\alpha = r$.

It also extends to hierarchical tree structures. For professions or other interests, a grid is perhaps not the right type of model. Instead, we may consider all professions to form a hierarchy (“scientific” vs. “arts” vs. “finance” etc., further subdivided into different sciences, subcategories, etc.). Each person is located at a leaf of this hierarchy, and generates $\Theta(\log^2 n)$ links to other people. The probability of linking to a given node decreases in the distance between the two nodes in the tree. It can be shown that for an appropriate choice of distribution, we still obtain local routing, and the distribution is closely related to the one studied above, in that it satisfies “uniformity over scales”.

Kleinberg [218] proves a more general result for a certain class of set systems which includes metric spaces that “grow uniformly”. Whether a similar result holds for arbitrary metric spaces is currently open. Watts
et al. [336] suggest a model wherein multiple metric spaces (or, in their case, hierarchies) combine to create the link structure. The idea is that links could be caused by proximity in a physical space (like the grid, above), or in the space of research topics, or others. Thus, the distance between two individuals is their minimum distance over all of the hierarchies, and individuals know about the hierarchy structure in their routing decisions. (See also the discussion by Killworth and Bernard [215].) Unfortunately, the result is not a metric, and it is not clear whether results can be proved formally about this model; Watts et al. based their investigation on simulations.

7.4 Further Reading

Since the original work by Kleinberg [217] on the algorithmic models, many subsequent papers have studied generalizations of the models in different respects. Excellent surveys of both the results described above and much of the follow-up work are given by Kleinberg [219] and Fraigniaud [148].

One of the interesting directions is how much extra knowledge the members of the network need in order to find (nearly) shortest paths. While we showed above that the greedy algorithms finds paths of length $O(\log^2 n)$, these are not necessarily shortest for every pair of nodes. Lebhar and Schabanel [234] analyze the effects of performing some depth-first exploration of the neighborhood of the current message holder before actually forwarding the message. They show that this results in much shorter paths of length $O(\log n (\log \log n)^2)$ (for our model of one random link). Manku, Naor, and Wieder [246] show that when the average degree of nodes is $O(\log n)$, by generating each long-range link at distance $d$ independently with probability $1/d$, a mere 1-step lookahead leads to paths of length $O(\log n / \log \log n)$.

Another popular direction for extensions is to change the underlying graph structure from a grid, or the precise long-range link distribution. One of the most comprehensive studies of this topic is by Nguyen and Martel [282], who analyze the diameter of the above model in $D$ dimensions, and show that short paths exist for any exponent $\alpha < 2D$, while the diameter is polynomial for $\alpha > 2D$. (The case $\alpha = 2D$ is currently open.) They also relate these diameter bounds and greedy routing more generally to expansion properties.

The Watts-Strogatz model for small-world networks can be extended to settings beyond social networks. For example, Menczer [251] posits a model for generation of web content and links wherein the metric is obtained by textual similarity of web pages (using standard bag-of-words models), and links follow a power-law distribution in this metric. He compares the model to real-world data, and infers that it implies searchability of the Web by greedy algorithms.
Chapter 8

Diffusion of Information and Influence in Social Networks

In the first few chapters, we were studying networks mostly as “enhancing” the information stored in them. In the previous chapter, for the first time, we considered networks as active entities forwarding information between their nodes. We continue the study of such dynamics in this and the next chapter, focusing on the interactions between the structure of a network and its role in disseminating information, diseases, data, etc. As some examples of the types of “things” that spread through a network, consider the following:

**Biological Networks**: Diseases and mutations. Dominant mutations result in more and more affected entities from generation to generation, until almost all of the species is affected. Diseases spread through contact, while mutations spread through natural selection in reaction to evolutionary constraints.

**Technological Networks**: Viruses and worms, but also power failures or car traffic. Power failures in a power grid increase the load on other generating stations causing them to also break down, etc [21, 22, 335]. Congestion on freeways leads to similar behavior, spreading to nearby freeways or roads by virtue of drivers’ evasive strategies. Viruses or worms directly use human or technological connections to spread from one host to another.

**Social Networks**: Rumors or group behavior such as riots. A rumor spreads from person to person until a large subset of people are aware of it. Similarly, peer pressure or “safety in numbers” entices people to participate in behavior shared by others, such as rioting [308, 170].

**Strategies**: as a special case of social networks, some trends, such as cell phone usage or car sizes, spreads as a result of “strategy optimization” [49, 125, 263, 344]. For example, to not be at a disadvantage on the road with bigger cars like SUVs around, people themselves need bigger cars. Similarly, once more others have adopted a new innovation, such as a new media format, it may become beneficial or necessary to adopt it also [188].

The common thread of all of these examples is that a new behavior or piece of information starts out with one or a few individuals, and gradually propagates through the network to (sometimes) reach a much larger number eventually.

In order to analyze these types of dynamics, we need to first define a model for the effect of the network on individuals’ behavior. Here, we can draw on several decades worth of work in sociology, biology, and economics.

### 8.1 The Bass Model for Innovation Diffusion

While there had been many empirical studies of the diffusion of technological innovations in society (see, e.g., [96, 303], or [66] for a study of the relationship between innovation diffusion and strength of ties [169]),
perhaps the first descriptive mathematical model is due to Bass [33]. It is similar to standard models for epidemic diseases [24, 115]. The presentation here is inspired by that of Jackson [190].

The population is assumed to be infinite and mixed, i.e., each individual is equally likely to interact with each other individual. Two different types of “pressure” affect individuals: mass communication such as advertising, and word-of-mouth interactions with others who have adopted the innovation already.

Specifically, assume that an individual has not yet adopted the innovation. In each “time step”, he will be convinced by the “mass marketing” with probability $p$. In addition, he may meet another individual who has already adopted the innovation. In that case, he will adopt the innovation with probability $q$. (Let’s assume that $q = (1 - p)q'$, i.e., $q$ accounts only for the case when the individual has not adopted the product due to mass marketing.) If at time $t$, an $N(t)$ fraction of the population have already adopted the innovation, then the meeting probability is $N(t)$ for any individual. Thus, the individual’s probability of adopting is $p + q \cdot N(t)$. Since the fraction of individuals for which this calculation applies is $1 - N(t)$, we obtain that the increase in the fraction of adopting individuals at time $t$ is $(1 - N(t)) \cdot (p + q \cdot N(t))$. This leads to the differential equation

$$\frac{dN(t)}{dt} = (1 - N(t)) \cdot (p + q \cdot N(t)).$$

This differential equation can then be solved, and the effects of varying $p$ and $q$ (as well as other assumptions) can be studied. Following the initial work of Bass [33], there has been a long line of extensions and variation. See, for instance, [243, 326].

### 8.2 Schelling’s model for segregation

The Bass model assumes a homogeneous population, and furthermore assumes that each individual is equally likely to meet each other individual. Thus, it does not take any network structure into account. One of the first models to explicitly consider network structure in such behavior at a population scale was Schelling’s model for segregation [308]. Schelling was motivated by the question: why is it that most neighborhoods are very uniform (racially, and in other respects), even though most people profess that they would prefer to live in a diverse community?

Schelling proposed the following model: Assume that roughly $\frac{n^2}{2}$ individuals live on an $n \times n$ grid. Each node wants to make sure to not be completely isolated (the odd person in a community): formally, if less than an $\varepsilon$ fraction of $v$’s neighbors (in a small ball around $v$) are of the same color as $v$, then $v$ is unhappy and moves to a spot where it is not unhappy (say, the closest, or a random, such spot).

What Schelling observed was that even with a relatively small value of $\varepsilon \approx \frac{1}{3}$, neighborhoods end up mostly segregated: when the process quiesces, about $\frac{4}{5}$ of each node’s neighbors are of the same color as the node itself. While this result has been verified experimentally, it appears that no result formally proves this observation as of yet. It is also not known how the result relates with the topology of the neighborhood graph.

### 8.3 Granovetter’s threshold model

Schelling’s model addresses one shortcoming of the Bass model: the assumption of a uniformly mixed population. However, it still assumes that all individuals act the same. Granovetter [170], motivated by the study of outbreaks of riots or other collective behavior, notices that cities or neighborhoods with very similar statistical demographics often tend to exhibit very different overall behavior of a riot or similar event. Thus, he concludes that an accurate description of group behavior is impossible in terms of merely high-level statistical data, and needs to take individuals’ tendencies into account.

[170] therefore proposes a threshold model which has formed the basis of much subsequent work. The simplest version can be described as follows: Each individual (node) has a threshold $t_v$. If $t_v$ other nodes are already active (have adopted the behavior, such as starting to riot, or using a cell phone), then $v$ joins (becomes active).
As an example, consider the case of 100 people present in a city square in a tense environment, where the threshold of node \( v_i \) is \( t_{v_0} = 0, t_{v_1} = 1, \ldots, t_{v_{99}} = 99 \). In this case, \( v_0 \) becomes active first, followed by \( v_1 \), etc., until the entire population ends up active. On the other hand, if we make a small change and set \( t_{v_1} = 2 \), only node 0, and no other node, becomes active, even though 99 of the 100 nodes have the same thresholds as before. While this example may be a bit contrived, it clearly shows that the outcome of such a process can change dramatically even when the setup is almost identical.

One of the first questions we may wish to answer about this model is: given the nodes’ thresholds, can we predict how many nodes will be active in the end. We can define \( F(x) = |\{v \mid t_v < x\}| \). Then, the number of active nodes is the smallest \( x \) with \( F(x) \leq x \). Figure 8.1 shows the plot of \( F(x) \) vs. \( x \). The vertical lines denote the number of nodes active after 0, 1, ... time steps. The number of nodes that will be active in the next time step is then \( F(F(x)) \), which we can determine pictorially by moving horizontally to the diagonal, and then vertically to the next intersection with \( F(x) \). The process then quiesces at the first point for which the function \( F \) crosses the diagonal, as that is a fixed point.

![Figure 8.1: The threshold process and its dependency on \( F(x) \)](image)

The threshold values we used in our introductory example were clearly somewhat contrived. Instead, we may wonder what happens when the values are drawn from a more “reasonable” distribution. For instance, we may assume that the thresholds are distributed as Gaussians, with mean, say, \( \frac{n}{4} \) and variance \( \sigma^2 \). We could then investigate what is the result of varying \( \sigma^2 \), i.e., making the thresholds more or less sharply concentrated.

For small \( \sigma \), almost no nodes end up active, because almost no thresholds will be 0 or close to 0. As \( \sigma \) is increased, there is a phase transition from almost no one active to almost everyone active. A further increase of \( \sigma \) leads to a slight drop, steadying around \( \frac{n}{2} \).

Granovetter’s threshold model is clearly very simple. We could make it more interesting and realistic in several ways:

**Graphs**: We assumed that all nodes affect all other nodes. In a large social network, this is certainly not true. We may consider a graph in which nodes can only affect their neighbors.
Weights: The model of graphs can be further generalized. In our model so far, each node is influenced equally by all of its neighbors. Instead, each edge in the graph could have a weight associated with it, representing how much influence one endpoint has over the other. (Edges may be directed, and have different weights in both directions.)

Probabilistic activations: In real world scenarios, some nodes may fail in activating their neighbors. This can be modeled using probabilities on activations. The resulting model may be somewhat different from Granovetter’s. (A simple model is described by Goldenberg et al. [163, 162].)

Non-monotonic behavior: In many scenarios, a large rate of “participation” may discourage nodes from being active. For instance, a new fashion trend, initially spread by copying behavior, may become uninteresting to the trend-setters if many others are following it. As a result, in addition to the activation threshold \( t_{1,v} \), we would have a deactivation threshold \( t_{2,v} > t_{1,v} \): if more than \( t_{2,v} \) nodes are active, then \( v \) becomes inactive again. In its most general version, this model subsumes cellular automata and the game of life (and hence, can simulate computation).

