Lecture 1: Introduction to Course & Introduction to Chernoff Bounds

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Outline

Introduction

Topics and Syllabus

A (Very) Brief Reminder of Probability Theory

Basic Examples

Introduction to Chernoff Bounds

What? Randomised Algorithms utilise random bits to compute their output.

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Why? Randomised Algorithms often provide an efficient (and elegant!) solution or approximation to a problem that is costly (or impossible) to solve deterministically.

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Why? Randomised Algorithms often provide an efficient (and elegant!) solution or approximation to a problem that is costly (or impossible) to solve deterministically.

But often: simple algorithm at the cost of a sophisticated analysis!

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Why? Randomised Algorithms often provide an efficient (and elegant!) solution or approximation to a problem that is costly (or impossible) to solve deterministically.

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- Donald E. Knuth (in Randomization and Religion)



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How? This course aims to strengthen your knowledge of probability theory and apply this to analyse examples of randomised algorithms.

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How? This course aims to strengthen your knowledge of probability theory and apply this to analyse examples of randomised algorithms.

What if I (initially) don't care about randomised algorithms?

Many of the techniques in this course (Markov Chains, Concentration of Measure, Spectral Theory) are very relevant to other popular areas of research and employment such as Data Science and Machine Learning.

Some stuff you should know...

In this course we will assume some basic knowledge of probability:

- random variable
- computing expectations and variances
- notions of independence and conditional probabilities
- "general" idea of how to compute probabilities (manipulating, counting and estimating)



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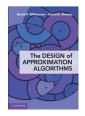
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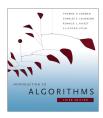


You should also be familiar with basic computer science, mathematics knowledge such as:

- graphs
- basic algorithms (sorting, graph algorithms etc.)
- matrices, norms and vectors







- (*) Michael Mitzenmacher and Eli Upfal. Probability and Computing: Randomized Algorithms and Probabilistic Analysis, Cambridge University Press, 2nd edition, 2017
- David P. Williamson and David B. Shmoys. The Design of Approximation Algorithms, Cambridge University Press, 2011
- Cormen, T.H., Leiserson, C.D., Rivest, R.L. and Stein, C. Introduction to Algorithms. MIT Press (3rd ed.), 2009 (We will adopt some of the labels (e.g., Theorem 35.6) from this book in Lectures 6-10)

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Lectures 2-5 focus on probabilistic tools and techniques.

2-3 Concentration (Lectures)

 Concept of Concentration; Recap of Markov and Chebyshev; Chernoff Bounds and Applications; Extensions: Hoeffding's Inequality and Method of Bounded Differences; Applications.

4 Markov Chains and Mixing Times (Lecture)

 Recap; Stopping and Hitting Times; Properties of Markov Chains; Convergence to Stationary Distribution; Variation Distance and Mixing Time

5 Hitting Times and Application to 2-SAT (Lecture)

 Reversible Markov Chains and Random Walks on Graphs; Cover Times and Hitting Times on Graphs (Example: Paths and Grids); A Randomised Algorithm for 2-SAT Algorithm

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Lectures 6-8 introduce linear programming, a (mostly) deterministic but very powerful technique to solve various optimisation problems.

6–7 Linear Programming (Lectures)

Introduction to Linear Programming, Applications, Standard and Slack Forms, Simplex Algorithm, Finding an Initial Solution, Fundamental Theorem of Linear Programming

8 Travelling Salesman Problem (Interactive Demo)

Hardness of the general TSP problem, Formulating TSP as an integer program; Classical TSP instance from 1954; Branch & Bound Technique to solve integer programs using linear programs

7

We then see how we can efficiently combine linear programming with randomised techniques, in particular, rounding:

9-10 Randomised Approximation Algorithms (Lectures)

 MAX-3-CNF and Guessing, Vertex-Cover and Deterministic Rounding of Linear Program, Set-Cover and Randomised Rounding, Concluding Example: MAX-CNF and Hybrid Algorithm We then see how we can efficiently combine linear programming with randomised techniques, in particular, rounding:

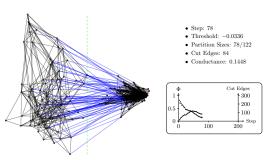
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Lectures 11-12 cover a more advanced topic with ML flavour:

11–12 Spectral Graph Theory and Spectral Clustering (Lectures)

 Eigenvalues, Eigenvectors and Spectrum; Visualising Graphs; Expansion; Cheeger's Inequality; Clustering and Examples; Analysing Mixing Times



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Recap: Probability Space

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Components of the Probability Space $(\Omega, \Sigma, \mathbf{P})$ -

- The Sample Space Ω contains all the possible, mutually exclusive outcomes $\omega_1, \omega_2, \ldots$ of the experiment.
- The Event Space Σ is the power-set of Ω containing events, which are combinations of outcomes (subsets of Ω including \emptyset and Ω).
- The Probability Measure **P** is a function from Σ to \mathbb{R} satisfying
 - $\begin{array}{ll} \text{(i)} & 0 \leq \textbf{P}\left[\,\mathcal{E}\,\right] \leq 1,\,\text{for all }\mathcal{E} \in \Sigma \\ \text{(ii)} & \textbf{P}\left[\,\Omega\,\right] = 1 \end{array}$

 - (iii) If $\mathcal{E}_1, \mathcal{E}_2, \ldots \in \Sigma$ are pairwise disjoint $(\mathcal{E}_i \cap \mathcal{E}_j = \emptyset)$ for all $i \neq j$) then

$$\mathbf{P}\left[\bigcup_{i=1}^{\infty} \mathcal{E}_i\right] = \sum_{i=1}^{\infty} \mathbf{P}\left[\mathcal{E}_i\right].$$

A random variable X on $(\Omega, \Sigma, \mathbf{P})$ is a function $X : \Omega \to \mathbb{R}$ mapping each sample "outcome" to a real number.

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Examples of random variables

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• The indicator random variable $\mathbf{1}_{\mathcal{E}}$ of an event $\mathcal{E} \in \Sigma$ given by

$$\mathbf{1}_{\mathcal{E}}(\omega) = egin{cases} 1 & \text{if } \omega \in \mathcal{E} \\ 0 & \text{otherwise}. \end{cases}$$

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■ The number of sixes of two dice throws $X_1, X_2 \in \{1, 2, ..., 6\}$ is

$$\mathbf{1}_{X_1=6} + \mathbf{1}_{X_2=6}$$

Union Bound -

Let $\mathcal{E}_1, \dots, \mathcal{E}_n$ be a collection of events in Σ . Then

$$\mathbf{P}\left[\bigcup_{i=1}^{n} \mathcal{E}_{i}\right] \leq \sum_{i=1}^{n} \mathbf{P}\left[\mathcal{E}_{i}\right].$$

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A Proof using Indicator Random Variables:

1. Let $\mathbf{1}_{\mathcal{E}_i}$ be the random variable that takes value 1 if \mathcal{E}_i holds, 0 otherwise

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- 3. It is clear that $\mathbf{1}_{\bigcup_{i=1}^{n} \mathcal{E}_{i}} \leq \sum_{i=1}^{n} \mathbf{1}_{\mathcal{E}_{i}}$ (Check this)
- 4. Taking expectation completes the proof.

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MAX-CUT Problem ——

■ Given: Undirected graph G = (V, E)

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MAX-CUT Problem -

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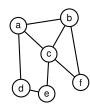
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14

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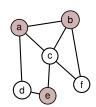
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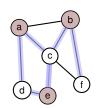
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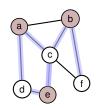
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Applications:



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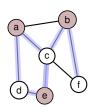
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- network or chip design
- machine learning
- statistical physics



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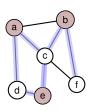
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Comments:

- MAX-CUT is NP-hard
- It is different from the clustering problem, where we want to find a sparse cut
- Note that the MIN-CUT problem is solvable in polynomial time!



$$S = \{a, b, e\}$$

 $e(S, S^c) = 6$

RANDMAXCUT(G)

This kind of "random guessing" will appear often in this course!

1: Start with $S \leftarrow \emptyset$

2: **For** each $v \in V$, add v to S with probability 1/2

3: Return S

RANDMAXCUT(G)

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Proposition ————

RANDMAXCUT(G) gives a 2-approximation using time O(n).

RANDMaxCut(G)

- 1: Start with $S \leftarrow \emptyset$
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Ratio between optimal and expected value of our solution is \leq 2 (more on this in Lecture 9)

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15

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Question:

- 1. What is the sample space Ω here?
- 2. Which quantity do we need to analyse?

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Proof:

• We need to analyse the expectation of $e(S, S^c)$:

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$$\mathbf{E}\left[\,e\left(\mathcal{S},\mathcal{S}^{c}\right)\,\right]$$

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RANDMAXCUT(G)

- 1: Start with $S \leftarrow \emptyset$
- 2: **For** each $v \in V$, add v to S with probability 1/2
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Ratio between optimal and expected value of our solution is \leq 2 (more on this in Lecture 9)

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Later: learn stronger tools that imply concentration around the expectation!

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Source: https://www.express.co.uk/life-style/life/567954/Discount-codes-money-saving-vouchers-coupons-mum

Coupon Collector Problem –

Suppose that there are n coupons to be collected from the cereal box. Every morning you open a new cereal box and get one coupon. We assume that each coupon appears with the same probability in the box.



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Exercise ([Ex. 1.11])



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16

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1. Prove it takes $n \sum_{k=1}^{n} \frac{1}{k} \approx n \log n$ expected boxes to collect all coupons



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- 1. Prove it takes $n \sum_{k=1}^{n} \frac{1}{k} \approx n \log n$ expected boxes to collect all coupons
- 2. Use Union Bound to prove that the probability it takes more than $n \log n + cn$ boxes to collect all n coupons is $< e^{-c}$.

Hint: It is useful to remember that $1 - x \le e^{-x}$ for all x

Outline

Introduction

Topics and Syllabus

A (Very) Brief Reminder of Probability Theory

Basic Examples

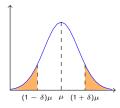
Introduction to Chernoff Bounds

Concentration Inequalities

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- It gives us the best of two worlds:
 - 1. Randomised Algorithms: Easy to Design and Implement
 - 2. Deterministic Algorithms: They do what they claim

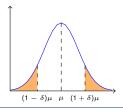


Chernoff Bounds: A Tool for Concentration (1952)

- Chernoffs bounds are "strong" bounds on the tail probabilities of sums of independent random variables
- random variables can be discrete (or continuous)
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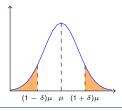


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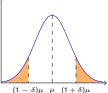
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 - Randomised Algorithms and Statstics
 - Random Projections and Dimensionality Reduction
 - Complexity Theory and Learning Theory (e.g., PAC-learning)



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Recap: Markov and Chebyshev

Markov's Inequality

If X is a non-negative random variable, then for any a > 0,

$$P[X \ge a] \le E[X]/a$$
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Chebyshev's inequality (or Markov) can be obtained by chosing $f(X) := (X - \mu)^2$ (or f(X) := X, respectively).

From Markov and Chebyshev to Chernoff

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We can consider the first, second, third and more moments! That is the basic idea behind the Chernoff Bounds

Our First Chernoff Bound

Chernoff Bounds (General Form, Upper Tail) =

Suppose X_1, \ldots, X_n are independent Bernoulli random variables with parameter p_i . Let $X = X_1 + \ldots + X_n$ and $\mu = \mathbf{E}[X] = \sum_{i=1}^n p_i$. Then, for any $\delta > 0$ it holds that

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By substitution, this implies that for any $t > \mu$,

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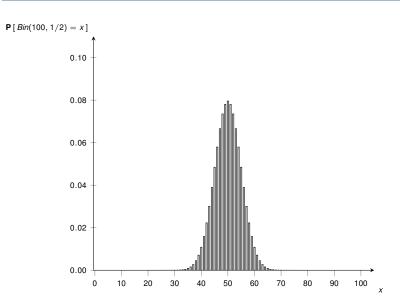
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What about a concrete value of
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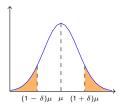
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Chernoff bound yields a much better result (but needs independence!)