Deriving thresholds: So far, we assumed that the thresholds are known. Actually obtaining or estimating these data is an interesting question, and may be partially amenable to data mining techniques. (In the absence of more information, we may assume that thresholds are random, or all identical, but this is only a very crude approximation.) One interesting way to derive thresholds in the strategic settings described earlier is to cast them in terms of the underlying game played by the nodes, e.g., weighing the costs and benefits of a larger car depending on the number of other people having large vs. small cars.

Starting or preventing cascades: Understanding a model is usually only the first step towards acting on it. In the above scenarios, we may be interested in containing the spread of an epidemic, computer virus, or power failure, or promote the adoption of a product.

8.4 The Contagion model

In investigating the problem of starting or preventing cascades, a first question we may want to answer is how many nodes it takes to infect the entire graph. This leads us to a question considered by Morris [263]. His model is as follows: the graph \( G \) is assumed to be infinite, but each vertex has finite degree. All vertices have the same threshold \( p \in [0,1] \): here, this means that they become active if at least a \( p \)-fraction of their neighbors is active. (Notice that the number of neighbors may be different for different nodes, but the fraction is the same).

These thresholds can be derived through a coordination game as follows [125]: Suppose that we have the following payoff matrix with \( q \geq \frac{1}{2} \):

<table>
<thead>
<tr>
<th></th>
<th>inactive</th>
<th>active</th>
</tr>
</thead>
<tbody>
<tr>
<td>inactive</td>
<td>( 1-q, 1-q )</td>
<td>( a, b )</td>
</tr>
<tr>
<td>active</td>
<td>( b, a )</td>
<td>( q, q )</td>
</tr>
</tbody>
</table>

meaning that if both players are inactive, they each get a payoff of \( 1-q \) (for instance, by communicating with some old technology). If both are active, they each get a payoff of \( q \geq 1-q \), having switched to the new technology. When player 1 is using the old technology, and player 2 the new technology, their respective payoffs are \( a \) and \( b \); usually, we would assume that \( a < 1-q \) and \( b < q \), since the players may be less able to communicate or share technology. (An important special case is \( a = b = 0 \), meaning that the players cannot communicate using different technology or languages.) Then, a node \( v \) with \( d \) neighbors, of whom \( r \) are active, will prefer to be active if and only if \( rq + (d-r)b > ra + (d-r)(1-q) \), or \( r > \frac{1-q}{1-(a+b)} \cdot d \). Thus, we have derived \( \frac{1-q}{1-(a+b)} \cdot d \) as an activation threshold.

A set \( X \subseteq V \) is called contagious, if starting from \( X \), all nodes are activated eventually. The contagion threshold \( t(G) \) is the supremum of all \( p \) such that there exists a finite contagious set for threshold \( p \).
In looking for obstacles to a set $X$ being contagious, we notice that each node that is not infected must have at least a $1 - p$ fraction of its edges to other uninfected nodes. Hence, if $X^*$ denotes the set of all nodes reached from $X$, then $V \setminus X^*$ is is active previously, and all nodes have threshold $p$. Define $f_p^k(S) := f_p(f_p^{k-1}(S))$ inductively to be the set of nodes active after $k$ steps (with $f_p^0(S) := S$, of course).

1. Assume that $c(G) < p$. By definition of the community threshold, this means that every cofinite set $S$ contains an infinite $(1 - p)$-community.

Let $X$ be an arbitrary finite set. Then, $V \setminus X$ is co-finite, and thus contains an infinite $(1 - p)$-community $C$. Because $C \cap X = \emptyset$, induction on the number of steps $k$ shows that no node from $C$ is ever active. Hence, $X$ cannot be contagious, and because $X$ was an arbitrary finite set, we have proved that $t(G) < p$.

2. Assume that $t(G) < p$. Let $S$ be a co-finite set, so $V \setminus S$ is finite. Therefore, by assumption, it is not contagious for parameter $p$. Let $Y := \bigcup_{k \geq 0} f_p^k(V \setminus S)$. Then, $Y$ is a $(1 - p)$-community, as each node of $Y$ must have strictly less than a $p$ fraction of its edges to nodes from $Y$. Moreover, $Y$ is infinite, for otherwise, $Y \cup S$ would be a finite contagious set, a contradiction. Hence, every co-finite set $S$ contains such an infinite $(1 - p)$-community $Y$.

Next, we would like to know how large contagion thresholds can possibly be supported by any graphs. We will show that $\frac{1}{2}$ is the limit:

**Theorem 8.2** There is no graph $G$ with $t(G) > \frac{1}{2}$.

**Proof.** For a set $S$, let $\phi(S) := |\delta(S)|$ be the number of edges leaving $S$. For $i \geq 0$, let $S_i$ denote the set of infected nodes after $i$ rounds. Observe that if at any point $S_i = S_{i+1}$, then clearly the infection process has terminated, as no new nodes can ever be activated.

Let $p > \frac{1}{2}$ be arbitrary, and $S_0$ a finite set starting out infected. Consider the sequence $\phi(S_i)$, $i \geq 0$, the number of edges leaving the infected sets in each iteration. When node $v$ moves from $S_i$ to $S_{i+1}$ (i.e., $v$ gets infected in the $i^{th}$ step), node $v$ has strictly more edges into $S_i$ than into $S_i$, because $p > \frac{1}{2}$. In the $(i+1)^{th}$ iteration, all the edges from $v$ to $S_i$ are not cut any more, whereas those from $v$ to $S_{i+1}$ are cut. Summing over all $v$, $\phi$ strictly decreases in each iteration in which nodes become infected. Because $S_0$ is a finite set, and each node has finite degree, $\phi(S_0)$ too must be finite. Therefore, $\phi$ can only decrease a finite number of times, and $S_k = S_{k+1}$ for some $k$. As only finitely many nodes become infected each round, only finitely many nodes will become infected eventually. This holds for any finite set $S_0$, and any $p > \frac{1}{2}$; therefore, the contagion threshold $t(G)$ cannot exceed $\frac{1}{2}$.

---

1 A cofinite set is a set whose complement is finite.
In view of the above fact, a natural question to ask is which graphs have \( t(G) \) close to \( \frac{1}{2} \). We begin with some definitions.

A labeling of a graph \( G \) is a bijection \( \lambda : \mathbb{N} \to V \). A labeling \( \lambda \) for \( G \) is \( p \)-inductive with parameter \( k_0 \in \mathbb{N} \) if for all \( k \geq k_0 \), each node \( \lambda(k) \) has at least a \( p \)-fraction of edges to \( \lambda(0), \ldots, \lambda(k-1) \). We define \( \ell(G) \) to be the supremum over all \( p \) such that \( G \) has a \( p \)-inductive labeling. This new measure is again equal to the contagion and community thresholds.

**Theorem 8.3** \( \ell(G) = t(G) \).

**Proof.** Start with nodes \( \lambda(0), \ldots, \lambda(k_0) \) infected. Then, by definition of \( p \)-inductiveness, all subsequent nodes will become infected at threshold \( p \). Similarly, the order in which nodes become infected can be used as a labeling, which is easily seen to be \( p \)-inductive.

A BFS-labeling is a labeling \( \lambda \) for graph \( G \) consistent with the BFS traversal of \( G \) from some finite node set \( X \). It is \( \delta \)-smooth for parameters \( p \) and \( k_0 \) if all nodes \( \lambda(k) \) for \( k \geq k_0 \) have between a \((p - \delta)\)-fraction and a \( p \)-fraction of edges to \( \lambda(0), \ldots, \lambda(k-1) \). Also, recall that a graph \( G \) has subexponential growth if for all \( c > 1 \) and all finite sets \( X \), the ball sizes are \( |B_r(X)| = o(c^r) \). (Where \( B_r(X) := \{ u \mid d(u, X) \leq r \} \) is the set of all nodes at distance at most \( r \) from \( X \).)

**Theorem 8.4** If a graph \( G \) has subexponential growth and a \( \delta \)-smooth BFS-labeling, then \( t(G) \geq \frac{1}{2} - \delta \).

Notice that this is not an “if and only if” statement. For instance, consider the grid. For purposes of infection, this graph acts very much like the grid (as most edges are grid edges). However, as we discussed in Section 7.3, this graph has exponential growth. Indeed, one might consider a definition of “small world” graphs to have both exponential growth and \( t(G) \) close to \( \frac{1}{2} \).

**Proof.** We will show that the labeling must have \( p = \frac{1}{2} \). Then, because the labeling is \( \frac{1}{2} - \delta \) inductive, Theorem 8.3 implies the result. We prove the statement by contraposition: if \( p < \frac{1}{2} \), then the graph must have exponential growth. Consider a BFS consistent with \( \lambda \). For each \( v \), let \( f(v) \) be the number of edges between \( v \) and nodes \( u \) with lower labels, and \( g(v) \) the number of edges between \( v \) and nodes \( u \) with higher labels.

By definition, \( f(v) \leq p(f(v) + g(v)) \), so \( g(v) \geq \frac{1-p}{p} \cdot f(v) \). Let \( L_r = \{ u \mid d(X, v) = r \} \) be the \( r \)th layer of the BFS from \( X \). We denote the edges within \( L_r \) by \( I_r \), and the edges from \( L_{r-1} \) to \( L_r \) by \( F_r \). Then, \( \sum_{v \in L_r} f(v) = |I_r| + |F_r| \) (counting each edge towards the value of the higher endpoint), and \( \sum_{v \in L_r} g(v) = |I_r| + |F_{r+1}| \) (counting edges for lower endpoints). Combining this with the inequality between \( f(v) \) and \( g(v) \), we obtain that \( |F_r| + |F_{r+1}| \geq \left( \frac{1-p}{p} \right) \cdot (|I_r| + |F_r|) \), which we can rearrange to obtain

\[
|F_{r+1}| \geq \frac{1-2p}{p} \cdot |I_r| + \frac{1-p}{p} \cdot |F_r| \geq \frac{1-p}{p} \cdot |F_r|,
\]

because \( \frac{1-2p}{p} \geq 0 \). Therefore, the sizes of \( F_r \) grow exponentially, and since the nodes have finite degrees, the number of vertices in the \( r \)th layer must also grow exponentially.

### 8.5 A non-monotone infection model

In the previous section, we began an investigation of models for the spread of infections on graphs. We assumed the following notion of monotonicity: once a node becomes infected, it stays infected. For certain types of diseases or behaviors, this model may be very accurate. For others, nodes will reevaluate their choice in every time step, based on the choices of their neighbors. A natural modification of the previous model posits that if in the previous time step, at least a \( p \) fraction of a node’s neighbors are infected, then the node will be infected in the next step, and otherwise, it will not.
In thinking about this new definition, a first question is whether any infinite graphs can be infected starting from a finite set of infected nodes. It is pretty easy to see that the answer is “Yes”: if we start with two adjacent infected nodes on the infinite line, and \( p < \frac{1}{2} \), then the entire line will eventually become infected. In fact, for infinite graphs, Morris [263] showed the following lemma:

**Lemma 8.5** If the (infinite) graph \( G \) has a (finite) contagious set with threshold \( p \) in the monotone model, then it also has a finite contagious set with the same threshold in the non-monotone case. The starting set in the non-monotone case will be the union of the starting set of the monotone case and its neighborhood.

Notice that our example of an infinite graph with finite infectious set crucially used that \( p < \frac{1}{2} \). It seems much more difficult to come up with examples having \( p \geq \frac{1}{2} \). In fact, is not even obvious how to construct arbitrarily large finite graphs that can be infected by a small set of nodes.