Randomised Algorithms

Lecture 2: Concentration Inequalities, Application to Balls-into-Bins

Thomas Sauerwald (tms41@cam.ac.uk)





Outline

How to Derive Chernoff Bounds

Application 1: Balls into Bins

Appendix: More on Moment Generating Functions (non-examinable)

Recipe -

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- 3. Optimise value of λ to obtain best tail bound

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$$2. \ \mathbf{E}\left[\, \boldsymbol{e}^{\lambda X} \, \right] = \mathbf{E}\left[\, \boldsymbol{e}^{\lambda \, \sum_{i=1}^n \, X_i} \, \right] \underset{\text{indep}}{=} \, \prod_{i=1}^n \mathbf{E}\left[\, \boldsymbol{e}^{\lambda X_i} \, \right]$$

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Suppose X_1,\ldots,X_n are independent Bernoulli random variables with parameter p_i . Let $X=X_1+\ldots+X_n$ and $\mu=\mathbf{E}[X]=\sum_{i=1}^n p_i$. Then, for any $\delta>0$ it holds that

$$\mathbf{P}[X \geq (1+\delta)\mu] \leq \left[\frac{e^{\delta}}{(1+\delta)^{(1+\delta)}}\right]^{\mu}.$$

Proof:

1. For $\lambda > 0$,

$$\mathbf{P}\left[\,X \geq (1+\delta)\mu\,\right] \underset{e^{\lambda X} \text{ is inor}}{=} \mathbf{P}\left[\,e^{\lambda X} \geq e^{\lambda(1+\delta)\mu}\,\right] \underset{\mathsf{Markov}}{\leq} e^{-\lambda(1+\delta)\mu} \mathbf{E}\left[\,e^{\lambda X}\,\right]$$

$$2. \ \mathbf{E}\left[\, \boldsymbol{e}^{\lambda X} \, \right] = \mathbf{E}\left[\, \boldsymbol{e}^{\lambda \, \sum_{i=1}^n \, X_i} \, \right] \underset{\text{indep}}{=} \, \prod_{i=1}^n \mathbf{E}\left[\, \boldsymbol{e}^{\lambda X_i} \, \right]$$

3.

$$\mathbf{E}\left[e^{\lambda X_i}\right] = e^{\lambda} p_i + (1 - p_i) = 1 + p_i(e^{\lambda} - 1) \leq e^{\lambda p_i(e^{\lambda} - 1)}$$

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4. Putting all together

$$\mathbf{P}\left[X \geq (1+\delta)\mu\right] \leq e^{-\lambda(1+\delta)\mu} \prod_{i=1}^{n} e^{p_i(e^{\lambda}-1)} = e^{-\lambda(1+\delta)\mu} e^{\mu(e^{\lambda}-1)}$$

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5. Choose $\lambda = \log(1 + \delta) > 0$ to get the result.

Chernoff Bounds: Lower Tails

We can also use Chernoff Bounds to show a random variable is **not too** small compared to its mean:

Chernoff Bounds (General Form, Lower Tail) —

Suppose X_1,\ldots,X_n are independent Bernoulli random variables with parameter p_i . Let $X=X_1+\ldots+X_n$ and $\mu=\mathbf{E}[X]=\sum_{i=1}^n p_i$. Then, for any $0<\delta<1$ it holds that

$$\mathbf{P}[X \leq (1-\delta)\mu] \leq \left[\frac{e^{-\delta}}{(1-\delta)^{1-\delta}}\right]^{\mu},$$

and thus, by substitution, for any $t < \mu$,

$$\mathbf{P}[X \leq t] \leq e^{-\mu} \left(\frac{e\mu}{t}\right)^t.$$

Exercise on Supervision Sheet

Hint: multiply both sides by -1 and repeat the proof of the Chernoff Bound

Nicer Chernoff Bounds

Nicer Chernoff Bounds



Suppose $X_1, ..., X_n$ are independent Bernoulli random variables with parameter p_i . Let $X = X_1 + ... + X_n$ and $\mu = \mathbf{E}[X] = \sum_{i=1}^n p_i$. Then,

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■ For all t > 0,

$$P[X \ge E[X] + t] \le e^{-2t^2/n}$$

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$$\mathbf{P}[X \ge (1+\delta)\mathbf{E}[X]] \le \exp\left(-\frac{\delta^2\mathbf{E}[X]}{3}\right)$$

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All upper tail bounds hold even under a relaxed independence assumption: For all $1 \le i \le n$ and $x_1, x_2, \dots, x_{i-1} \in \{0, 1\}$,

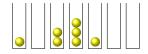
$$P[X_i = 1 \mid X_1 = x_1, \dots, X_{i-1} = x_{i-1}] \le p_i.$$

Outline

How to Derive Chernoff Bounds

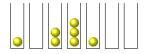
Application 1: Balls into Bins

Appendix: More on Moment Generating Functions (non-examinable)



Balls into Bins Model -

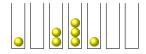
You have m balls and n bins. Each ball is allocated in a bin picked independently and uniformly at random.



Balls into Bins Model -

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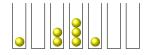
A very natural but also rich mathematical model



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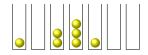
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Exercise: Think about the relation between the Balls into Bins Model and the Coupon Collector Problem.



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$$\boxed{ \mathbf{P}[X \geq t] \leq e^{-\mu} (e\mu/t)^t }$$

By the Chernoff Bound,

$$P[X \ge 6 \log n] \le e^{-2 \log n} \left(\frac{2e \log n}{6 \log n}\right)^{6 \log n} \le e^{-2 \log n} = n^{-2}$$



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An event \mathcal{E} (that implicitly depends on an input parameter n) occurs whp if $\mathbf{P}[\mathcal{E}] \to 1$ as $n \to \infty$.

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- By pigeonhole principle, the max loaded bin receives at least 2 log n balls.
 Hence our bound is pretty sharp.

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• One can prove that whp at least one bin receives at least $c \log n / \log \log n$ balls, for some constant c > 0.

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For any $m \ge n$, we can improve this by sampling two bins in each step and then assign the ball into the bin with lesser load.

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Conclusions

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 \Rightarrow for m = n this gives a maximum load of $\log_2 \log n + \Theta(1)$ w.p. 1 - 1/n.

This is called the **power of two choices**: It is a common technique to improve the performance of randomised algorithms (covered in Chapter 17 of the textbook by Mitzenmacher and Upfal)

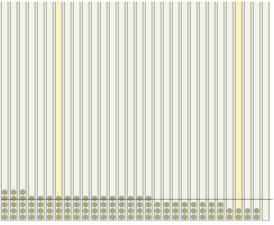
ACM Paris Kanellakis Theory and Practice Award 2020



For "the discovery and analysis of balanced allocations, known as the power of two choices, and their extensive applications to practice."

"These include i-Google's web index, Akamai's overlay routing network, and highly reliable distributed data storage systems used by Microsoft and Dropbox, which are all based on variants of the power of two choices paradigm. There are many other software systems that use balanced allocations as an important ingredient."

Simulation



Sampled two bins u.a.r.

https://www.dimitrioslos.com/balls_and_bins/visualiser.html

Outline

How to Derive Chernoff Bounds

Application 1: Balls into Bins

Appendix: More on Moment Generating Functions (non-examinable)

Moment-Generating Function ——

The moment-generating function of a random variable *X* is

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Lemma

- 1. If X and Y are two r.v.'s with $M_X(t) = M_Y(t)$ for all $t \in (-\delta, +\delta)$ for some $\delta > 0$, then the distributions X and Y are identical.
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- 1. If X and Y are two r.v.'s with $M_X(t) = M_Y(t)$ for all $t \in (-\delta, +\delta)$ for some $\delta > 0$, then the distributions X and Y are identical.
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$$M_{X+Y}(t) = M_X(t) \cdot M_Y(t).$$

Proof of 2:

$$M_{X+Y}(t) = \mathbf{E} \left[e^{t(X+Y)} \right] = \mathbf{E} \left[e^{tX} \cdot e^{tY} \right] \stackrel{(!)}{=} \mathbf{E} \left[e^{tX} \right] \cdot \mathbf{E} \left[e^{tY} \right] = M_X(t) M_Y(t) \quad \Box$$

Randomised Algorithms

Lecture 3: Concentration Inequalities, Application to Quick-Sort, Extensions

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Application 2: Randomised QuickSort

Extensions of Chernoff Bounds

Applications of Method of Bounded Differences

Appendix: More on Moment Generating Functions (non-examinable)