For the monotone case, the question with finite graphs is easy: take a star graph with threshold 1/2, and start with just the center node infected. However, in the non-monotone case, this graph will end up oscillating with period 2, between the center node infected, and all leaf nodes infected. In fact, this is far from accidental: the following is a corollary of slightly more general results of Goles and Olivos [164] and Poljak and Sura [291].

**Theorem 8.6 ([164],[291])** In the non-monotone case for a finite graph, the infection either converges, or oscillates with a period of 2.

Returning to the question of the existence of small infectious sets for arbitrarily large graphs, Berger [35] answered it in the affirmative for \( p = \frac{1}{2} \).

**Theorem 8.7 (Berger [35])** There are arbitrarily large graphs that can be infected (in the non-monotone case) with \( p = \frac{1}{2} \) and a starting set \( S \) with \( |S| \leq 18 \).

At this point, it is open whether such finite infectious sets exist for \( p > \frac{1}{2} \). However, we have the following partial result, showing that arbitrarily large \( p \) will not work.

**Lemma 8.8** It is impossible to infect arbitrarily large graphs if \( p > 3/4 \).

**Proof.** We denote by \( A_t \) the set of nodes active in round \( t \), and write \( I_t := A_t \cap A_{t-1} \) for the set of nodes active both in rounds \( t \) and \( t-1 \). Then, let \( S_t := \bigcup_{i \leq t} I_t \) be the set of nodes which were active in any two consecutive rounds up to and including round \( t \). Finally, let \( \delta(S_t) \) be the set of edges leaving the set \( S_t \), and \( \sigma_t := |S_t| + |\delta(S_t)| \).

We will show that if we start with \( c \) nodes active, then \( \sigma_t = O(c^2) \) for all \( t \). For this purpose, we show that \( \sigma_1 = O(c^2) \), and that \( \sigma \) is a non-increasing function of \( t \).

For the case \( t = 1 \), notice that \( S_1 = A_0 \cap A_1 \subseteq A_0 \). Thus, \( |S_1| \leq c \). Since each \( v \in S_1 \) is active at time \( t = 1 \), its degree can be at most \( \frac{c}{2} \) (for at least \( \frac{c}{4} \) of its neighbors must be in \( A_0 \)). Thus, the sum of all degrees in \( S_1 \) is at most \( c \cdot \frac{c}{2} \), which is \( O(c^2) \).

To show that \( \sigma_{t+1} \leq \sigma_t \), consider any node \( v \in S_{t+1} \setminus S_t \). Due to the addition of \( v \), the \( |S_{t+1}| \) term of \( \sigma_{t+1} \) increases by one. On the other hand, because \( v \in A_{t+1} \cap A_t \), more than \( \frac{c}{2} \) of \( v \)'s neighbors were active at time \( t \), and also at time \( t-1 \). Thus, more than half of \( v \)'s neighbors are active both at times \( t \) and \( t-1 \), i.e., they are in \( A_t \cap A_{t-1} \subseteq S_t \). For each of those neighbors \( u \in I_t \), there was an edge \((u, v) \in \delta(S_t)\), whereas \((u, v) \notin \delta(S_{t+1})\) any more. For each neighbor \( u \notin I_t \), we introduced at most one new edge into \( \delta(S_{t+1}) \) (or none, if \( u \) was also added to \( S_{t+1} \)). Since there are strictly fewer neighbors not in \( I_t \) than neighbors in \( I_t \), \( |\delta(S_t)| \) decreases by at least 1 for each node \( v \in S_{t+1} \setminus S_t \). Therefore, \( \sigma \) cannot increase overall. \( \blacksquare \)
8.6 Causing a large spread of an infection

So far, the results we studied were existential. We showed that there are large graphs that can be infected (in various models) by a small set of initial nodes. From a practical point of view, a much more relevant question is how to reach as many nodes as possible in a given graph. Naturally, this question has applications in using network effects for the purpose of marketing. More formally, we can look at the following question: Given a number \( k \), which set \( S \) with \( |S| \leq k \) will eventually infect as many nodes of \( G \) as possible?

As an optimization problem motivated by viral marketing, this question was first suggested by Domingos and Richardson [119, 301], who studied it first [119] for a very general model in terms of Markov Random Fields (see Section 4.1), and then for a much more restricted linear model [301]. Unluckily, even in the Morris Contagion model from Section 8.4, this problem not only is NP-hard, but also hard to approximate.

**Lemma 8.9** Unless \( P = NP \), the above problem cannot be approximated to within \( O(n^{1-\epsilon}) \) for any \( \epsilon > 0 \).

**Proof.** We prove this with a reduction from Set Cover: Given sets \( S_1, \ldots, S_m \), each a subset of \( \{1, \ldots, n\} \), and a number \( k \in \mathbb{N} \), are there \( k \) sets \( S_i \) whose union is the entire set?

For the reduction, we create a directed graph as follows: Let \( \{s_1, s_2, \ldots, s_m\} \) be nodes corresponding to the \( m \) subsets and \( \{u_1, u_2, \ldots, u_n\} \) be nodes corresponding to the \( n \) elements. Our construction will ensure that \( u_i \) becomes active when at least one of the nodes corresponding to sets containing \( u_i \) is active. We achieve this by connecting each \( s_k \) to the \( u_i \)’s in it, and setting a threshold of \( 1/m \) for each \( u_i \). Next, for a large constant \( c \), we add \( N = nc \) more nodes \( \{x_1, x_2, \ldots, x_N\} \). Each \( x_j \) is connected to all of the nodes \( u_i \), and it becomes active only when all of the \( u_i \) are (i.e., the \( x_j \) have a threshold of 1). The construction is depicted graphically in Figure 8.2.

![Figure 8.2: Construction for proving hardness of approximation](image)

If there are at most \( k \) sets that cover all elements, then activating the nodes corresponding to these \( k \) sets will activate all of the nodes \( u_i \), and thus also all of the \( x_j \). In total, at least \( N + n + k \) nodes will be active. Conversely, if there is no set cover of size \( k \), then no targeted set will activate all of the \( u_i \), and hence none of the \( x_j \) will become active (unless targeted). (Here, we are using implicitly that \( k < n \). If \( k \geq n \), then it is trivial to decide if all sets can be activated, for a Greedy Algorithm that adds sets to activate at least one more node in each iteration will work.) In particular, fewer than \( n + k \) nodes are active in the end. If an algorithm could approximate the problem within \( n^{1-\epsilon} \) for any \( \epsilon \), it could distinguish between the cases where \( N + n + k \) nodes are active in the end, and where fewer than \( n + k \) are. But this would solve the underlying instance of Set Cover, and therefore is impossible assuming \( P \neq NP \).
Notice that this proof can be modified to deal with uniform thresholds, by adding more nodes.

In the Morris model, the problem thus turns out to be completely intractable. Yet, we may be interested in finding models that are more amenable to approximation. One small modification that turns out to make a significant difference is to assume that the thresholds are uniformly (and independently) random instead of deterministic. Thus, we obtain a model where each edge has a weight \( w_e \geq 0 \) such that \( \sum_{e \in V} w_e \leq 1 \). Each node \( v \) independently chooses a threshold \( \theta_v \in [0, 1] \) uniformly at random. The goal is to choose a set \( S \) with \( |S| \leq k \) reaching as large a set as possible in expectation (over the random choices of \( \theta_v \)). We define \( f(S) \) to be objective function, i.e., the expected size of the finally infected set if we start with set \( S \) infected. The following theorem is the main result for this section:

**Theorem 8.10** There is a \( 1 - \frac{1}{e} \) approximation algorithm for the problem of selecting the set \( S \) maximizing \( f(S) \).

In fact, a simple greedy algorithm will give us the desired guarantees:

**Algorithm 5** The simple greedy algorithm
1: Start with \( S = \emptyset \).
2: for \( k \) iterations do
3: Add to \( S \) the node \( v \) maximizing \( f(S + v) - f(S) \).

The proof of the performance guarantee consists of three parts, captured by the following lemmas.

**Lemma 8.11** A node \( v \) approximately maximizing \( f(S + v) - f(S) \) can be found in polynomial time.

**Lemma 8.12** \( f \) is a non-negative, monotone and submodular function of \( S \).

Recall that a function on sets is monotone if \( f(S') \geq f(S) \) whenever \( S' \supseteq S \), and submodular (having diminishing returns) if \( f(S' \cup \{x\}) - f(S') \leq f(S \cup \{x\}) - f(S) \) whenever \( S' \supseteq S \). Equivalently, submodularity is characterized by the condition that \( f(S) + f(T) \geq f(S \cap T) + f(S \cup T) \) for all sets \( S, T \). Also notice that the non-negativity and monotonicity of \( f \) are obvious.

**Lemma 8.13** (Nemhauser/Wolsey/Fischer [102, 268]) If \( f \) is a non-negative, monotone and submodular function, then the greedy algorithm is a \( (1 - \frac{1}{e}) \)-approximation for the problem of maximizing \( f(S) \) subject to the constraint that \( |S| = k \).

We begin by proving the last lemma, which is naturally useful for other problems as well (such as Set Cover).

**Proof of Lemma 8.13.** Let \( v_1, v_2, \ldots, v_k \) be the nodes selected by the greedy algorithm (in the order in which they were selected), and denote \( S_i = \{ v_1, v_2, \ldots, v_i \} \). Then, the marginal benefit derived from the addition of element \( v_i \) is \( \delta_i = f(S_i) - f(S_{i-1}) \). Let \( T \) be the optimal solution, and \( W_i = T \cup S_i \).

First, the monotonicity of \( f \) implies that \( f(T) \leq f(W_i) \) for all \( i \). Because the algorithm chooses to add the best available node in the \( (i + 1)\)st iteration, and the benefit of any elements added later cannot be greater by submodularity, the total objective value for the set \( W_i \) is at most \( f(W_i) \leq f(S_i) + k\delta_{i+1} \), and thus also \( f(T) \leq f(S_i) + k\delta_{i+1} \).

Solving this for \( \delta_{i+1} \), and using that \( f(S_{i+1}) = f(S_i) + \delta_{i+1} \) now shows that \( f(S_{i+1}) \geq f(S_i) + \frac{1}{k} \cdot (f(T) - f(S_i)) \). We will prove by induction that \( f(S_i) \geq (1 - (1 - \frac{1}{e})^i) \cdot f(T) \). The base case \( i = 0 \) is trivial.
For the inductive step from \( i \) to \( i + 1 \), we use the above inequality to write

\[
\begin{align*}
f(S_{i+1}) &\geq f(S_i) + \frac{1}{k} \cdot (f(T) - f(S_i)) \\
&= \left(1 - \frac{1}{k}\right) f(S_i) + \frac{1}{k} \cdot f(T) \\
&\overset{IH}{\geq} \left(1 - \frac{1}{k}\right)(1 - (1 - \frac{1}{k})^i) \cdot f(T) + \frac{1}{k} \cdot f(T) \\
&= (1 - \frac{1}{k}) - (1 - \frac{1}{k})^{i+1} + \frac{1}{k} \cdot f(T) \\
&= (1 - (1 - \frac{1}{k})^{i+1}) \cdot f(T),
\end{align*}
\]

completing the inductive proof. Using the fact that \((1 - \frac{1}{k})^i \geq 1/e\) for all \( i \leq k \), we obtain that the algorithm is a \((1 - 1/e)\)-approximation. 

In order to prove Lemma 8.12, we first observe that monotonicity of \( f \) is trivial. We then focus on a simpler model (proposed by Goldenberg et al. [163, 162] in the context of viral marketing) to illustrate the concepts of the proof.