QuickSort

```
QUICKSORT (Input A[1], A[2], \ldots, A[n])
 1: Pick an element from the array, the so-called pivot
2: If n = 0 or n = 1 then
3.
        return A
4: else
5:
        Create two subarrays A_1 and A_2 (without the pivot) such that:
           A_1 contains the elements that are smaller than the pivot
6.
            A<sub>2</sub> contains the elements that are greater (or equal) than the pivot
7:
        QUICKSORT(A_1)
8.
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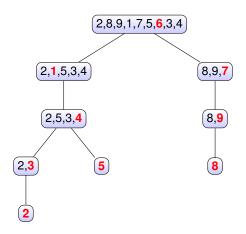
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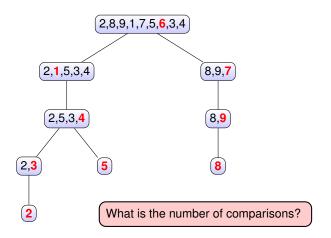
- Example: Let A = (2, 8, 9, 1, 7, 5, 6, 3, 4) with A[7] = 6 as pivot. ⇒ $A_1 = (2, 1, 5, 3, 4)$ and $A_2 = (8, 9, 7)$
- Worst-Case Complexity (number of comparisons) is Θ(n²), while Average-Case Complexity is O(n log n).

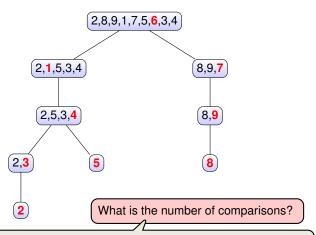
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- Example: Let A = (2, 8, 9, 1, 7, 5, 6, 3, 4) with A[7] = 6 as pivot. $\Rightarrow A_1 = (2, 1, 5, 3, 4)$ and $A_2 = (8, 9, 7)$
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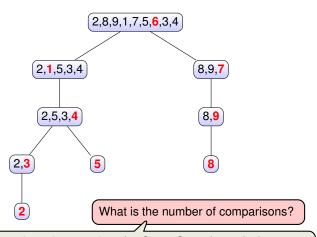
We will now give a proof of this "well-known" result!







Note that the number of comparison by QUICKSORT is equivalent to the sum of the depths of all nodes in the tree (why?).



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$$0+1+1+2+2+3+3+3+4=19.$$

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This should be your standard answer in this course ©

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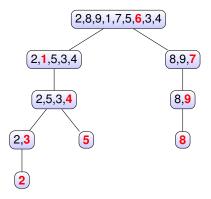
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4. Actually, we will prove sth slightly stronger:

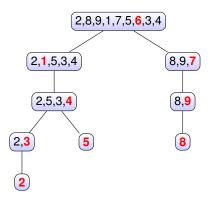
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Let P be a path from the root to the deepest level of some element

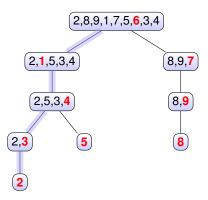
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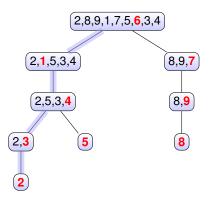
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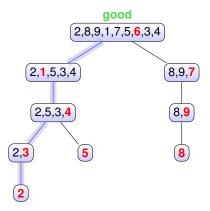
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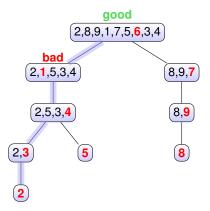
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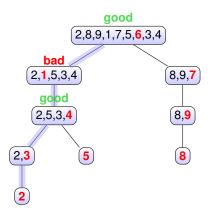
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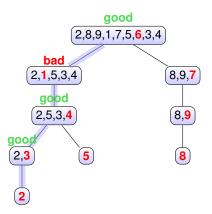
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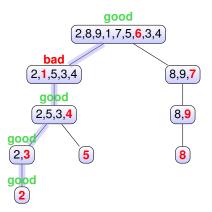
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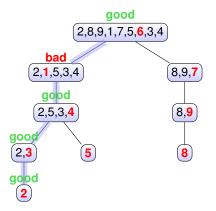
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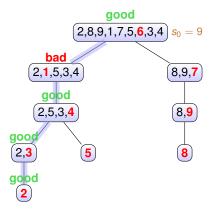
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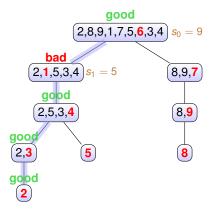
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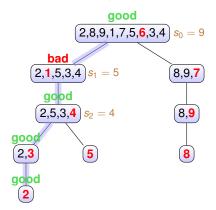
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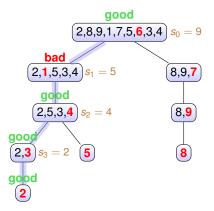
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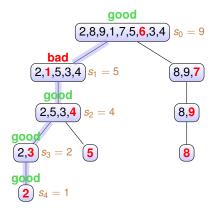
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• We start with $s_0 = n$

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This even holds always,

i.e., deterministically!

⇒ There are at most $T = \frac{\log n}{\log(3/2)} < 5 \log n$ many good nodes on any path P.

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 - \Rightarrow number of **bad** nodes in the first 24 log *n* levels is more than 21 log *n*.

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 - Assume $|P| \ge C \log n$ for C := 24
 - \Rightarrow number of **bad** nodes in the first 24 log *n* levels is more than 21 log *n*.

Let us now upper bound the probability that this "bad event" happens!

• Consider the first 24 log *n* nodes of *P* to the deepest level of element *i*.

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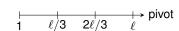
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Question: Edge Case: What if the path P does not reach level j?

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Question: Edge Case: What if the path P does not reach level j?

Answer: We can then simply define X_i as 0 (deterministically).

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- Well-known: any comparison-based sorting algorithm needs $\Omega(n \log n)$
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Exercise: [Ex 2-3.6] Our upper bound of $O(n \log n)$ whp also immediately implies a $O(n \log n)$ bound on the expected number of comparisons!

- It is possible to deterministically find the best pivot element that divides the array into two subarrays of the same size.
- The latter requires to compute the median of the array in linear time, which is not easy...
- The presented randomised algorithm for QUICKSORT is much easier to implement!

Outline

Application 2: Randomised QuickSort

Extensions of Chernoff Bounds

Applications of Method of Bounded Differences

Appendix: More on Moment Generating Functions (non-examinable)

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Hoeffding's Extension Lemma ———

Let X be a random variable with mean 0 such that $a \leq X \leq b$. Then for all $\lambda \in \mathbb{R}$,

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We omit the proof of this lemma!

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Hoeffding's Inequality

Let X_1,\ldots,X_n be independent random variables with mean μ_i such that $a_i \leq X_i \leq b_i$. Let $X = X_1 + \ldots + X_n$, and let $\mu = \mathbf{E}[X] = \sum_{i=1}^n \mu_i$. Then for any t > 0,

$$\mathbf{P}\left[X \geq \mu + t\right] \leq \exp\left(-\frac{2t^2}{\sum_{i=1}^{n}(b_i - a_i)^2}\right),\,$$

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Proof Outline (skipped):

■ Let
$$X_i' = X_i - \mu_i$$
 and $X' = X_1' + ... + X_n'$, then $P[X \ge \mu + t] = P[X' \ge t]$

•
$$\mathbf{P}[X' \ge t] \le e^{-\lambda t} \prod_{i=1}^n \mathbf{E}\left[e^{\lambda X_i'}\right] \le \exp\left[-\lambda t + \frac{\lambda^2}{8} \sum_{i=1}^n (b_i - a_i)^2\right]$$

■ Choose $\lambda = \frac{4t}{\sum_{i=1}^{n} (b_i - a_i)^2}$ to get the result.

This is not "magic" – you just need to optimise λ !

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In all those cases (and more) we can easily prove concentration of $f(X_1, ..., X_n)$ around its mean by the so-called **Method of Bounded Differences**.

A function f is called Lipschitz with parameters $\mathbf{c} = (c_1, \dots, c_n)$ if for all $i = 1, 2, \dots, n$,

$$|f(x_1, x_2, \ldots, x_{i-1}, \mathbf{x}_i, x_{i+1}, \ldots, x_n) - f(x_1, x_2, \ldots, x_{i-1}, \mathbf{x}_i, x_{i+1}, \ldots, x_n)| \leq c_i,$$

where x_i and \tilde{x}_i are in the domain of the *i*-th coordinate.

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- The proof is omitted here (it requires the concept of martingales).

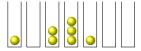
Outline

Application 2: Randomised QuickSort

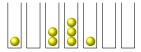
Extensions of Chernoff Bounds

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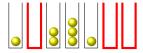
Appendix: More on Moment Generating Functions (non-examinable)



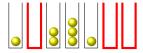
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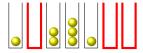


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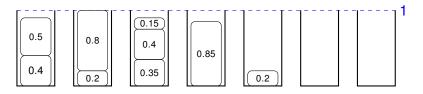
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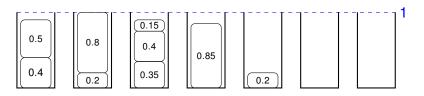
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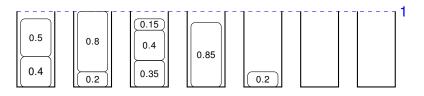
This is a decent bound, but for some values of m it is far from tight and stronger bounds are possible through a refined analysis.



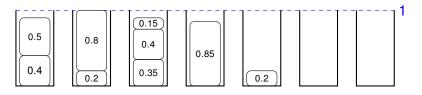
• We are given *n* items of sizes in the unit interval [0, 1]



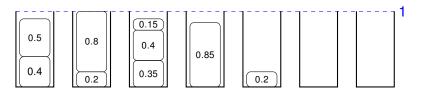
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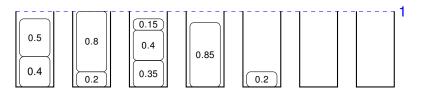
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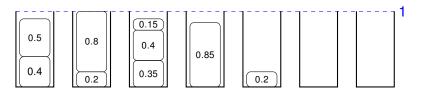
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Application 4: Bin Packing



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This is a typical example where proving concentration is much easier than calculating (or estimating) the expectation!

Outline

Application 2: Randomised QuickSort

Extensions of Chernoff Bounds

Applications of Method of Bounded Differences

Appendix: More on Moment Generating Functions (non-examinable)

Moment-Generating Function ——

The moment-generating function of a random variable *X* is

$$M_X(t) = \mathbf{E}\left[e^{tX}\right], \quad \text{where } t \in \mathbb{R}.$$

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Lemma

- 1. If X and Y are two r.v.'s with $M_X(t) = M_Y(t)$ for all $t \in (-\delta, +\delta)$ for some $\delta > 0$, then the distributions X and Y are identical.
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$$M_{X+Y}(t) = M_X(t) \cdot M_Y(t).$$

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Proof of 2:

$$M_{X+Y}(t) = \mathbf{E} \left[e^{t(X+Y)} \right] = \mathbf{E} \left[e^{tX} \cdot e^{tY} \right] \stackrel{(!)}{=} \mathbf{E} \left[e^{tX} \right] \cdot \mathbf{E} \left[e^{tY} \right] = M_X(t) M_Y(t) \quad \Box$$

Randomised Algorithms

Lecture 4: Markov Chains and Mixing Times

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Recap of Markov Chain Basics

Irreducibility, Periodicity and Convergence

Total Variation Distance and Mixing Times

Application 1: Card Shuffling

Application 2: Markov Chain Monte Carlo (non-examin.)

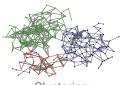
Appendix: Remarks on Mixing Time (non-examin.)



Broadcasting



Broadcasting



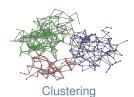
Clustering



Broadcasting



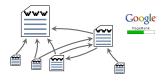
Ranking Websites



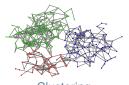
4. Markov Chains and Mixing Times © T. Sauerwald



Broadcasting



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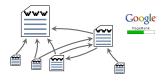


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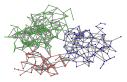




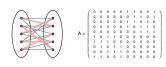
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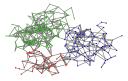




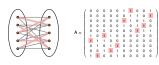
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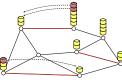
Sampling and Optimisation



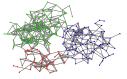
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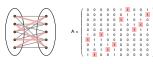
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Load Balancing



Clustering



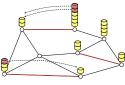
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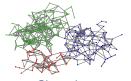
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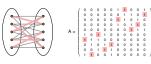
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Particle Processes

Markov Chain (Discrete Time and State, Time Homogeneous)

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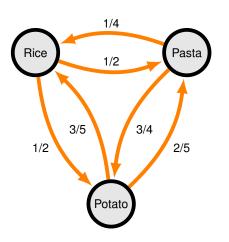
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• For all $0 \le t_1 < t_2, x \in \Omega$,

$$\mathbf{P}[X_{t_2} = x] = \sum_{y \in \Omega} \mathbf{P}[X_{t_2} = x \mid X_{t_1} = y] \cdot \mathbf{P}[X_{t_1} = y].$$

What does a Markov Chain Look Like?

Example: the carbohydrate served with lunch in the college cafeteria.



This has transition matrix:

$$P = \begin{bmatrix} \text{Rice} & \text{Pasta} & \text{Potato} \\ 0 & 1/2 & 1/2 \\ 1/4 & 0 & 3/4 \\ 3/5 & 2/5 & 0 \end{bmatrix} \begin{array}{c} \text{Rice} \\ \text{Pasta} \\ \text{Potato} \\ \text{Potato} \\ \end{array}$$



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• The Markov Property and line above imply that for any $t \ge 0$

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- \Rightarrow can replace ρ by any (load) vector and view P as a balancing matrix!

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Some distinguish between $\tau_y^+ = \min\{t \ge 1 \colon X_t = y\}$ and $\tau_y = \min\{t \ge 0 \colon X_t = y\}$

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A Useful Identity ———

Hitting times are the solution to a set of linear equations:

$$h(x,y) \stackrel{\mathsf{Markov} \ \mathsf{Prop.}}{=} 1 + \sum_{z \in \Omega \setminus \{y\}} P(x,z) \cdot h(z,y) \qquad \forall x \neq y \in \Omega.$$

Outline

Recap of Markov Chain Basics

Irreducibility, Periodicity and Convergence

Total Variation Distance and Mixing Times

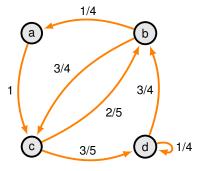
Application 1: Card Shuffling

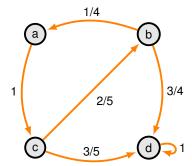
Application 2: Markov Chain Monte Carlo (non-examin.)

Appendix: Remarks on Mixing Time (non-examin.)

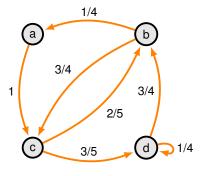
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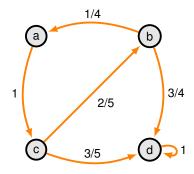
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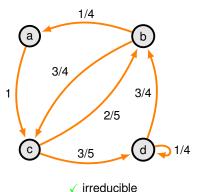


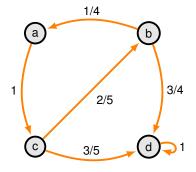




Question: Which of the two chains (if any) are irreducible?

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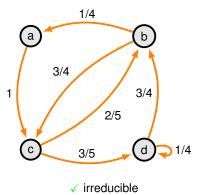


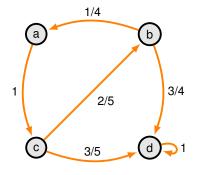
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Finite Hitting Time Theorem —

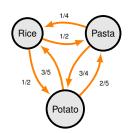
For any states x and y of a finite irreducible Markov Chain $h(x, y) < \infty$.

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$$\left(\frac{4}{13}, \frac{4}{13}, \frac{5}{13}\right) \cdot \begin{pmatrix} 0 & 1/2 & 1/2 \\ 1/4 & 0 & 3/4 \\ 3/5 & 2/5 & 0 \end{pmatrix} = \left(\frac{4}{13}, \frac{4}{13}, \frac{5}{13}\right)$$



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Rice
$$\frac{1}{1/2}$$

$$\frac{3}{13}$$
Potato

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Pasta

2/5

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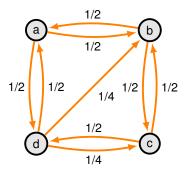
Existence and Uniqueness of a Positive Stationary Distribution -

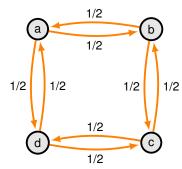
Let *P* be finite, irreducible M.C., then there exists a unique probability distribution π on Ω such that $\pi = \pi P$ and $\pi(x) = 1/h(x, x) > 0$, $\forall x \in \Omega$.

• A Markov Chain is aperiodic if for all $x \in \Omega$, $gcd\{t \ge 1 : P^t(x,x) > 0\} = 1$.

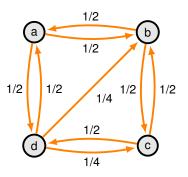
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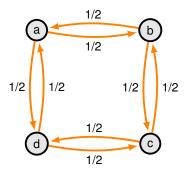
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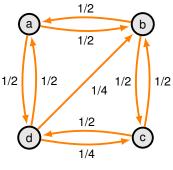






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1/2

1/2

1/2

1/2

1/2

1/2

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1/2

1/2

Convergence Theorem —

Let P be any finite, irreducible, aperiodic Markov Chain with stationary distribution π . Then for any $x,y\in\Omega$,

$$\lim_{t\to\infty}P^t(x,y)=\pi(y).$$

 ${\sf Ergodic} = {\sf Irreducible} + {\sf Aperiodic}$

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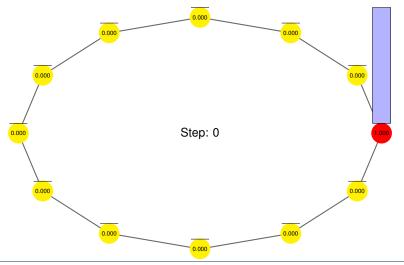
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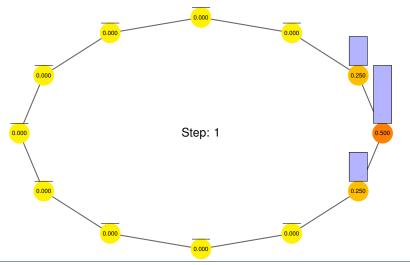
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 We will prove a simpler version of the Convergence Theorem after introducing Spectral Graph Theory.

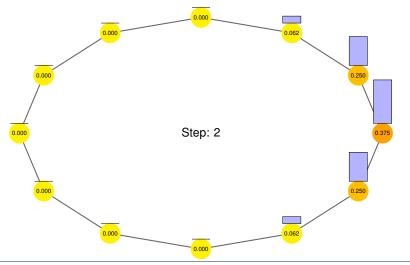
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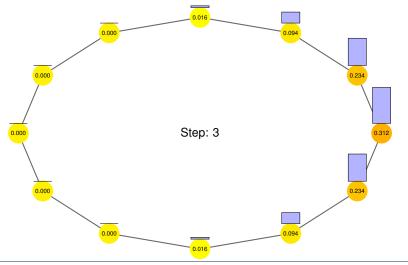
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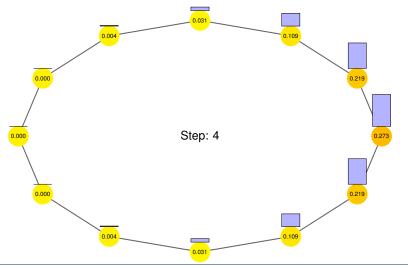
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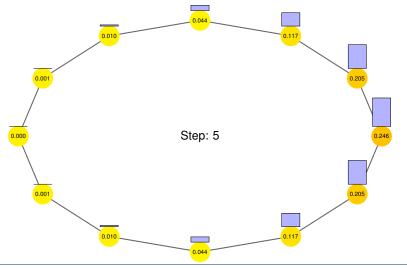
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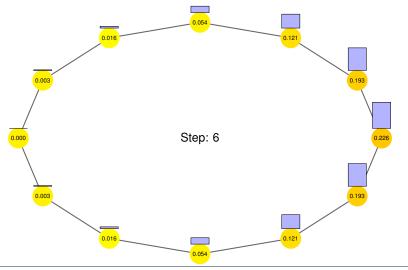
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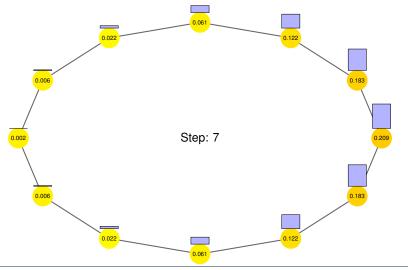
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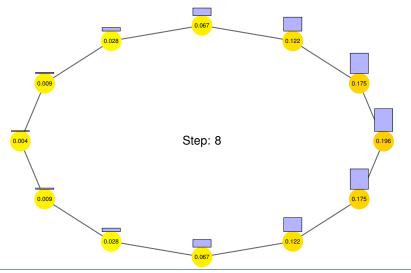
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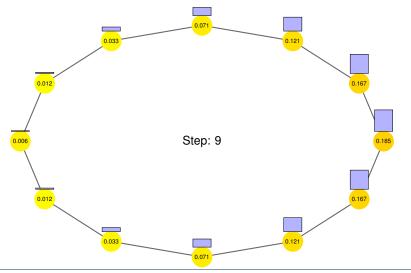
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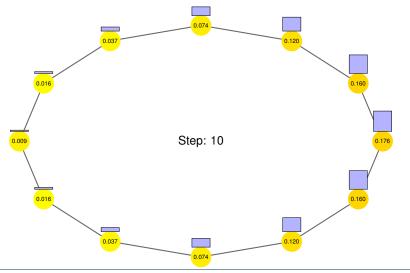
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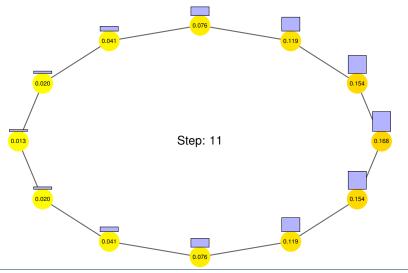
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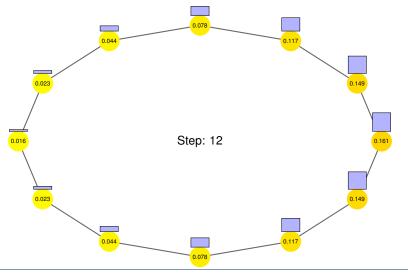
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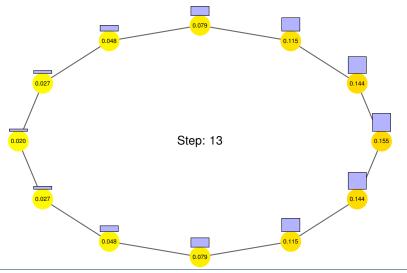
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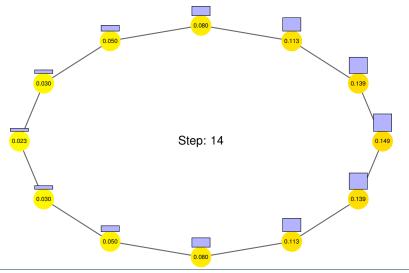