**Definition 8.14 (Independent Cascade Model)** Given a complete directed graph \( G \) with edge probabilities \( p_{e} \), we consider an infection model where once a node \( u \) becomes active, it infects a neighboring node \( v \) with probability \( p_{u,v} \). If the attempt succeeds, \( v \) becomes active; \( u \), however, does not get to try infecting \( v \) again.

We can observe that for each edge \( e \), we can decide ahead of time (randomly) if the activation attempt will succeed when/if it happens, with probability \( p_{e} \). We observe that, in the graph \( G \) of “successful edges”, the nodes active in the end are exactly the ones reachable from \( S \).

As a first step, we will show that the number of reachable nodes in a given graph \( G \) is a submodular function. We define \( f_{G}(S) \) to be the number of nodes reachable from \( S \) in \( G \), and prove

**Claim 8.15** \( f_{G} \) is submodular for all \( G \).

**Proof.** We need to show that \( f_{G}(S \cup \{x\}) - f_{G}(S) \geq f_{G}(S' \cup \{x\}) - f_{G}(S') \) whenever \( S \subseteq S' \). We write \( R_{G}(S) \) for the set of all nodes reachable from \( S \) in \( G \). Then, \( f_{G}(S) = |R_{G}(S)| \), and because any node reachable from \( S \cup \{x\} \), but not from \( S \), must be reachable from \( x \), we can observe that

\[
|R_{G}(S \cup \{x\})| - |R_{G}(S)| = |R_{G}(S \cup \{x\}) \setminus R_{G}(S)| = |R_{G}(\{x\}) \setminus R_{G}(S)|.
\]

By monotonicity of \( R_{G} \), we have that \( R_{G}(\{x\}) \setminus R_{G}(S) \supseteq R_{G}(\{x\}) \setminus R_{G}(S') \), so \( |R_{G}(\{x\}) \setminus R_{G}(S)| \geq |R_{G}(\{x\}) \setminus R_{G}(S')| \), which proves submodularity of \( f_{G} \).

From the submodularity for any fixed graph, we would like to derive that the same holds for the function \( f \), over randomly chosen graphs. To that end, we use the following useful way of rewriting a random variable’s expectation \( E[X] \): If \( \{\mathcal{E}_1, \mathcal{E}_2, \ldots, \mathcal{E}_m\} \) is a partition of the probability space \( \Omega \), and \( X \) is a random variable on \( \Omega \), then:

\[
E[X] = \sum_{i} \text{Prob}[\mathcal{E}_i] \cdot E[X | \mathcal{E}_i].
\]

In our case, the random variable \( X \) is the number of nodes reached from \( S \), and \( f(S) = E[X] \) is its expectation. Let \( G_1, \ldots, G_m \) be an enumeration of all graphs on \( n \) nodes (note that \( m \) is large). Notice that the events \( \mathcal{E}_i = [G = G_i] \) form a partition of the probability space we are considering. Thus, the above identity implies that \( f(S) = \sum_{i} \text{Prob}[G = G_i] \cdot f_{G_i}(S) \). Because all the coefficients \( \text{Prob}[G = G_i] \) are non-negative, and each \( f_{G_i} \) is submodular, the submodularity of \( f \) now follows from the following

**Fact 8.16** A non-negative linear combination of submodular functions is itself submodular.
Proof. Let \( \alpha_i \geq 0 \) for all \( i \), and \( f = \sum_i \alpha_i f_i \). If each \( f_i \) is submodular, then for all \( S \subseteq S' \),
\[
\begin{align*}
f(S \cup \{x\}) - f(S) &= \sum_i \alpha_i f_i(S \cup \{x\}) - \sum_i \alpha_i f_i(S) \\
&= \sum_i \alpha_i (f_i(S \cup \{x\}) - f_i(S)) \\
&\geq \sum_i \alpha_i (f_i(S' \cup \{x\}) - f_i(S')) \\
&= f(S' \cup \{x\}) - f(S').
\end{align*}
\]

Notice that we did not use a lot of properties of the Independent Cascade Process. In fact, we proved the following stronger lemma:

Lemma 8.17 If the activation process is such that we can define a distribution on graphs \( \{G_i\} \) such that \( f(S) \) equals the expected number of nodes reachable from \( S \) under the distribution, then \( f \) is submodular.

We would like to use this lemma to prove submodularity for the Linear Threshold Model that we started out with. In order to do that, we need to come up with a distribution over graphs such that the requirements of the lemma are met. In this case, it is not as obvious which distribution to choose, but we will be able to show that the following model works:

Definition 8.18 (Random Graph Model) Each node \( v \) has at most one incoming edge which emanates from \( u \) with probability \( w(u,v) \), the weight of the influence of \( u \) on \( v \). (Recall that \( \sum_u w(u,v) \leq 1 \) for all \( v \).) With probability \( 1 - \sum_u w(u,v) \), \( v \) has no incoming edge.

Claim 8.19 This model is equivalent to the Threshold Model, in the sense that the expected number of nodes reached is the same.

Proof. We will prove by induction on \( t \) that for each time step \( t \geq 0 \), and any node sets \( T \subseteq T' \), the probability that exactly \( T \) is active at time \( t \) and \( T' \) at time \( t+1 \) is the same in both the threshold and the random graph processes.

In the base case \( t = 0 \), the probability is 1 for the pair \((\emptyset, S)\) for both processes, and 0 for all other pairs, because both processes start with only the selected set \( S \) active.

For the inductive step, assume that \( T \subseteq T' \) are the active sets at time \( t - 1 \) and \( t \), and consider some node \( v \notin T' \). We investigate the probability that \( v \) becomes active at time \( t + 1 \) in either process.

In the threshold model, because \( v \) was not active at time \( t \), we know that \( \theta_v \geq \sum_{u \in T} w(u,v) \). However, subject to that, \( \theta_v \) is still uniformly random, so by the Principle of Deferred Decisions, we can re-choose \( \theta_v \) uniformly at random from the interval \((\sum_{u \in T} w(u,v), 1] \).

\( v \) becomes active at time \( t+1 \) iff \( \theta_v \leq \sum_{u \in T} w(u,v) \). This happens with probability, \( \frac{\sum_{u \in T'} w(u,v)}{1 - \sum_{u \in T} w(u,v)} \).

Now, let us look at the probability that \( v \notin T' \) becomes active at time \( t+1 \) in the random graph process. Because \( v \) is not active at time \( t \), we know that \( v \)'s incoming edge, if any, does not come from \( T \). Thus, \( v \) becomes active at time \( t + 1 \) if the edge comes from \( T' \setminus T \).

The probability that \( v \)'s edge does not come from \( T \) is \( 1 - \sum_{u \in T} w(u,v) \), and the probability that it comes from \( T' \setminus T \) is \( \sum_{u \in T' \setminus T} w(u,v) \). Hence, the conditional probability of \( v \) becoming active is \( \frac{\sum_{u \in T' \setminus T} w(u,v)}{1 - \sum_{u \in T} w(u,v)} \).

So for any individual nodes, the probability of activation at time \( t+1 \) is the same under both processes. As these decisions are independent in both processes (thresholds resp. edges are chosen independently), the probability that exactly all nodes from \( W \) become active is the same under both processes for any set \( W \).

\[
\begin{align*}
\text{Prob}[\{T', T''\} \text{ active at time } (t, t+1)] &= \sum_{T'} \text{Prob}[\langle T, T'\rangle \text{ active at time } (t-1, t)] \cdot \text{Prob}[\text{exactly } T'' \setminus T' \text{ becomes active } | \text{ } T' \text{ active at time } t]
\end{align*}
\]
is the same for both. Hence, by induction, the processes behave the same, and in particular, the final expected number of activated nodes is the same.

Notice that combining all of the above, we have established Lemma 8.12, i.e., the fact that \( f \) is submodular.

It remains to show how to find, in polynomial time, a node giving largest marginal gain at any point. In fact, it is now known whether or how one can find such a node. However, we can establish a slightly weaker fact, which is sufficient for our purposes. As stated in Lemma 8.11, we can find, in polynomial time, a \((1 - \epsilon)\)-approximately best node to add.

**Proof of Lemma 8.11.** First, we describe an algorithm for estimating the marginal gain of adding a node. If we can get sufficiently accurate estimates for each node, then a choice based on these estimates will give us an approximately best node.

The algorithm is simply to simulate the random process multiple times, and take the average of the number of additional nodes reached over all these simulations. This average will certainly in the limit converge to the expectation, but the question is how quickly, i.e., how often is “multiple” times? If some “gain values” had very low probability of occurring, but were extremely high, then we may need a large number of simulations to get a good estimate. Fortunately, in our case, the values we sample have a natural upper bound of \( n \), the number of nodes. We will see that the number of iterations will not need to be very high as a result.

Formally, we will use the Chernoff-Hoeffding Bound, as given by the following theorem.

**Theorem 8.20 (Chernoff-Hoeffding Bound [182, 248])** If \( X_1, \ldots, X_m \) are independent random variables with \( 0 \leq X_i \leq b_i \) for all \( i \), and \( X = \sum_i X_i \), and \( \mu = E[X] \), then for all \( \Delta \geq 0 \),

\[
\Pr[|X - \mu| \geq \Delta] \leq 2e^{-\frac{\Delta^2}{2m}}.
\]

In our case, we let \( X_i \) be the outcome of the \( i \)th simulation. Thus, \( 0 \leq X_i \leq n \) for all \( i \), so \( b_i = n \). Also, because at least one node is active in each outcome (the start node itself), we have that \( X_i \geq 1 \), and thus \( \mu \geq m \). Choosing \( \Delta = \epsilon \mu \), we obtain that

\[
\Pr[|X - \mu| \geq \epsilon \mu] \leq \Pr[|X - \mu| \geq cm] \leq 2e^{-\frac{(cm)^2}{m \epsilon^2}} = 2e^{-\frac{c^2m}{n \epsilon^2}}.
\]

Thus, if we are aiming for a \((1 \pm \epsilon)\)-approximation with probability at least \( 1 - \delta \), it is sufficient to require that \( 2e^{-\frac{c^2m}{n \epsilon^2}} \leq \delta \). Solving for \( m \) gives that \( m > \frac{n^2}{\epsilon^2} \log \frac{2}{\delta} \) simulations are sufficient. By a Union Bound over all (at most \( n \)) iterations of the greedy algorithm, and all \( n \) nodes in each iteration, all simulations have error at most \( \epsilon \) with probability at least \( 1 - n^2 \delta \).

If we want, for instance, success probability at least \( 1 - \frac{1}{n^2} \), we can choose \( \delta = \frac{1}{n^2} \), and then, we need to run \( O\left(\frac{n^2}{\epsilon^2} \log n\right) \) iterations to get a \((1 \pm \epsilon)\) accurate estimate.

It still remains to verify that we actually obtain a close to best node when picking a node based on \( \epsilon \)-estimates. In the worst case, the picked node’s gain is over-estimated by a factor of \((1 + \epsilon)\), while the true best node’s gain is under-estimated by a factor of \((1 - \epsilon)\). But because the picked node appeared to be better, its gain must have been within a factor of \( \frac{(1 - \epsilon)(1 + \epsilon)}{1 + \epsilon} \geq 1 - 2\epsilon \) of the best node’s. So we have a \((1 - 2\epsilon)\)-approximate best node in each iteration, in polynomial time.