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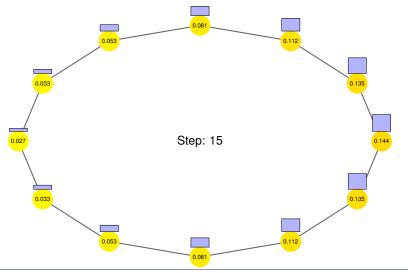
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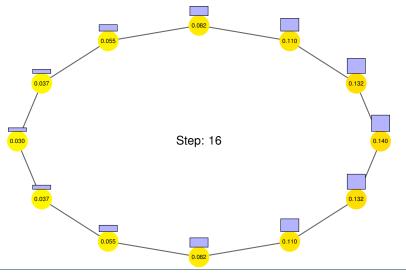
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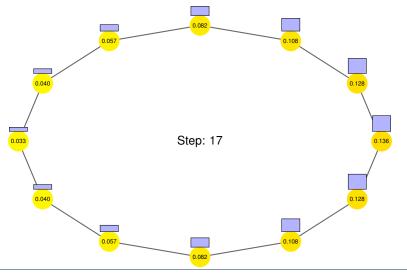
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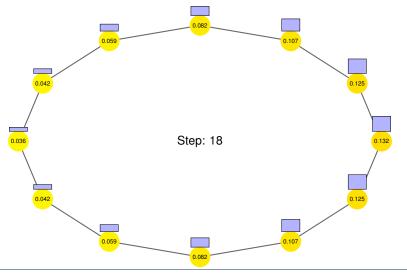
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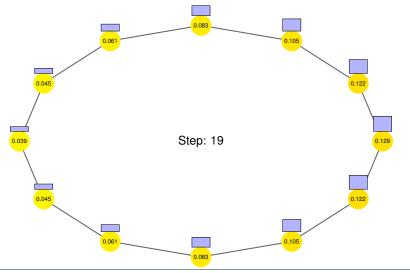
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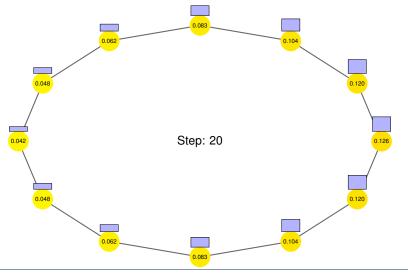
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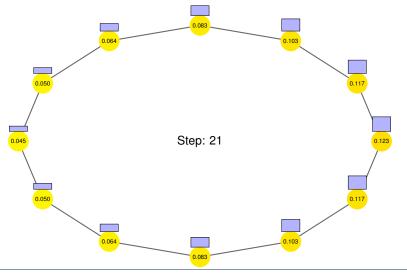
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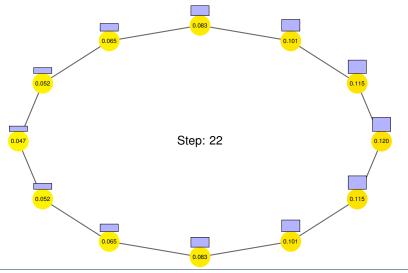
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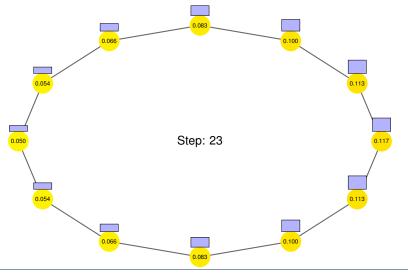
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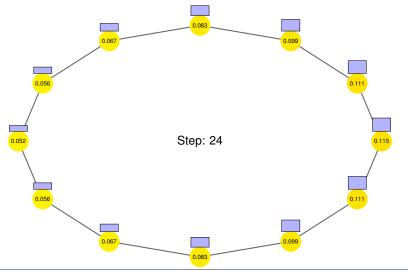
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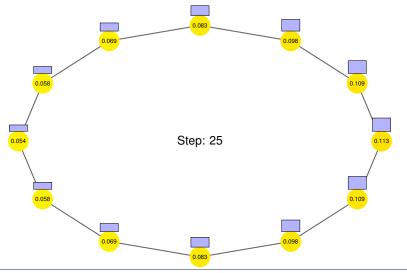
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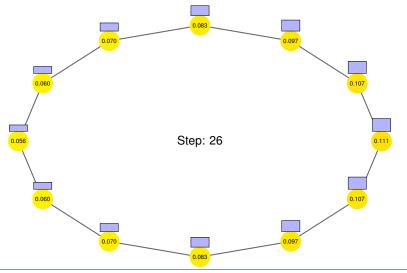
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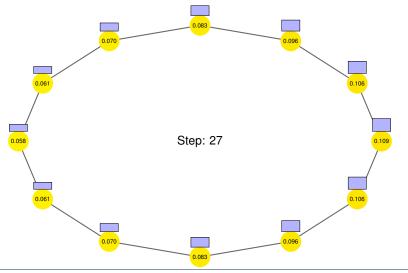
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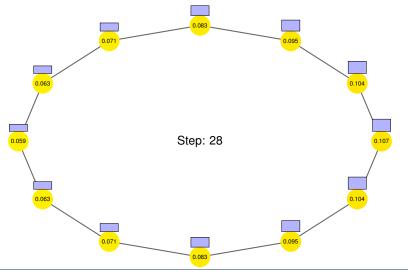
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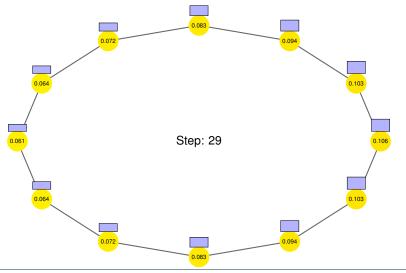
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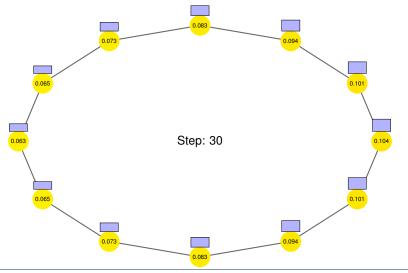
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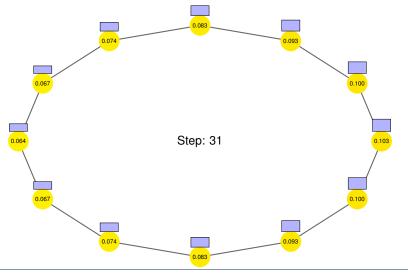
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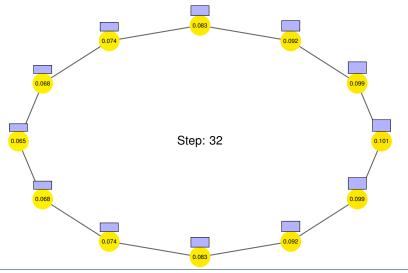
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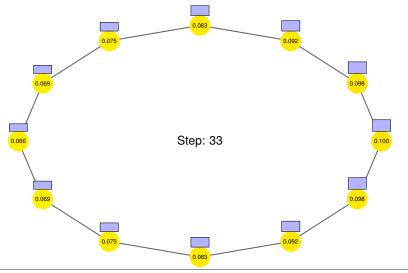
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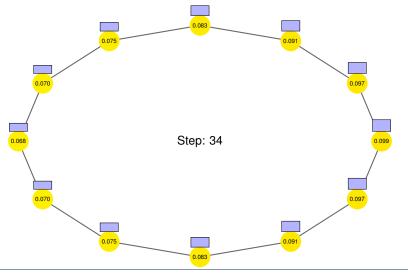
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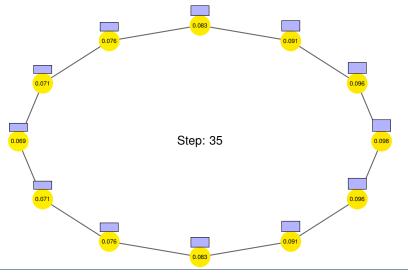
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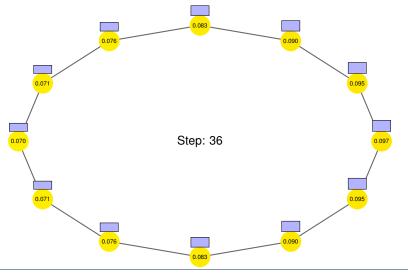
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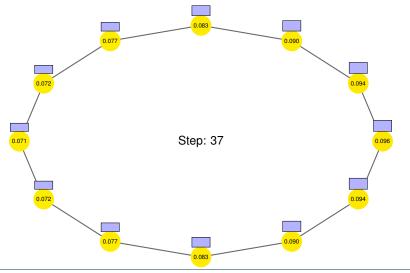
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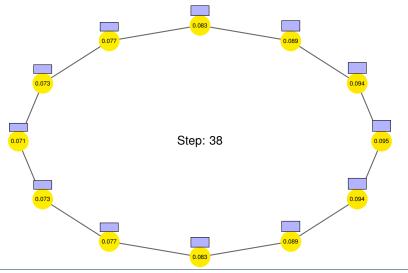
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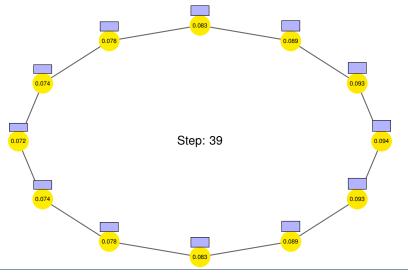
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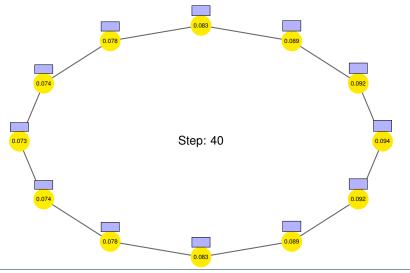
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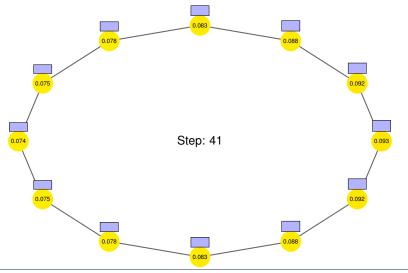
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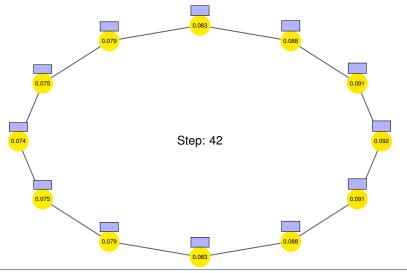
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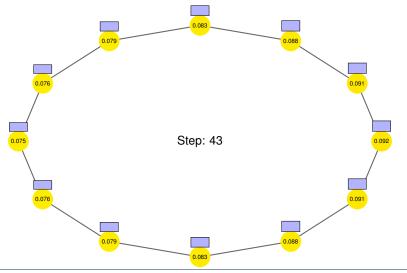
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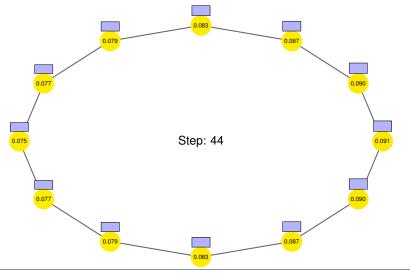
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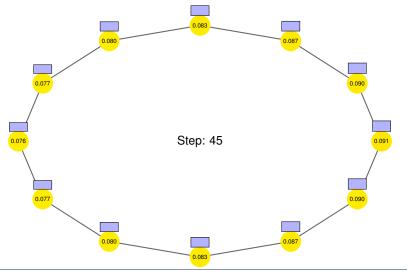
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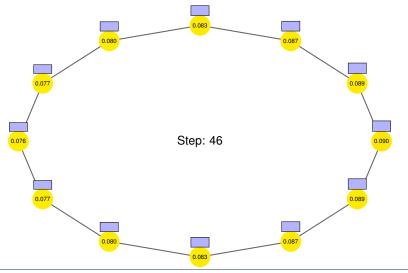
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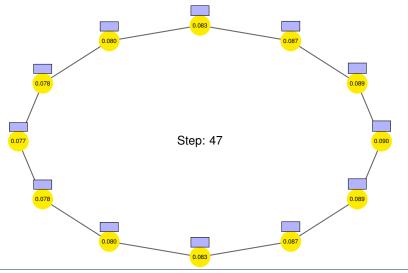
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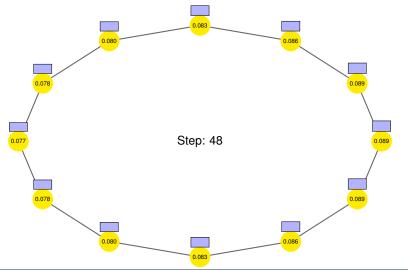
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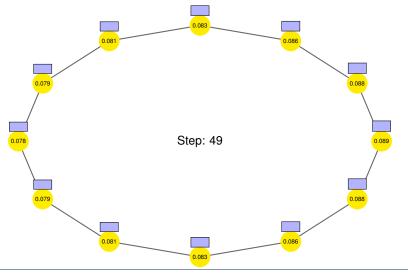


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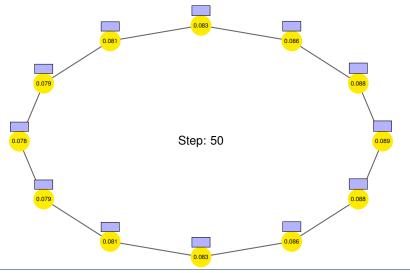
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Outline

Recap of Markov Chain Basics

Irreducibility, Periodicity and Convergence

Total Variation Distance and Mixing Times

Application 1: Card Shuffling

Application 2: Markov Chain Monte Carlo (non-examin.)

Appendix: Remarks on Mixing Time (non-examin.)

Loaded Dice -

• You are presented three loaded (unfair) dice A, B, C:

X	1	2	3	4	5	6
P[A=x]	1/3	1/12	1/12	1/12	1/12	1/3
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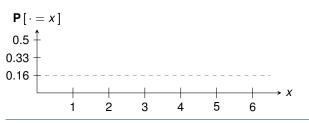
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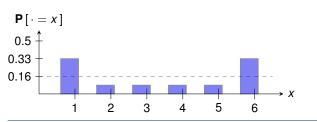
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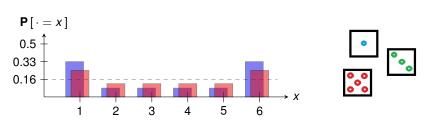
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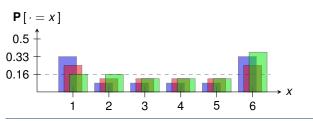
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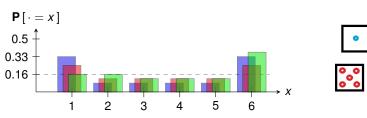
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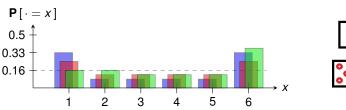
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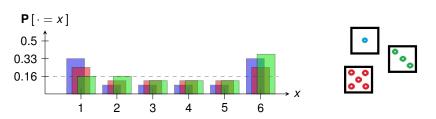
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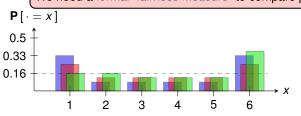
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We need a formal "fairness measure" to compare probability distributions!





The Total Variation Distance between two probability distributions μ and η on a countable state space Ω is given by

$$\|\mu - \eta\|_{tv} = \frac{1}{2} \sum_{\omega \in \Omega} |\mu(\omega) - \eta(\omega)|.$$

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Thus

$$\|D - B\|_{tv} = \|D - C\|_{tv}$$
 and $\|D - B\|_{tv}, \|D - C\|_{tv} < \|D - A\|_{tv}.$

So A is the least "fair", however B and C are equally "fair" (in TV distance).

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We will see a similar result later after introducing spectral techniques (Lecture 12)!

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See final slides for some comments on why we choose 1/4.

Outline

Recap of Markov Chain Basics

Irreducibility, Periodicity and Convergence

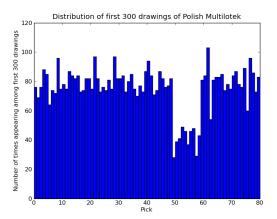
Total Variation Distance and Mixing Times

Application 1: Card Shuffling

Application 2: Markov Chain Monte Carlo (non-examin.)

Appendix: Remarks on Mixing Time (non-examin.)

Experiment Gone Wrong...



Thanks to Krzysztof Onak (pointer) and Eric Price (graph)

Source: Slides by Ronitt Rubinfeld



Source: wikipedia

How long does it take to shuffle a deck of 52 cards?



Source: wikipedia

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Persi Diaconis (Professor of Statistics and former Magician)

Source: www.soundcloud.com



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One of the leading experts in the field who has related card shuffling to many other mathematical problems.

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Here we will focus on one shuffling scheme which is easy to analyse.

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What is Card Shuffling?



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The Card Shuffling Markov Chain

TOPTORANDOMSHUFFLE (Input: A pile of *n* cards)

- 1: **For** t = 1, 2, ...
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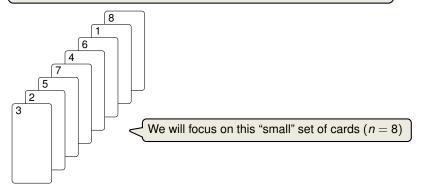
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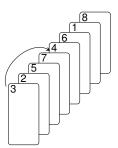
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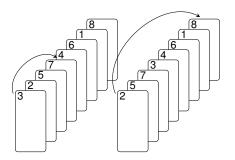
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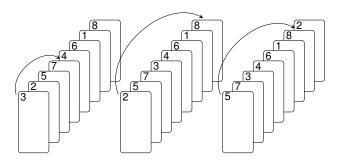
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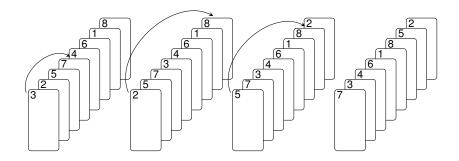
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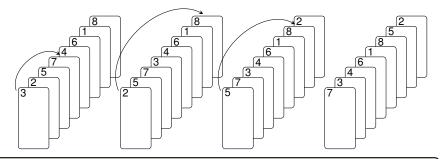


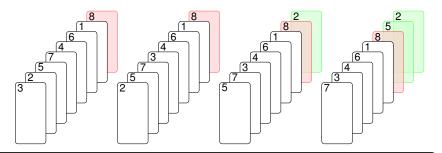


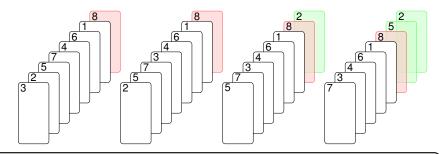


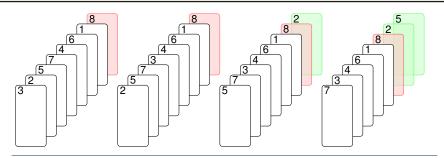


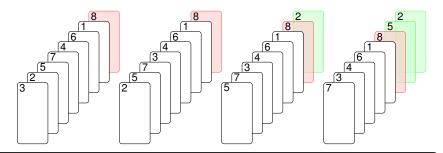


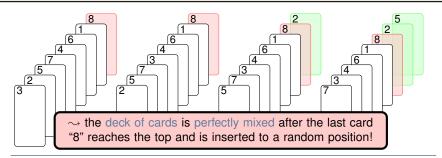


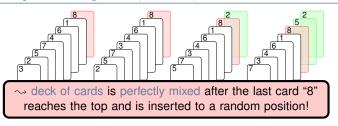


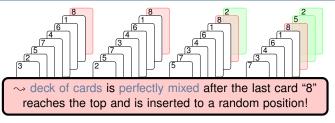




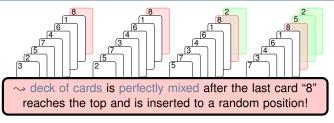




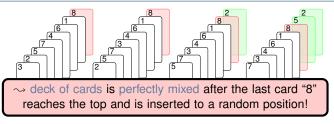




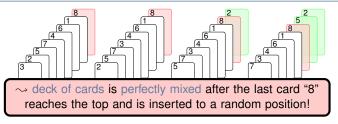
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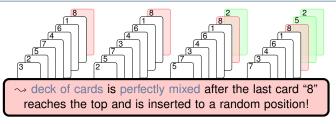
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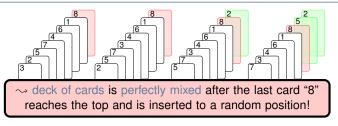
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- At the last position, card "n" moves up with probability $\frac{1}{n}$ at each step
- At the second last position, card "n" moves up with probability $\frac{2}{n}$



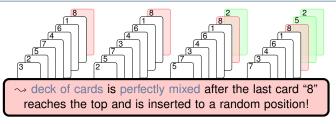
- How long does it take for the last card "n" to become top card?
- At the last position, card "n" moves up with probability $\frac{1}{n}$ at each step
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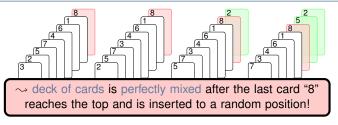
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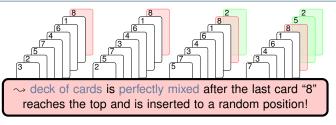


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Using the so-called coupling method, one could prove $t_{mix} \leq n \log n$.

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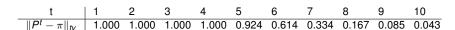


Figure: Total Variation Distance for *t* riffle shuffles of 52 cards.

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The Annals of Applied Probability 1992, Vol. 2, No. 2, 294-313

TRAILING THE DOVETAIL SHUFFLE TO ITS LAIR

By Dave Bayer 1 and Persi Diaconis 2

Columbia University and Harvard University

We analyze the most commonly used method for shuffling cards. The main result is a simple expression for the chance of any arrangement after any number of shuffles. This is used to give sharp bounds on the approach to randomness: $\frac{3}{2} \log_2 n + \theta$ shuffles are necessary and sufficient to mix up n cards.

Key ingredients are the analysis of a card trick and the determination of the idempotents of a natural commutative subalgebra in the symmetric group algebra.

t	1	2	3	4	5	6	7	8	9	10
$\ P^t - \pi\ _{tv}$	1.000	1.000	1.000	1.000	0.924	0.614	0.334	0.167	0.085	0.043

Figure: Total Variation Distance for *t* riffle shuffles of 52 cards.

Outline

Recap of Markov Chain Basics

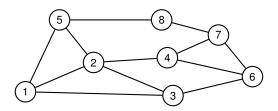
Irreducibility, Periodicity and Convergence

Total Variation Distance and Mixing Times

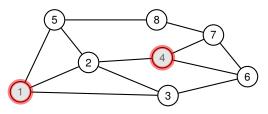
Application 1: Card Shuffling

Application 2: Markov Chain Monte Carlo (non-examin.)

Appendix: Remarks on Mixing Time (non-examin.)

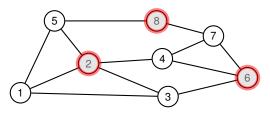


Independent Set —



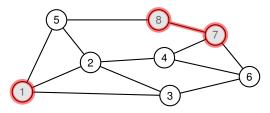
 $S = \{1, 4\}$ is an independent set \checkmark

Independent Set -



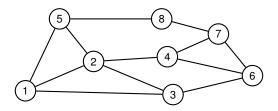
 $S = \{2, 6, 8\}$ is an independent set $\sqrt{}$

Independent Set -

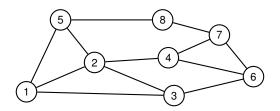


 $S = \{1, 7, 8\}$ is **not** an independent set \times

Independent Set -



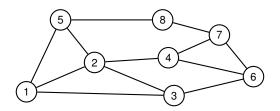
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Given an undirected graph G = (V, E), an independent set is a subset $S \subseteq V$ such that there are no two vertices $u, v \in S$ with $\{u, v\} \in E(G)$.

How can we take a sample from the space of all independent sets?

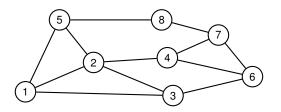


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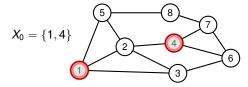
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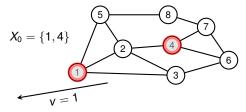
We can use a generic Markov Chain Monte Carlo approach to tackle this problem!

- 1: Let X_0 be an arbitrary independent set in G
- 2: **For** t = 0, 1, 2, ...
- 3: Pick a vertex $v \in V(G)$ uniformly at random
- 4: If $v \in X_t$ then $X_{t+1} \leftarrow X_t \setminus \{v\}$
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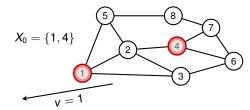
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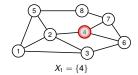


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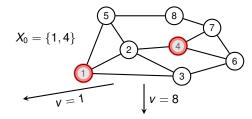


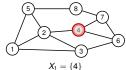
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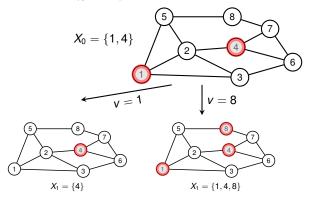


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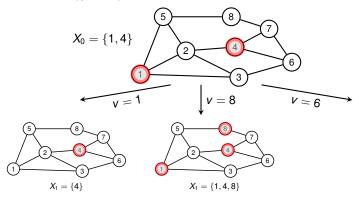




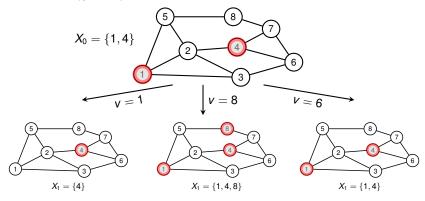
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not covered here, see the textbook by Mitzenmacher and Upfal

Outline

Recap of Markov Chain Basics

Irreducibility, Periodicity and Convergence

Total Variation Distance and Mixing Times

Application 1: Card Shuffling

Application 2: Markov Chain Monte Carlo (non-examin.)

Appendix: Remarks on Mixing Time (non-examin.)

• One can prove $\max_X \|P_X^t - \pi\|_{tv}$ is non-increasing in t (this means if the chain is " ϵ -mixed" at step t, then this also holds in future steps) [Mitzenmacher, Upfal, 12.3]

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Remark: This freedom on how to pick ϵ relies on the sub-multiplicative property of a (version) of the variation distance. First, let

$$d(t) := \max_{x} \left\| P_{x}^{t} - \pi \right\|_{t_{t}}$$

be the variation distance after t steps when starting from the worst state. Further, define

$$\overline{d}(t) := \max_{\mu,\nu} \left\| P_{\mu}^{t} - P_{\nu}^{t} \right\|_{t\nu}.$$

These quantities are related by the following double inequality

$$d(t) < \overline{d}(t) < 2d(t)$$
.

Further, $\overline{d}(t)$ is sub-multiplicative, that is for any s, t > 1,

$$\overline{d}(s+t) < \overline{d}(s) \cdot \overline{d}(t)$$
.

Hence for any fixed 0 $< \epsilon < \delta < 1/2$ it follows from the above that

$$\tau(\epsilon) \leq \left\lceil \frac{\ln \epsilon}{\ln(2\delta)} \right\rceil \tau(\delta).$$

In particular, for any $\epsilon < 1/4$

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This 2 is the reason why we ultimately need $\epsilon < 1/2$ in this derivation. On the other hand, see [Exercise (4/5).8] why $\epsilon < 1/2$ is also necessary.

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In particular, for any $\epsilon < 1/4$

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Hence smaller constants $\epsilon < 1/4$ only increase the mixing time by some constant factor.

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Randomised Algorithms

Lecture 5: Random Walks, Hitting Times and Application to 2-SAT

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Application 3: Ehrenfest Chain and Hypercubes

Random Walks on Graphs, Hitting Times and Cover Times

Random Walks on Paths and Grids

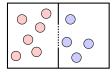
SAT and a Randomised Algorithm for 2-SAT

Ehrenfest Model ———

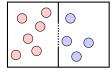
 A simple model for the exchange of molecules between two boxes

Ehrenfest Model -

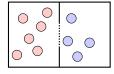
 A simple model for the exchange of molecules between two boxes



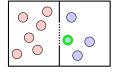
- A simple model for the exchange of molecules between two boxes
- We have *d* particles



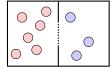
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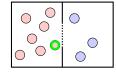
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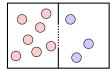


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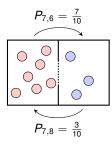
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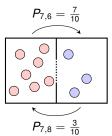
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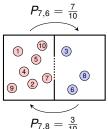


Let us now enlarge the state space by looking at each particle individually!

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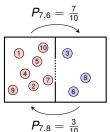


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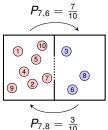
Random Walk on the Hypercube

The Ehrenfest Markov Chain

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Random Walk on the Hypercube —

• For each particle an indicator variable $\Rightarrow \Omega = \{0, 1\}^d$

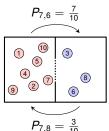


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(Non-Lazy) Random Walk on the Hypercube

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■ At each step *t* = 0, 1, 2 . . .

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- At each step t = 0, 1, 2 . . .
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Lazy Random Walk (2nd Version) -

- At each step t = 0, 1, 2...
 - Pick a random coordinate in [d]
 - Set coordinate to {0, 1} uniformly.

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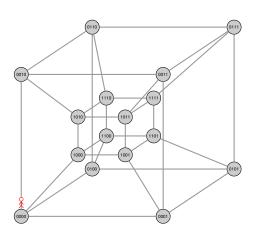
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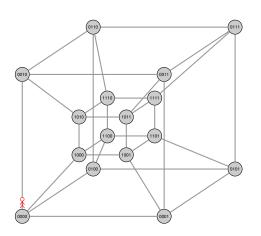
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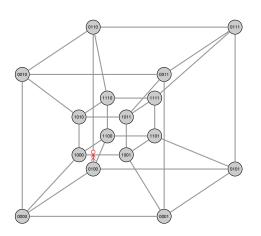


These two chains are equivalent!

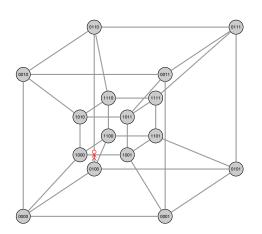




t	Coord.	X_t			
)	2	0	0	0	
1		0	?	0	Ī



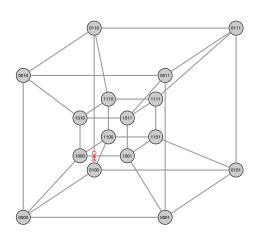
t	Coord.	X_t			
0	2	0	0	0	(
1		0	1	0	(



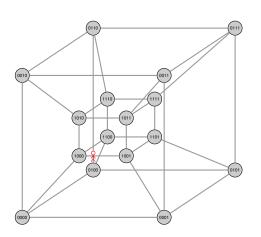
t	Coord.	
)	2	0
1	3	0
2		0

0	0	0	0
0	1	0	0
0	1	?	0

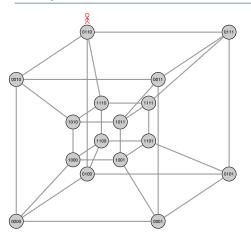