By essentially mimicking our earlier proof of the Nemhauser-Wolsey-Fischer theorem, we can show the following stronger version:

**Theorem 8.21** If we choose a \((1 - \epsilon)\)-approximate best node in each iteration of the greedy algorithm, then we get a \( 1 - \frac{1}{\epsilon} - \epsilon' \) approximation algorithm, where \( \epsilon' \to 0 \) as \( \epsilon \to 0 \), polynomially fast.
8.7 Further Reading

We just showed that two natural models for influence propagation — the independent cascade model of Goldenberg et al. [163, 162] and the linear threshold model of Granovetter [170] — lead to submodular overall influence functions, and thus a $1 - 1/e$ approximation algorithm. A natural question is then which other models permit such approximations (or other approximation guarantees).

Several recent papers make progress in this direction. [210] shows that a generalization of the cascade model still leads to submodularity. Here, the probability of activation success $p_{u,v}(A)$ along an edge $(u,v)$ may depend on which other nodes $A$ have previously tried (and failed) to activate $v$. However, the probability $p_{u,v}(A)$ must be non-increasing in $A$, i.e., the marginal probability of success decreases. [210] shows that this condition is sufficient to yield submodularity. Interestingly, it is also shown that for this model, there may be no graph distribution yielding the same expected number of active nodes. Hence, [210] uses a different technique for proving submodularity.

The result from [210] is further generalized in a beautiful recent result of Mossel and Roch [265]. They consider a generalization of the threshold model, where each node $v$ has a local threshold function $f_v(A)$, and becomes active when $f_v(A) \geq \theta_v$, where $A$ is the set of previously active nodes. Thus, the linear threshold model is the special case where $f_v(A) = \sum_{u \in A} w_{u,v}$. [265] resolves a conjecture of [209] by showing that if each $f_v$ is non-negative, monotone increasing, and submodular, then so is the overall activation function. The proof uses a significantly more intricate coupling in the style of our proof of Claim 8.19.

An interesting question raised by [210] is which influence models do allow a proof of submodularity via an equivalent distribution of graphs. This question is answered in a recent paper by Salek et al. [305]. For a function $g$ defined on sets, let $g_v$ denote the “discrete derivative” $g_v(S) := g(S \cup \{v\}) - g(S)$, and inductively define $g_{T_v}(S) := g_T(S \cup \{v\}) - g_T(S)$. (It is not hard to see that this definition does not depend on which element of a set is chosen first.) [305] shows that there is an equivalent distribution over graphs if and only if the following three conditions hold simultaneously:

1. $g$ itself is non-negative.
2. All discrete derivatives $g_T$ with $|T|$ odd are non-negative.
3. All discrete derivatives $g_T$ with $|T| > 0$ even are non-positive.

A somewhat different influence model was considered by Even-Dar and Shapira [132]. They consider the voter model [92, 184]: in each iteration, each node $v$ simultaneously chooses a random neighbor according to some distribution (e.g., uniformly), and copies the state of that neighbor. In the limit, for a connected graph, this will result in all individuals adopting the same behavior. By showing equivalence to a Markov Chain, [132] shows that the influences of nodes factor, and therefore, a simple greedy algorithm is optimal. Even if different nodes have different activation costs, there is an FPTAS for the problem, based on the FPTAS for Knapsack [329].

As we mentioned, when thresholds at nodes are fixed, the problem becomes harder to approximate. In face even if all nodes have the same constant threshold $\tau \geq 2$, the problem of finding a minimum-size set of nodes that eventually activates the entire graph is NP-complete, a fact shown for $\tau \geq 3$ by Dreyer and Roberts [197] and for $\tau = 2$ (as well as majority and some other rules) by Chen [82]. In fact, [82] showed approximation hardness results for this version of the problem.

A direction which has received a fair share of attention recently is multiple competing influences. The motivation here is viral marketing, when different companies want to use word-of-mouth in the same social network to promote competing products. Each node can only adopt (and recommend) one product, and the main question is how to deal with the possible timing issues and create a well-defined and natural model. Probably the first paper offering explicit analysis of such games is one by Dubey et al. [122], which studies equilibria of such competition in the simple linear model of [301]. Subsequently, several other papers considered the aspect of competition. Bharathi et al. [42] augmented the cascade model with a timing component, which leads to a simple tie-breaking for choosing which product a node adapts. One of the main results of [42] is that the benefit of a set for the last player is still a submodular function; this also implies that the “Price of Competition” is at most a factor of 2, i.e., regardless of how many companies are
competing, the expected reach of these influences combined is at least half as much as if they all pooled and coordinated their resources. Karnes et al. [72] also consider the problem of the last player in this setting, but without a timing component. Instead, they propose two tie breaking rules: one distance-based (with an additional distance notion introduced for the problem), and one by uniform random choice from among all active neighbors.

The extent of diffusions (or epidemic diseases) through a social network has also received attention in the physics and economics communities (e.g., [191, 192, 272, 286, 287, 345]). There, the assumption is usually (implicitly) that the social network is uniformly random subject to its degree distribution, and that the initial set of infected individuals is also random. The interest in this line of work is then in “tipping point” or “phase transition” behaviors, i.e., at what settings of parameters does the reach of the influence go from nearly negligible to significant.

As we saw implicitly in Section 8.4, there can be multiple “equilibrium” states of the system. In particular, it is possible for some of the nodes to end up active, and others inactive. In terms of technology, this means that some are still using the old (inferior) technology, while others have switched to the new and superior one. This coexistence disappears if we modify the model as follows: with small probability, nodes choose the wrong response to their neighbors’ choices. Thus, sometimes, a node will become active even though it would be superior (in terms of the payoffs defined in Section 8.4) to remain inactive, or vice versa. In this case, the system always converges to the state in which all nodes use the superior technology [200]. An interesting question is then how quickly this convergence happens. Following up on work of Ellison [125], which analyzed the special case of complete graphs and $k$-nearest neighbor graphs on the line or cycle (for $n \gg k$), Montanari and Saberi [262] showed that the convergence time is asymptotically characterized by the tilted cutwidth of the graph, a certain isoperimetric property. In particular, graphs that expand well tend to have slow convergence, while globally poorly expanding graph have fast convergence. Note that this stands in stark contrast to the spread of diseases or models like the cascade model, where expansion tends to imply fast spread.

In our optimization problem so far, we have assumed that the sole goal is to maximize the number of individuals adopting the innovation. In a more realistic viral marketing setting, these adoptions will correspond to sales of an item, and the company will have the ability to set different prices at different times, or even for different consumers. This version has been studied by Hartline et al. [176]. They show that the following “Influence and Exploit” strategies are within a factor $\frac{1}{3}$ of optimal: first, give the product for free to a most influential subset; then, choose prices for the remaining bidders in random order to maximize revenue, ignoring their network effects. As a second step, they show that the resulting problem of selecting the most influential set leads to a submodular objective function, though not necessarily a monotone one. Using an 0.4-approximation algorithm of Feige et al. [140], they thus obtain a constant-factor approximation algorithm. A slightly modified model was studied by Arthur et al. [18]: in their model, the seller cannot choose arbitrarily whom to offer the product to. Instead, the product spreads by recommendations, but the seller can decide on a price for each individual once a recommendation occurs. In this model, Arthur et al. prove competitive guarantees based on the type of buyer model.

In our derivation of activation thresholds in Section 8.4, we assumed implicitly that each node adopted only one technology. Instead, a node $v$ may have the choice of adopting both the old and new technologies, at an extra cost of $d_v\phi$, for some parameter $\phi$ (where $d_v$ is $v$’s degree). In return, $v$ gets communication benefits of $1 - q$ or $q$ for all neighbors. For instance, a person could learn multiple languages, maintain multiple messaging softwares, or multiple data formats, such as CDs and MP3s. Immorlica et al. [188] considered this question. In their most basic model, they assume $a = b = 0$, i.e., no communication is possible with differing technologies. They investigate the regions for the parameters $\phi, q$ which allow the new technology to spread through the network. Interestingly, this region is not always convex. They then consider the generalization to partial compatibility, i.e., values $\phi, q > 0$, and its impact on the diffusion of the new innovation, again noticing non-monotonic behavior.

The technique we used in the proof of Theorem 8.10 is quite general: prove that the function to be optimized is non-negative, monotone and submodular; then, the greedy algorithm will be a $1 - 1/e$ approximation. This technique has been applied in a number of scenarios such as selecting sensors, variables to observe, or samples to draw [104, 224, 225, 226, 239]. In different scenarios, we may be interested in some-
what different objectives or constraints. For instance, we saw above that when sales prices are taken into account, the natural objective function is submodular, but not monotone. As another examples, Krause et al. [227] show that a sensor placement and scheduling problem for coverage can be expressed as maximizing the minimum of multiple submodular functions \( f_i(A_i) \), where the sets \( A_i \) must be disjoint, and their union satisfies a size constraint. [227], among others, gives a 1/6-approximation algorithm for this problem.

The frequency with which submodular objective functions appear in many settings has led to a lot of work on approximating them with various constraints. In a classic paper, Iwata et al. [189] show that for any submodular function \( f \), finding a set \( S \) minimizing \( f(S) \) can be achieved in polynomial time. (Notice, however, that we did not include any size constraints on \( S \).) For the maximization problem, there has been a lot of recent interest in generalizing the result of Nemhauser et al. [102, 268]. Sviridenko [320] shows how to obtain a \( 1 - 1/e \) approximation for monotone submodular functions even when different elements have different inclusion costs, and a total budget is given. Calinescu et al. [70, 331] give a beautiful \( 1 - 1/e \) approximation algorithm for the case when the solution is constrained to be from a given matroid\(^2\). If the solution is required to belong to \( k \) different matroids, then a recent result of Lee et al. [236] shows how to approximate the best solution to within a factor of \( 1/(k + \epsilon) \) for monotone functions.

For non-monotone submodular functions, we previously mentioned the 0.4-approximation algorithm of Feige et al. [140]. Lee et al. [235] consider the more general constraint of restricting solutions to \( k \) different matroids, and give a \( \frac{1}{k+2+1/k+\epsilon} \) approximation for this problem, as well as a \( 1/5 - \epsilon \) approximation when solutions must meet \( k \) Knapsack constraints. Vondrák [332] provides some lower bounds on how well these problems can be approximated.

In order to maximize the diffusion of an innovation in a network, it is necessary to first learn the influence that individuals have over each other. So far, we have assumed that these data are given precisely, which is certainly unrealistic. Indeed, they must be learned first, presumably from past observed data. This question has recently been investigated by several papers, with different objectives and results [15, 166, 304]. The underlying theme is to use a Maximum Likelihood Estimation of the network, with different constraints on what edges or probabilities are allowed.

---

\(^2\)A matroid is a downward-closed set system with the following exchange property: If \( S, T \) are two sets in the set system with \( |S| < |T| \), then there is an element \( u \in T \setminus S \) such that \( S \cup \{u\} \) is also in the matroid. In other words, elements of sets can be exchanged one by one.
Chapter 9

Epidemic Algorithms

In the previous chapter, we considered examples of processes that naturally spread through networks in an epidemic fashion, such as information or behaviors, or — of course — epidemic diseases. One intuitive observation about such spreading processes is that they are very resilient: even if a few individuals are immune or removed from the network, the epidemic tends to reach a similarly large share of the network. This resilience, while undesirable for actual epidemics, is very desirable for communication and computation in computer networks: even if a few nodes are faulty, we would like for information to diffuse to the rest of the network, and for computation to succeed. The first to propose simulating the behavior of epidemics as a principle for systems design were Demers et al. [113], in the context of maintaining replicated databases. Their idea was that updates to the databases would not be distributed directly to all other copies, but passed around like a rumor, with nodes randomly choosing others to exchange updates with.

Since the appearance of [113], there has been a large body of work on epidemic algorithms as a primitive for computer networks. Among the first to propose a system based on these approaches were van Renesse et al. with the design of Astrolabe [327, 328]. Here, we will investigate some of the analysis underlying two key primitives of distributed computation: spreading a single message to everyone, and computing the average of the values held by all nodes in the network. These lie at the core of many more complex network tasks.