 X_t



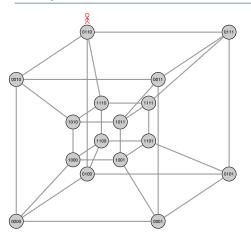
t .	Coord.	X		
)	2	0	0	
	3	0	1	
2		0	1	



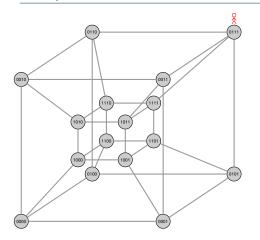
t	Coord.	X_t			
0	2	0	0	0	0
1	3	0	1	0	0
2	3	0	1	0	0
3		0	1	?	0



t	Coord.	X_t			
)	2	0	0	0	0
I	3	0	1	0	0
2	3	0	1	0	0
3		0	1	1	Λ

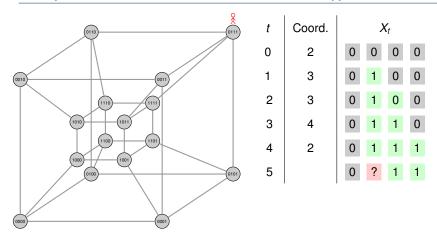


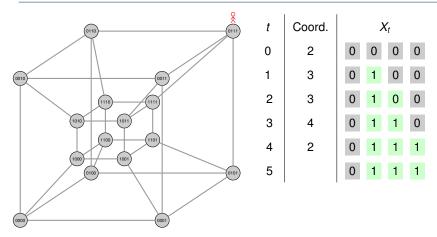
t	Coord.	X_t			
)	2	0	0	0	(
1	3	0	1	0	(
2	3	0	1	0	(
3	4	0	1	1	(

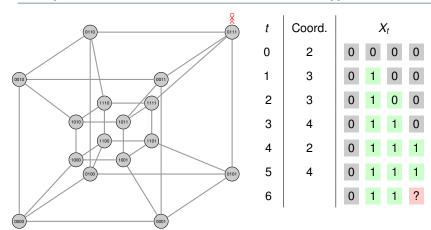


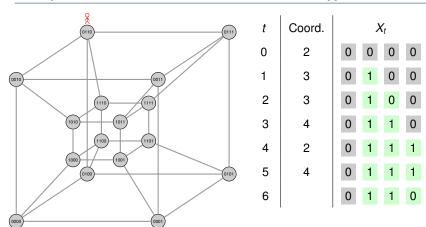
t	Coord.	
0	2	0
1	3	0
2	3	0
3	4	0
_		

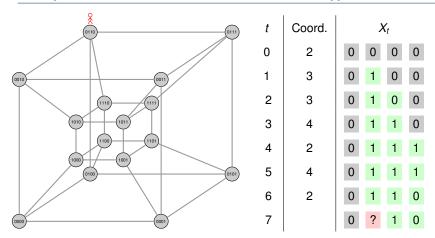
X_t						
0	0	0	0			
0	1	0	0			
0	1	0	0			
0	1	1	0			
0	1	1	1			

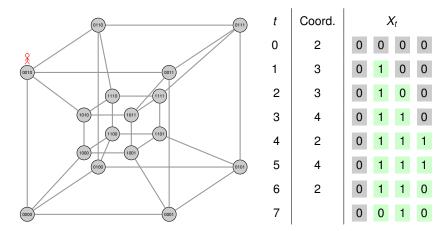


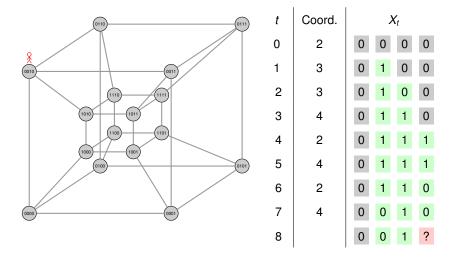


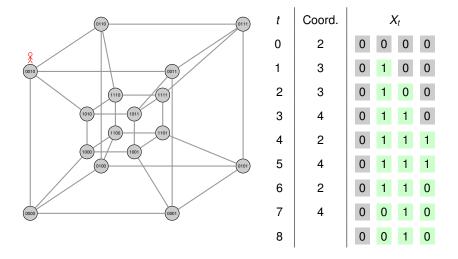


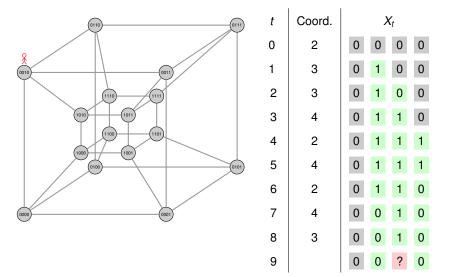


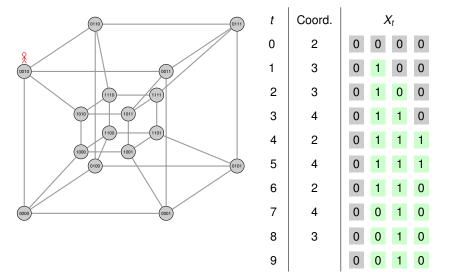


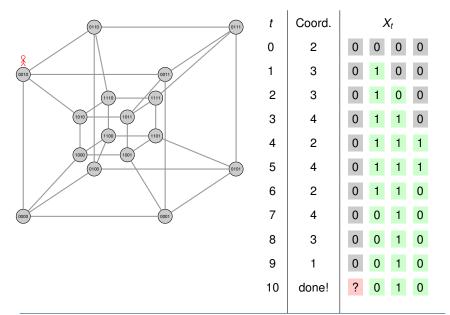


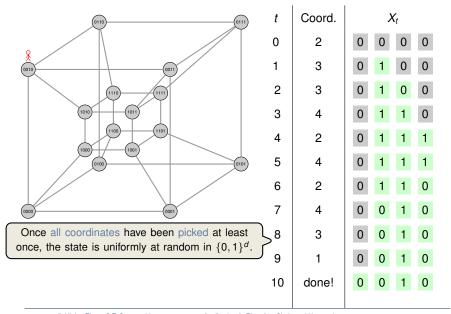


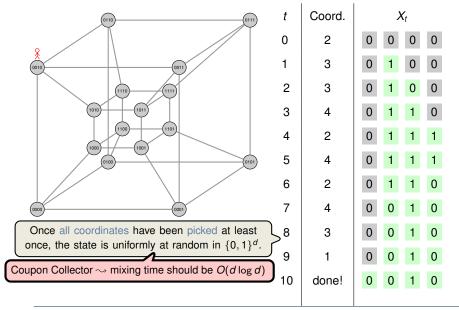


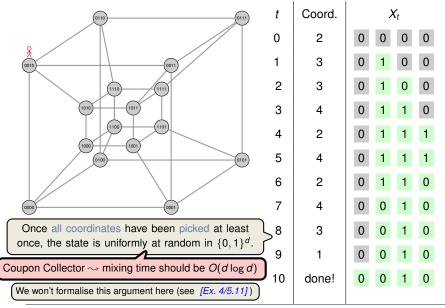




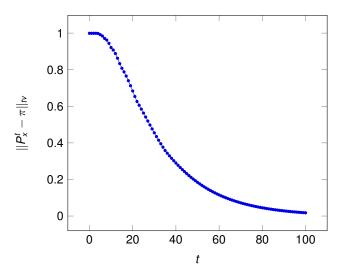




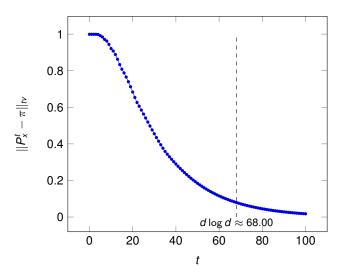




Total Variation Distance of Random Walk on Hypercube (d = 22)



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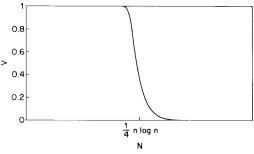


Fig. 1. The variation distance V as a function of N, for $n = 10^{12}$.

Source: "Asymptotic analysis of a random walk on a hypercube with many dimensions", P. Diaconis, R.L. Graham, J.A. Morrison; Random Structures & Algorithms, 1990.



53

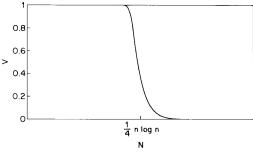


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- The variation distance exhibits a so-called cut-off phenomena:



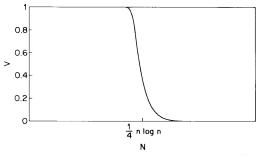


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 - Distance remains close to its maximum value 1 until step $\frac{1}{4}n \log n \Theta(n)$
 - Then distance moves close to 0 before step $\frac{1}{4}n \log n + \Theta(n)$

Outline

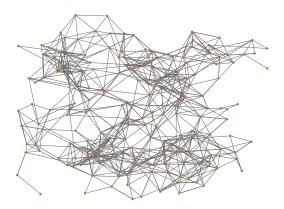
Application 3: Ehrenfest Chain and Hypercubes

Random Walks on Graphs, Hitting Times and Cover Times

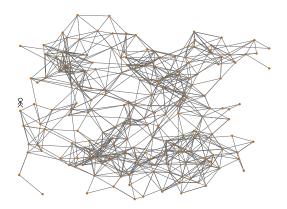
Random Walks on Paths and Grids

SAT and a Randomised Algorithm for 2-SAT

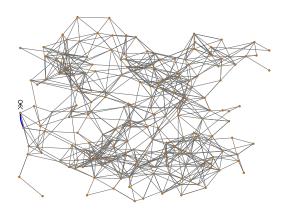
$$P(u,v) = egin{cases} rac{1}{\deg(u)} & ext{if } \{u,v\} \in E, \ 0 & ext{if } \{u,v\}
ot\in E. \end{cases}, \qquad ext{and} \qquad \pi(u) = rac{\deg(u)}{2|E|}$$



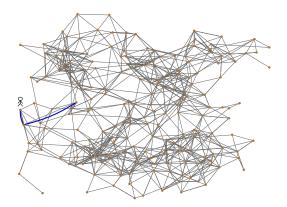
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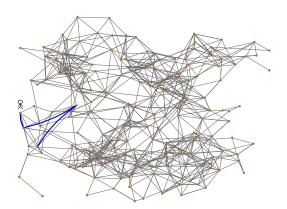
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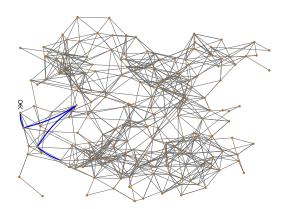
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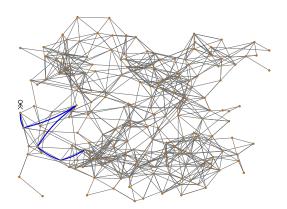
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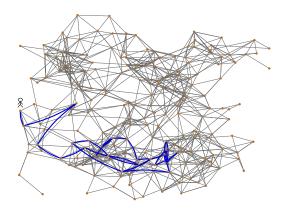
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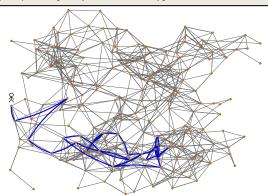
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A Simple Random Walk (SRW) on a graph G is a Markov chain on V(G) with

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Recall: $h(u, v) = \mathbf{E}_u[\min\{t \ge 1 : X_t = v\}]$ is the hitting time of v from u.



The Lazy Random Walk (LRW) on G given by $\widetilde{P} = (P + I)/2$,

$$\widetilde{P}_{u,v} = egin{cases} rac{1}{2\,\deg(u)} & \text{if } \{u,v\} \in E, \\ rac{1}{2} & \text{if } u = v, \\ 0 & \text{otherwise.} \end{cases}$$

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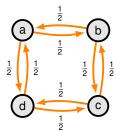
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Fact: For any graph G the LRW on G is aperiodic.

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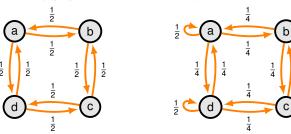


SRW on C4, Periodic

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SRW on C₄, Periodic

LRW on C₄, Aperiodic

Outline

Application 3: Ehrenfest Chain and Hypercubes

Random Walks on Graphs, Hitting Times and Cover Times

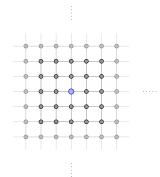
Random Walks on Paths and Grids

SAT and a Randomised Algorithm for 2-SAT

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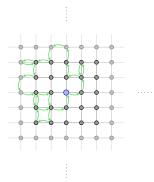
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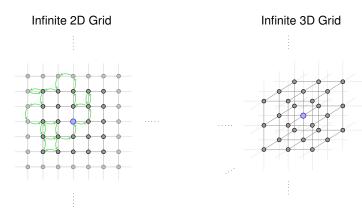


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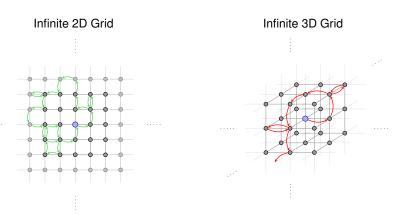
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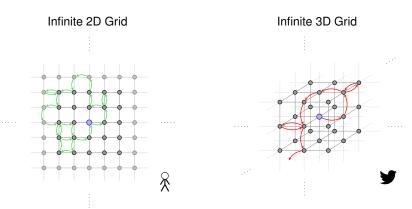
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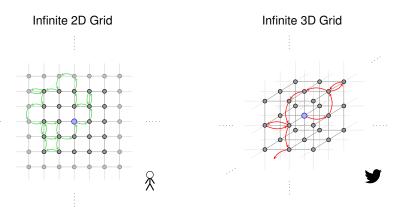


Will a random walk always return to the origin?



"A drunk man will find his way home, but a drunk bird may get lost forever."

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But for any regular (finite) graph, the expected return time to u is $1/\pi(u) = n$

SRW Random Walk on Two-Dimensional Grids: Animation

The *n*-path P_n is the graph with $V(P_n) = [0, n], E(P_n) = \{\{i, j\} : j = i + 1\}.$

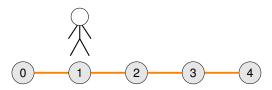


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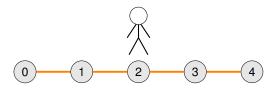
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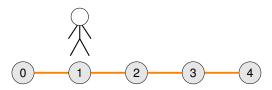
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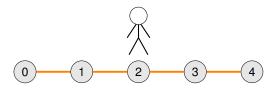
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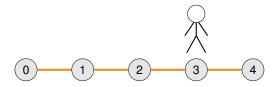
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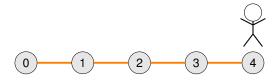
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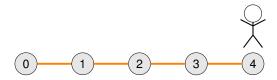
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Exercise: [Exercise 4/5.15] What happens for the LRW on P_n ?

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Outline

Application 3: Ehrenfest Chain and Hypercubes

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SAT and a Randomised Algorithm for 2-SAT

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 - → Model checking and hardware/software verification
 - → Design of experiments
 - → Classical planning
 - $\rightarrow \dots$

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F T T F F F T T

$$0$$

$$1$$

$$2$$

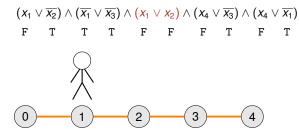
$$3$$

$\alpha =$	(T,	T,	F,	T)).

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F

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t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
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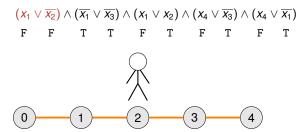
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F F T T F T F T
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F O 1 2 3

α	= ((Τ,	Т,	F,	T)	١.

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
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$\alpha -$	ι.,	Ι,	٠,	- /	,.

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$$F \quad F \quad T \quad F \quad T \quad F \quad T$$

$$0 \quad 1 \quad 2 \quad 3 \quad 4$$

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1	F	Т	F	F
2	Т	Т	F	F

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee \overline{x_3}) \wedge (x_4 \vee \overline{x_1})$$
T F F T T T F F
$$(x_1 \vee \overline{x_2}) \wedge (x_4 \vee \overline{x_3}) \wedge (x_4 \vee \overline{x_1})$$

$\alpha =$	(T,	Т,	F,	T)	١.
$\alpha -$	ι.,	Ι,	٠,	- /	,.

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	Т	F	F
2	Т	Т	F	F

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee \overline{x_3}) \wedge (x_4 \vee \overline{x_1})$$
T F F T T T T T F

$\alpha =$	(T,	Т,	F,	T)).	

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	Т	F	F
2	Т	Т	F	F
3	Т	Т	F	Т

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

Example 1 : Solution Found

$\alpha =$	(T,	Т,	F,	T)	١.
------------	-----	----	----	----	----

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	T	F	F
2	Т	T	F	F
3	Т	T	F	Т

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee x_3) \wedge (x_4 \vee \overline{x_1})$$
F T T F F F F T

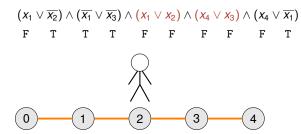
$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee x_3) \wedge (x_4 \vee \overline{x_1})$$
F T T T T F F F F T

$\alpha = (1, F, F, 1)$.	$\alpha =$	(T,	F,	F,	T)	١.
---------------------------	------------	-----	----	----	----	----

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.