9.1 Spreading One Message

To get a feel for the benefits of epidemic algorithms, let us begin by considering an obvious alternative, also used in many systems. To disseminate a message to all nodes, the nodes could build (a priori) a tree, and then forward the message along the edges of the tree. So long as the height of the tree is logarithmic, and the degrees bounded by a constant (e.g., for a complete binary tree), the message will reach all nodes in $O(\log n)$ steps. The best completion time (time by which the last node receives the message) is achieved by a binomial tree. Notice that $\Omega(\log n)$ is also an obvious lower bound: in each round, the number of nodes having the information can at most double.

The problem with any tree is that it is not fault-tolerant. If one node fails, then none of the nodes in the subtree rooted at it will receive the message. In order to achieve more fault-tolerance, we will need more redundancy in messages sent. One could make the graph 2-connected, or more highly connected. But in order to achieve higher fault-tolerance, the graph structures will have to become more complicated, and thus difficult to compute and maintain. Instead, we will see that simple randomized protocols based on the behavior of epidemics achieve essentially the same time bounds, and are naturally fault-tolerant. Such algorithms are usually called epidemic algorithms or gossip algorithms, and are based on information exchanges between random node pairs.

For the time being, we assume that every pair of nodes can in principle communicate, i.e., that the communication graph is complete. (We will discuss some pointers to literature avoiding this assumption at the end of the chapter.) In all versions of gossip protocols we consider here, we assume that there are
synchronous rounds, and in each round, each node calls a random other node and exchanges information. These exchanges can be divided into three categories:

1. In push gossip, each (informed) sender randomly picks a partner and sends a message.
2. In pull gossip, each node picks a random partner, and receives a message from the partner (if it has one).
3. In push & pull gossip, each node picks a random partner, and both forwards and receives a message.

Naturally, push & pull gossip makes most sense when multiple messages are being exchanged, as otherwise, no informed node ever needs to pull, and no uninformed node can ever push. For instance, exchanging information in both directions is useful in synchronizing a replicated distributed database periodically, the original application of [113].

In deciding whether to use gossip as a primitive, an important consideration is its speed: how quickly do all nodes receive the message? We expect the time to be roughly $O(\log n)$, as the number of nodes having the message roughly doubles until most nodes already have the message.

**Proposition 9.1** With high probability, it takes $O(\log n)$ rounds to get a message to everyone.

**Proof.** We first analyze the behavior until more than $n/3$ nodes have the message. So consider a time $t$ when at most $m \leq n/2$ nodes have the message. Then, each of the $m$ messages sent has probability at least $2/3$ of reaching a previously uninformed node, so in expectation, at least $2m/3$ sent messages reach previously uninformed nodes.

Unluckily, this does not quite guarantee that $2m/3$ new nodes will be informed, as some of these messages will reach the same node. We want to upper-bound how frequently this happens. For any given pair of messages, the probability that they both go to the same uninformed node is $O(1/2^x)$, where $x$ is the number of newly infected nodes.

Let the random variable $X_t$ be the number of informed nodes at time $t$. By our arguments above, $X_0 = 1$, and $E[X_{t+1} | X_t = m] = \frac{2m}{3}$. By induction, $E[X_t] \geq \frac{17}{12}m$. By increasing the constant in the $O(\log n)$ time bound, we can improve the probability to $1 - \frac{1}{n^c}$, for any $c$.

A more careful analysis is performed in [151, 290], where the precise constants and lower-order terms are worked out to guarantee a high-probability bound.

In the proof of Proposition 9.1, notice the following: From the point at which $n/3 = \Omega(n)$ nodes have the message until the point when all nodes have the message, there are $\Omega(\log n)$ rounds, and in each such round, each informed node sends the message to some other node, for a total of $\Omega(n \log n)$ messages. On the other hand, using any tree would only require sending $O(n)$ messages. Thus, gossip achieves its fault-tolerance with high redundancy: by using $\Omega(\log n)$ time as many messages as necessary.

Naturally, one wonders whether such an amount of redundancy is inherent in any fault-tolerant approach, or in any gossip-based approach. Karp et al. [203] give a partial answer to this question. They show that any protocol based on uniform gossip must send $\omega(n)$ messages, and, under additional restrictions on the protocol, $\Omega(n \log n)$ messages. Conversely, they analyze more carefully the “push-pull” version of gossip, and show that if the terminating point is chosen carefully, only $O(n \log \log n)$ messages are sent.

**Theorem 9.2 ([203])** If push-pull is run for $\log_3 n + O(\log \log n)$ rounds, then, w.h.p., all nodes have the message, and, in addition, the total number of messages sent is $O(n \log \log n)$.  

92
Proof. Solely for the purposes of analysis, we divide the algorithm into four phases.

1. **Startup phase**: This phase lasts until \( \log^4 n \) nodes have the message with high probability (probability at least \( 1 - n^{-c} \) for some \( c \)). We observe that within \( O(1) \) rounds, the number of informed nodes doubles with high probability (because each informed node will have called a previously uninformed one), and so the entire phase takes \( O(\log \log n) \) rounds for the goal of achieving \( \log^4 n \) informed nodes.

2. **Exponential growth phase**: This phase lasts until \( \frac{n}{\log n} \) nodes are informed. We write \( S_t \) for the number of nodes having the message at time \( t - 1 \), and \( m_t \) the number of messages that were sent in round \( t \). Because in expectation, each informed node calls one node, and is called by one, we have that \( E[m_t] = 2S_{t-1} \).

   In fact, \( m_t \) is sharply concentrated, i.e., unlikely to deviate far from its expectation. To see this, we let

   \[
   X_{uv} = \begin{cases} 
   1 & \text{if node } u \text{ calls } v \text{ in round } t \\
   0 & \text{otherwise.} 
   \end{cases}
   \]

   Then, \( m_t = \sum_{u \in S_{t-1}} \sum_v (X_{uv} + X_{vu}) \). Because \( m_t \) is a sum of negatively dependent \( \{0,1\} \) variables, (a different version of) Chernoff Bounds shows that \( m_t \) is sharply concentrated around its mean.

   **Theorem 9.3 (Chernoff Bounds [84, 266])** If \( X_1, X_2, \ldots, X_n \) are independent (or negatively correlated) \( \{0,1\} \) random variables, with \( X = \sum_i X_i \), \( \mu \geq E[X] = \sum_i \text{Prob}[X_i = 1] \), and \( \delta > 0 \), then

   \[
   \text{Prob}[X > (1 + \delta)\mu] < \left( \frac{e^\delta}{(1+\delta)^{1+\delta}} \right)^\mu.
   \]

   Similarly, if \( \mu \leq E[X] \), then

   \[
   \text{Prob}[X < (1 - \delta)\mu] < e^{-\frac{\delta^2}{2}\mu}.
   \]

   Some of these \( m_t \) messages will not succeed in informing new nodes. The reasons could be two-fold: either they reach an already informed node, or multiple messages reach the same uninformed node.

   The probability that a given message reaches an already informed node is \( S_{t-1}/n \), and the probability of a collision with any other message is at most \( m_t/n \). Hence, the expected number of informed nodes after \( t \) iterations is at least

   \[
   E[S_t] \geq S_{t-1} + m_t(1 - m_t/n - S_{t-1}/n) \\
   = S_{t-1}(1 + (2 \pm o(1/\log n))(1 - O(1/\log n))) \\
   = S_{t-1}(3 \pm O(1/\log n)).
   \]

   In the second step here, we used the fact that \( S_{t-1} \leq n/\log n \). Again, we can write the number of successful messages (those that inform a new node) as a sum of \( \{0,1\} \) random variables, and apply a Chernoff Bound, showing that the value of \( S_t \) will be sharply concentrated (to within a multiplicative \( 1 \pm 1/\log n \)) around its expectation. When all iterations have outcomes close to the expectation (which happens with high probability by a union bound over all these iterations), we reach \( n/\log n \) nodes in \( \log_3 n + O(\log \log n) \) rounds.

3. **Quadratic Shrinking Phase**: Once \( n/\log n \) nodes are reached, each uninformed node has sufficiently high probability of reaching an informed node in its Pull attempt, so we will ignore Push transmissions, and only use Pull in the analysis. We write \( U_t \) for the number of uninformed nodes after \( t \) iterations. Then, for any node that is uninformed in iteration \( t - 1 \) and makes a uniformly random Pull attempt, the probability that it stays uninformed is \( U_{t-1}/n \). As a result, the expected fraction of uninformed nodes after the \( t \text{th} \) iteration is \( E[U_t/n] \leq (U_{t-1}/n)^2 \), leading to quadratic shrinking. Again, the calls for different nodes are independent, so we can write \( U_t \) as a sum of independent 0-1 random variables, and so long as the number of uninformed nodes \( U_{t-1} \) is large enough \( (U_{t-1} \geq \sqrt{n \log^4 n}) \), quadratic shrinking will occur with high probability by Chernoff Bounds.

93
Because the number of uninformed nodes thus shrinks doubly exponentially, within $O(\log \log n)$ rounds, the number of remaining uninformed nodes is at most $\sqrt{n} \log^4 n$, with high probability.

4. **Finishing**: Finally, once all but $\sqrt{n} \log^4 n$ of the nodes are informed, each node has probability $1 - \frac{\log^4 n}{\sqrt{n}}$ of being informed during any one Pull call it makes, and successive rounds are independent. Hence, after three rounds, the probability of any one node not being informed is at most $\frac{\log^4 n}{\sqrt{n}}$, and by a union bound, all nodes are informed after the additional three rounds with probability at least $\frac{\log^4 n}{\sqrt{n}}$.

Taking all this together, we have proved that with high probability, all nodes will learn the message within $\log_3 n + O(\log \log n)$ rounds. To analyze the total number of messages sent, we notice that the startup, quadratic shrinking, and finishing phases take only $O(\log \log n)$ rounds total, so the total number of messages sent in those rounds is $O(n \log \log n)$. During the exponential growth phase, at most $n/\log n$ nodes are active, so no more than $n/\log n$ messages are sent each round, thus no more than $O(n)$ total. As a result, the total message complexity is $O(n \log \log n)$.

We should notice that the protocol, as presented here, is very vulnerable to faults. It needs to be run for exactly $\log_3 n + O(\log \log n)$ rounds. If the number of rounds is $\Omega((1 + \varepsilon) \log_3 n)$, then $\Omega(\varepsilon n \log \log n)$ messages are sent during those extra $\varepsilon \log_3 n$ rounds, violating the goal of $O(n \log \log n)$ messages. Conversely, if the number of rounds is $O((1 - \varepsilon) \log_3 n)$, then with high probability, some nodes will remain uninformed. Karp et al. [203] show how a more sophisticated “Median Counter” algorithm achieves the same kind of guarantee while being more resilient to failures.

A second question is how low we can push the number of messages sent. Of course, if we do away with gossip, and build a tree instead, then $n - 1$ messages suffice to inform all nodes. Even using uniform gossip, we can get away with $n - 1$ messages if we are willing to sacrifice the completion time. A protocol achieving this would be one where the initial owner of the message is the only one to transmit it, and does so whenever he calls (or is called by) an uninformed node. This turns the problem into a coupon collector problem: the message has reached everyone if everyone has called (or been called by) the initial message holder. This takes $\Theta(n \log n)$ rounds with high probability [266].

Obviously, this type of protocol goes completely against the idea of using the network in spreading a message. Karp et al. show that under “reasonable” restrictions on the protocol, $\omega(n)$ messages are necessary. The first result concerns address-oblivious protocols: in each round, a node must decide whether to transmit to its communication partner(s) without knowing their identity. In particular, nodes do not know if the communication partner already has the message. Notice that the Push-Pull protocol we analyzed above is address-oblivious.

**Theorem 9.4 ([203])** Any address-oblivious protocol has to send at least $\Omega(n \log \log n)$ messages to inform all nodes.

Even without the restriction, lower bounds can be proved, trading off the completion time and message complexity:

**Theorem 9.5 ([203])** Any protocol reaching all but $O(1)$ of the nodes in $O(\log n)$ rounds, and using uniform gossip, must send at least $\omega(n)$ messages.

An interesting question along the lines of reducing the number of messages sent is how many nodes are informed, and how many rounds it takes, when each node forwards the message $k$ times after receiving it, for some constant $k$ (such as $k = 2$).

### 9.2 Averaging using Gossip

So far, we have focused on gossip only as a means to disseminate a single message to all nodes. In reality, this is a simplistic abstraction. When each node has a message that is to be shared with all other nodes, we could
consider a protocol in which each node always forwards all messages it holds. In that case, we are essentially “super-imposing” the simple gossip protocol for all source messages, so all nodes will learn all messages in $O(\log n)$ rounds with high probability. Unfortunately, the message complexity will be $\Omega(n^2 \log \log n)$. However, if we assume that messages are atomic, i.e., cannot be aggregated, then there is not much beyond this bound that can be done — even dissemination along a tree will take $O(n^2)$ messages total.

On the other hand, we may be able to send significantly less data around if the messages can be aggregated. Instead of forwarding each message individually, we could then perform aggregation within the network. For instance, assume that each node $i$ just holds one number $x_i$ (the number of files, or available memory, or a temperature measurement), and we want to compute the average $\bar{x} = \frac{1}{n} \sum_i x_i$ of these numbers. Using a tree, values could be added up (and nodes counted) on the way to the root; the root could then compute the average, and distribute it back to all nodes using the tree. As we discussed above, using a tree is not a very fault-tolerant approach, so we want to use gossip instead.

Our gossip-based algorithm is very simple. Each node $i$ maintains only two values, a sum $s_i$, and a weight $w_i$. The sum is initialized to $s_i := x_i$, and the weight to $w_i := 1$. Then, in each round, each node executes the following protocol Push-Sum [206]:

Algorithm 6 Push-Sum
1: Send the pair $(s_i/2, w_j/2)$ to yourself and a uniformly randomly chosen other node.
2: Let $J_i$ be the set of all nodes that have sent a message to $i$ in this round.
3: The new values are $s'_i := s_i/2 + \sum_{j \in J_i} s_j/2$ and $w'_i := w_i/2 + \sum_{j \in J_i} w_j/2$.
4: Keep track of $s_i/w_i$ as approximation of the average.

One useful fact that we can notice right away about this protocol is its property of mass conservation: at any point of the execution, we always have that $\sum_i s_i = \sum_i x_i$, and $\sum_i w_i = n$. The sums and weights get redistributed, but not truly changed. We will prove the following theorem about the convergence of Push-Sum to the true average.

**Theorem 9.6** If all $x_i$ are non-negative, then within $O(\log n + \log \frac{1}{\delta} + \log \frac{1}{\epsilon})$ rounds, all estimates $s_i/w_i$ are within $(1 \pm \epsilon)$ of the true average $\bar{x} = \frac{1}{n} \sum_i x_i$, with probability at least $1 - \delta$.

In order to analyze this protocol, we will track what “share” of each node $j$’s number $x_j$ is currently contributing to each node’s sum $s_i$. So we look at a vector $\vec{v}_i$, in which the $j^{th}$ component $v_{i,j}$ denotes the fraction of node $j$’s value that currently is part of node $i$’s sum. This means that initially, we have $v_{i,i} = 1$ for all $i$, and $v_{i,j} = 0$ for all $i \neq j$. In this view, our protocol can be expressed as follows:

Algorithm 7 Push-Vector
1: Send the vector $\frac{1}{2} \vec{v}_i$ to yourself and a uniformly randomly chosen other node.
2: Let $J_i$ be the set of all nodes that have sent a message to $i$ in this round.
3: The new vector is $\vec{v}'_i := \frac{1}{2} \vec{v}_i + \sum_{j \in J_i} \frac{1}{2} \vec{v}_j$.

This new version of the protocol traces the old version in the following sense (as is easy to observe, and can be proved by induction on the time steps of the protocol):

**Fact 9.7** At any time during the execution of the protocol, and for any node $i$, we have that $s_i = \sum_j v_{i,j} x_j$, and $w_i = \sum_j v_{i,j}$.

Thus, the estimate that node $i$ has is $\frac{\sum_j v_{i,j} x_j}{\sum_j v_{i,j}}$, and we want to show that this quantity converges exponentially fast to $\bar{x}$. One way that we could guarantee convergence would be if all $v_{i,j}$ were equal, and thus equal to $1/n$. In that case, we would have exactly the true average. However, this is clearly too much to hope for: by standard Balls-in-Bins analysis [266], even if they were all equal at some point, one node would be likely to receive many calls (up to $\Omega(\log n)$), and others no call at all, resulting in new values $\Omega(\frac{\log n}{n})$.
vs. \( \frac{1}{2n} \). However, upon closer inspection, we notice that we do not need quite such a strong condition. It would be enough if, for a fixed \( i \), all \( v_{i,j} \) were the same. So we do not need a node’s value to be equally distributed among all other nodes; all we need is that a node has equally sized shares of everyone else’s values.

This motivates studying how fast the vectors \( \vec{v}_i \) converge to multiples of the all-ones vector \( \vec{1} \). In order to talk about this convergence, we use \( v_{t,i,j}, \vec{v}_{t,i}, s_{t,i}, w_{t,i} \) etc. to denote the values after \( t \) iterations of Push-Vector. For ease of notation, we will use the fact that \( w_{t,i} = \sum_j v_{t,i,j} \). We then measure the convergence in terms of the error \( \Delta_{t,i} = \max_j |\frac{v_{t,i,j}}{w_{t,i}} - \frac{1}{n}| \). We will prove the following two parts, giving the theorem together:

**Lemma 9.8**

1. The \( \Delta_{t,i} \) converge to 0 exponentially fast.

2. When the \( \Delta_{t,i} \) are small, the estimate of the average is good.

**Proof.** We first prove the (easier) second part of the lemma. Assume that at some point in time \( t \), the errors \( \Delta_{t,i} \) for all \( i \) are at most \( \epsilon/n \). Then, for each \( i \),

\[
\frac{\left| \sum_j \frac{v_{t,i,j} x_j}{w_{t,i}} - \bar{x} \right|}{|\bar{x}|} = n \cdot \frac{\left| \sum_j \left( \frac{v_{t,i,j}}{w_{t,i}} - \frac{1}{n} \right) x_j \right|}{|\sum_j x_j|} 
\leq \frac{n}{|\sum_j x_j|} \cdot (\max_j \left| \frac{v_{t,i,j}}{w_{t,i}} - \frac{1}{n} \right|) \cdot \sum_j |x_j| 
\leq \epsilon,
\]

where we used the triangle inequality in the numerator for the second step, and the bound on \( \Delta_{t,i} \) for the expression in parentheses in the third step. (We also were allowed to cancel the sums over \( x_j \) values because all \( x_j \) were assumed to be non-negative.) Notice that once we prove exponential convergence below, the time it takes to converge to error \( \epsilon/n \) is only by an additive \( O(\log n) \) larger than to converge to \( \epsilon \), so we were free to choose a value of \( \epsilon/n \) here.

To prove the first part of the lemma, we study a potential function, which measures how “close” to converged the system is. We have seen such functions before, for instance in the proof of Theorem 8.7. Here, our potential function will be the sum of variances of the vectors \( \vec{v}_{t,i} \). Formally, we define

\[
\Phi_t = \sum_{i,j} (v_{t,i,j} - \frac{w_{t,i}}{n})^2.
\]

We will show that this potential function decreases exponentially fast, and that a small value for it implies good convergence.

**Lemma 9.9** The conditional expectation of \( \Phi_t \) satisfies

\[
E[\Phi_{t+1} \mid \Phi_t = \phi] = (\frac{1}{2} - \frac{1}{2n})\phi.
\]

**Proof.** Consider the values \( v_{t,i,j}, w_{t,i} \) etc. at time \( t \), and let \( f(i) \) denote the random node called by node \( i \) in round \( t \). Then, with all random choices known, node \( i \)'s new vector \( \vec{v}_i' \) and weight \( w_i' \) are

\[
\vec{v}_i' = \frac{1}{2} \vec{v}_i + \frac{1}{2} \sum_{k : f(k) = i} \vec{v}_k,
\]

\[
w_i' = \frac{1}{2} w_i + \frac{1}{2} \sum_{k : f(k) = i} w_k.
\]
Plugging these values into the new potential $\Phi_{t+1}$, we obtain that

$$\Phi_{t+1} = \sum_{i,j} \left( \frac{1}{2} (v_{i,j} - \frac{w_i}{n}) + \frac{1}{2} \sum_{k : f(k) = i} (v_{k,j} - \frac{w_k}{n}) \right)^2$$

$$= \frac{1}{4} \sum_{i,j} (v_{i,j} - \frac{w_i}{n})^2 + \frac{1}{4} \sum_{i,j : k : f(k) = i} (v_{k,j} - \frac{w_k}{n})^2 + \frac{1}{2} \sum_{i,j : k : f(k) = i} (v_{i,j} - \frac{w_i}{n})(v_{k,j} - \frac{w_k}{n})$$

$$+ \frac{1}{2} \sum_{i,j : k : k' \neq k', f(k') = i} (v_{k,j} - \frac{w_k}{n})(v_{k',j} - \frac{w_{k'}}{n}).$$

Noticing that in the second sum, the term $(v_{k,j} - \frac{k}{n})^2$ appears exactly once for each $k$ (namely for the particular $i$ with $f(k) = i$), we see that the first two sums are precisely equal to $\frac{1}{2} \Phi_t$. Similarly, we can simplify the fourth sum by noticing that each pair $k, k'$ with $f(k) = f(k')$ will appear for exactly one $i$. So we simplify

$$\Phi_{t+1} = \frac{1}{2} \Phi_t + \frac{1}{2} \sum_{i,j,k} (v_{i,j} - \frac{w_i}{n})(v_{k,j} - \frac{w_k}{n}) \cdot \left[ f(k) = i \right]$$

$$+ \frac{1}{2} \sum_{j,k,k' : k \neq k'} (v_{k,j} - \frac{w_k}{n})(v_{k',j} - \frac{w_{k'}}{n}) \cdot \left[ f(k) = f(k') \right].$$

Here, we are using Iverson’s convention [168]: $\left[ f(k) = i \right] := \begin{cases} 1 & \text{if } f(k) = i \\ 0 & \text{otherwise} \end{cases}$.

In this form, the expectation of $\Phi_{t+1}$ is not too difficult to evaluate: we can use linearity of expectation, and notice that the only terms actually depending on the random choices are $\left[ f(k) = i \right]$ and $\left[ f(k) = f(k') \right]$. As they are $\{0, 1\}$ random variables, their expectation is exactly equal to the probability of being 1, which can be easily seen to be $1/n$ for both. (For the first one, this is obvious; for the second one, notice that for any choice of $f(k)$, the probability that $f(k') = f(k)$ is $1/n$, so the same holds overall.) Substituting all of these, we obtain that

$$E[\Phi_{t+1} \mid \Phi_t = \phi] = \frac{1}{2} \phi + \frac{1}{2} \sum_{i,j,k} (v_{i,j} - \frac{w_i}{n})(v_{k,j} - \frac{w_k}{n}) \cdot \text{Prob}[f(k) = i]$$

$$+ \frac{1}{2} \sum_{j,k,k' : k \neq k'} (v_{k,j} - \frac{w_k}{n})(v_{k',j} - \frac{w_{k'}}{n}) \cdot \text{Prob}[f(k) = f(k')]$$

$$= \frac{1}{2} \phi + \frac{1}{2n} \sum_{i,j,k} (v_{i,j} - \frac{w_i}{n})(v_{k,j} - \frac{w_k}{n}) + \frac{1}{2n} \sum_{j,k,k' : k \neq k'} (v_{k,j} - \frac{w_k}{n})(v_{k',j} - \frac{w_{k'}}{n})$$

$$= \frac{1}{2} \phi + \frac{1}{2n} \sum_{i,j,k} (v_{i,j} - \frac{w_i}{n})(v_{k,j} - \frac{w_k}{n}) - \frac{1}{2n} \sum_{j,k} (v_{k,j} - \frac{w_k}{n})^2$$

$$= \left( \frac{1}{2} - \frac{1}{2n} \right) \phi + \frac{1}{n} \sum_j \sum_i (v_{i,j} - \frac{w_i}{n}) \left( \sum_k w_k \right)$$

$$= \left( \frac{1}{2} - \frac{1}{2n} \right) \phi.$$

In the last step, we used mass conservation, which implied that $\sum_i v_{i,j} = 1$, and $\sum_i w_i = n$, so that the second term actually became 0. Two steps earlier, we made the last sum run over all pairs $k, k'$, and subtracted out the ones we had added in. We also notice that at that point, both sums are equal, so we added them up to form the first one. In summary, this proves the lemma about the conditional expectation.

By applying the lemma repeatedly, and using that $\Phi_0 \leq n$, we obtain that $E[\Phi_t] \leq n \cdot 2^{-t}$. Thus, after $t = \log n + \log \frac{1}{\epsilon}$ rounds, the expected potential is $E[\Phi_t] \leq \epsilon$. Markov’s Inequality [266] states that for any
At least $w_i,t$ events considered, and dividing the earlier bound by $w_i,t$ each node must have weight at least 2 worst halved in each round, and similarly while a node simply "holds" the message, we know that at time $\tau$ more specifically, a theorem by Frieze and Grimmett [151], after $\tau$ small is not quite the error measure $\Delta_{t,i}$ received the message after $\tau$ and $w_i,t$ could potentially be quite small. This happens for instance when node $i$ has not received a message from anyone in a while (unlikely, but possible), or when it did receive a message, it was from another node that had not received a message in a while. At this point, we can leverage the earlier analysis of the dissemination of a single message.

Look at all nodes at time $t - \tau$. (Notice that our choice of $t$ implies that $t \geq \tau$ so we are allowed to do that.) At that point, at least one node $i$ had weight at least 1. Consider the experiment in which this node has a "message", and look at when each node $i$ receives the message. By our previous analysis, and, more specifically, a theorem by Frieze and Grimmett [151], after $\tau = 4 \log n + \log \frac{2}{\delta}$ steps, all nodes have received the message after $\tau$ steps with probability at least $1 - \delta/2$. Because the weight in a message is at worst halved in each round, and similarly while a node simply “holds” the message, we know that at time $\tau$, each node must have weight at least $2^{-\tau}$ with probability at least $1 - \delta/2$. Taking a union bound over both events considered, and dividing the earlier bound by $w_i,t \geq 2^{-\tau}$, we obtain that at time $t$, with probability at least $1 - \delta$, the errors are at most $\Delta_{t,i} = \frac{w_i,t}{w_i,t} - \frac{1}{n} \leq \epsilon$.

Finally, a simple inductive proof shows that once the error at all nodes drops below $\epsilon$, it will stay below $\epsilon$ at all nodes deterministically, so the quality of the estimate never gets worse.

To complete the proof, we need to verify how large our $t$ is exactly. Substituting the values of $\epsilon$ and $\tau$ into $t$, we see that it is $O(\log n + \log \frac{1}{\epsilon} + \log \frac{1}{\delta})$. This completes the proof.

### 9.3 Further Reading

A recent survey covering some of the extensive work on gossip was written by Shah [313].

In our analysis of the time to spread a single message to all nodes, we assumed that the communication graph was complete, i.e., that each node could communicate with each other node. Several papers have worked on generalizing the analysis to the case when there is a given graph structure restricting which node pairs can communicate.

For the simplest uniform gossip protocol (where each informed node simply calls and informs a uniformly random neighbor), Elsässer and Sauerwald [127] give degree-dependent lower bounds of $\Omega(\log n)$ (where the base of the logarithm depends on the degree), and an upper bound of $O(n \log n)$. This matches a lower bound for the star graph obtained via the coupon collector problem, and is thus tight in the worst case. Naturally, for specific classes of graph, much better bounds than $\Theta(n \log n)$ are possible, even for the simple randomized broadcast protocol. For instance, Feige et al. [141] give tighter bounds for random graphs and hypercubes. Chierichetti et al. [85] use expansion-based analysis techniques to prove an upper bound of $O(\log^2 n)$ for graphs obtained from preferential attachment (see Section 6.1). Mosk-Aoyama and Shah [264] prove an upper bound of essentially $O(\frac{\log n}{\delta})$ for the time to spread a single message in an arbitrary graph; here, $\Phi$ is the conductance of the matrix $M$ determining the call probabilities. For their result, the authors assume that $M$ is double stochastic. Chierichetti et al. [86] extend this result to uniform gossip (which leads to matrices that are not doubly stochastic), and obtain bounds polynomial in $\frac{\log n}{\Phi^2}$. This bound is based on a breakthrough result on graph sampling by Spielman and Teng [316].

For specific classes of graphs, the lower bounds of Karp et al. [203] can also be strengthened. If the protocol is address oblivious, then for a star, it is easy to see that a message needs at least $\Omega(n \log n)$ rounds to reach all nodes. For $G(n,p)$ random graphs with $p$ sufficiently large, Elsässer [126] shows a lower bound of $\Omega(n \log n / \log \log n)$ on the number of messages required to inform all nodes. Interestingly, address obliviousness is crucial here; if the nodes have just enough memory to remember the last three other nodes they called, Elsässer and Sauerwald [128] show that the total number of messages can be reduced to $O(n \log \log n)$.
So far, our main goal in this context was to spread the message(s) to all nodes quickly. In many natural settings, nodes are embedded geographically, and it is important that nodes close to the source obtain the message first. This may be relevant if the message is an alarm or warning, or if it contains information more useful locally. The problem of spreading a message to nodes at distance $d$ within time poly-logarithmic in $d$ (independently of $n$) is studied in [208]. There, it is shown that if nodes call each other with non-uniform probability, decreasing polynomially in the distance, then such a poly-logarithmic spreading time is achieved. Both the distribution and analysis bear some similarity to the analysis of greedy routing in Small-World Networks in Section 7.3.

As shown in [208], inversely polynomial gossip distributions can be used as a substrate for more complex types of in-network computation, such as nearest network resources. The problem of resource discovery in networks has also received attention in the context of Peer-to-Peer networks; many different algorithms have been proposed (e.g., [175, 230]). Further protocols that can be implemented with inverse polynomial distributions are discussed in [207], which also proves lower bounds on how well uniform gossip could approximate the same problems.

Analyzing the speed of the diffusion of a message through a network under inversely polynomial distributions in the distance is also closely related to a problem studied in mathematics and physics under the name long-range percolation [310]. Here, a graph is generated by embedding the nodes in a metric space (usually, the $D$-dimensional grid or torus), and generating a certain number of edges per node, with probability $d^{-r}$ of reaching distance $d$. Thus, the model is essentially identical to the small-world model discussed in Section 7.3. Using techniques similar to [208], Benjamini and Berger [34] prove poly-logarithmic bounds on the diameter of such graphs. Some of these bounds were subsequently improved by Coppersmith et al. [100]. See also the paper by Nguyen and Martel [282] for a discussion.

In transitioning from single-message gossip to gossip-based aggregation, we briefly discussed the issue of gossiping multiple messages (which was already discussed in the paper by Demers et al. [113]). If we assume that each connection can only transmit a single message (instead of the naive approach discussed before, where all messages are always forwarded), then even recovering an optimal time bound is non-trivial. Deb et al. [110] propose a protocol based on network coding and uniform gossip, where messages are “aggregated” by forming random linear combinations. From these linear combinations, it is then possible to reconstruct the initial messages. The advantage is that nodes can communicate just one aggregated message at a time, and still eventually collect all the necessary information. The number of messages required is $O(k + \sqrt{k} \log(k) \log(n))$. This bound was subsequently improved by Fernandess and Malkhi [143] who gave a protocol not using network coding, instead carefully choosing which message to forward in any time step. They achieve an optimal bound of $O(k + \log n)$ rounds.

In our discussions of gossip-based averaging in Section 9.2, we assumed that the communication graph is complete. The problem of averaging (or more general aggregation of data) is also important in networks with communication restricted to neighbors. Depending on the degree of nodes, we could continue to assume that nodes can only communicate with one neighbor, or relax the condition and let nodes communicate with all neighbors simultaneously. For the latter case, [206] shows that by averaging appropriately weighted combinations of the values of all neighbors, a variant of the Push-Sum protocol computes the average, in time equal to the mixing time of a Markov Chain on the graph. Thus, the fastest convergence is achieved by choosing the multipliers so as to have the Markov Chain mix as quickly as possible. This problem has been studied by Boyd et al. [56], who show how to use convex programming to find the fastest mixing Markov Chain for a given graph.

While [206] shows that the mixing time of the fastest Markov Chain on a graph is an upper bound on the time to average using gossip, it does not rule out that much faster protocols could exist. A lower bound is provided by Boyd et al. [58]. Under a slightly more restrictive definition of “gossip algorithm”, they show that the mixing time of the fastest Markov Chain on the graph essentially provides a lower bound on any gossip algorithm as well. Consequently, in [57], they study the mixing time of the fastest Markov Chain on random geometric graphs (which are good models for sensor networks), and show a lower bound of $\Theta(1/r^2)$ for nodes with communication range $r$ embedded in a unit square.

Since this bound is fairly slow, and inefficient in terms of requiring $\Omega(n^2)$ rounds of communication, Dimakis et al. [118] propose a protocol called geographic gossip. In that protocol, nodes choose a uniformly
random point to route to, and then use greedy multi-hop routing to exchange messages with the closest node to that point. Messages are rejected with certain probabilities to avoid oversampling nodes with few nearby other nodes. For this protocol, [118] proves a bound of $O((n \log n)^{1.5})$ on the number of transmissions required. Assuming parallelism and an appropriate communication radius $r$, this corresponds to $O(1/r^{1.5})$ rounds in the previous model.

Gossip-based algorithms can be used as a substrate for more complex computations beyond message dissemination and averaging. These include computations of spanning trees [207], various sketches of data stored at nodes [206] and eigenvectors and eigenvalues of a matrix held in a distributed way at the nodes [211]. In addition to Astrolabe [327, 328], several system approaches have been proposed using gossip (see, e.g., [195, 194]).

Finally, while we motivated the gossip problem with an eye on its fault-tolerance via randomization, it is an interesting question how quickly a message could be spread over a given arbitrary network, under the constraint that each current message holder can only forward the message to one neighbor in any given round. This Optimal Broadcast problem is NP-complete. Ravi [296] gave an $O(\log^2 n)$ approximation, which was subsequently improved by Bar-Noy et al. [26] to $O(\log n)$.
Bibliography


104


105