$\alpha = 0$	(Τ,	F,	F,	T)	١.

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$\alpha =$	(T,	F,	F,	T)	١.

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$(x_1 \lor \overline{x_2}) \land (\overline{x_1} \lor \overline{x_3}) \land (x_1 \lor x_2) \land (x_4 \lor x_3) \land (x_4 \lor \overline{x_1})$$
F T T F F F T
$$0$$

$$0$$

$$1$$

$$2$$

$$3$$

$\alpha =$	(T,	F,	F,	T)).

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$\alpha = (T, F, F, T).$$

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	T

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$\alpha =$	(T,	F,	F,	T)).

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	T

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$\alpha =$	(T,	F,	F,	T)).

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	Т

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$\alpha =$	(T,	F, F	', T)	١.

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	Т
2	F	T	F	Т

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$\alpha = (1, F, F, 1)$.	$\alpha =$	(T,	F,	F,	T)	١.
---------------------------	------------	-----	----	----	----	----

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	Т
2	F	Т	F	T

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee x_3) \wedge (x_4 \vee \overline{x_1})$$

$$F \quad F \quad T \quad T \quad F \quad T \quad T \quad T$$

$$0 \quad 1 \quad 2 \quad 3 \quad 4$$

$\alpha =$	(T,	F,	F,	T)).

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	Т
2	F	T	F	T

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n² times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee x_3) \wedge (x_4 \vee \overline{x_1})$$

$$T \quad F \quad F \quad T \quad T \quad T \quad F \quad T \quad F$$

$$0 \quad 1 \quad 2 \quad 3 \quad 4$$

$\alpha =$	(T,	F,	F,	T)	١.
------------	-----	----	----	----	----

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	T
2	F	Т	F	T
3	Т	Т	F	Т

RANDOMISED-2-SAT (Input: a 2-SAT-Formula)

- 1: Start with an arbitrary truth assignment
- 2: Repeat up to 2n2 times
- 3: Pick an arbitrary unsatisfied clause
- 4: Choose a random literal and switch its value
- 5: If formula is satisfied then return "Satisfiable"
- 6: return "Unsatisfiable"
- Call each loop of (2) a step. Let A_i be the variable assignment at step i.
- Let α be any solution and $X_i = |\text{variable values shared by } A_i \text{ and } \alpha|$.

Example 2: (Another) Solution Found

$$(x_1 \vee \overline{x_2}) \wedge (\overline{x_1} \vee \overline{x_3}) \wedge (x_1 \vee x_2) \wedge (x_4 \vee x_3) \wedge (x_4 \vee \overline{x_1})$$

$$T \quad F \quad F \quad T \quad T \quad T \quad F \quad T \quad F$$

$$0 \quad 1 \quad 2 \quad 3 \quad 4$$

$\alpha =$	(T,	F,	F,	T)).

t	<i>X</i> ₁	<i>X</i> ₂	<i>X</i> ₃	<i>X</i> ₄
0	F	F	F	F
1	F	F	F	Т
2	F	T	F	T
3	Т	Т	F	Т

Expected iterations of (2) in RANDOMISED-2-SAT =

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Expected iterations of (2) in RANDOMISED-2-SAT -

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Expected iterations of (2) in RANDOMISED-2-SAT -

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

(i)
$$P[X_{i+1} = 1 \mid X_i = 0] = 1$$

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] \ge 1/2$

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] \ge 1/2$
- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2.$

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Proof: Fix any solution α , then for any $i \ge 0$ and $1 \le k \le n-1$,

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] \ge 1/2$
- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2$.

Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Proof: Fix any solution α , then for any $i \ge 0$ and $1 \le k \le n-1$,

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
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- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2$.

Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

Assume (pessimistically) that $X_0 = 0$ (none of our initial guesses is right).

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

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- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
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Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

Assume (pessimistically) that $X_0 = 0$ (none of our initial guesses is right).

The process X_i is complicated to describe in full; however by (i) - (iii) we can **bound** it by Y_i (SRW on the n-path from 0).

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Proof: Fix any solution α , then for any $i \ge 0$ and $1 \le k \le n-1$,

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] \ge 1/2$
- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2$.

Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

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The process X_i is complicated to describe in full; however by (i) - (iii) we can **bound** it by Y_i (SRW on the *n*-path from 0). This gives (see also [Ex 4/5.16])

$$\mathbf{E}[\text{time to find sol}] \leq \mathbf{E}_0[\min\{t : X_t = n\}] \leq \mathbf{E}_0[\min\{t : Y_t = n\}] = h(0, n) = n^2.$$

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Proof: Fix any solution α , then for any $i \ge 0$ and $1 \le k \le n-1$,

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] > 1/2$
- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2$.

Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

Assume (pessimistically) that $X_0 = 0$ (none of our initial guesses is right).

The process X_i is complicated to describe in full; however by (i) - (iii) we can **bound** it by Y_i (SRW on the *n*-path from 0). This gives (see also [Ex 4/5.16])

$$\mathbf{E}[\text{time to find sol}] \le \mathbf{E}_0[\min\{t : X_t = n\}] \le \mathbf{E}_0[\min\{t : Y_t = n\}] = h(0, n) = n^2.$$

Running for $2n^2$ steps and using Markov's inequality yields:

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before RANDOMISED-2-SAT outputs a valid solution is at most n^2 .

Proof: Fix any solution α , then for any $i \ge 0$ and $1 \le k \le n-1$,

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] > 1/2$
- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2.$

Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

Assume (pessimistically) that $X_0 = 0$ (none of our initial guesses is right).

The process X_i is complicated to describe in full; however by (i) - (iii) we can **bound** it by Y_i (SRW on the *n*-path from 0). This gives (see also [Ex 4/5.16])

$$\mathbf{E}[\text{time to find sol}] \le \mathbf{E}_0[\min\{t : X_t = n\}] \le \mathbf{E}_0[\min\{t : Y_t = n\}] = h(0, n) = n^2.$$

Running for 2n² steps and using Markov's inequality yields:

Proposition

Provided a solution exists, RANDOMISED-2-SAT will return a valid solution in $O(n^2)$ steps with probability at least 1/2.

Expected iterations of (2) in RANDOMISED-2-SAT

If the formula is satisfiable, then the expected number of steps before Randomised-2-SAT outputs a valid solution is at most n^2 .

Proof: Fix any solution α , then for any $i \ge 0$ and $1 \le k \le n-1$,

- (i) $P[X_{i+1} = 1 \mid X_i = 0] = 1$
- (ii) $P[X_{i+1} = k+1 \mid X_i = k] \ge 1/2$
- (iii) $P[X_{i+1} = k-1 \mid X_i = k] \leq 1/2$.

Notice that if $X_i = n$ then $A_i = \alpha$ thus solution found (may find another first).

Assume (pessimistically) that $X_0 = 0$ (none of our initial guesses is right).

The process X_i is complicated to describe in full; however by (i) - (iii) we can **bound** it by Y_i (SRW on the n-path from 0). This gives (see also [Ex 4/5.16])

E[time to find sol]
$$\leq$$
 E₀[min{ $t: X_t = n$ }] \leq **E**₀[min{ $t: Y_t = n$ }] = $h(0, n) = n^2$.



Exercise: (difficult, beyond this course) What happens to the above analysis if we apply the same algorithm to 3-SAT?

Boosting Lemma

Suppose a randomised algorithm succeeds with probability (at least) p. Then for any $C \ge 1$, $\lceil \frac{C}{p} \cdot \log n \rceil$ repetitions are sufficient to succeed (in at least one repetition) with probability at least $1 - n^{-C}$.

Boosting Lemma

Suppose a randomised algorithm succeeds with probability (at least) p. Then for any $C \geq 1$, $\lceil \frac{C}{p} \cdot \log n \rceil$ repetitions are sufficient to succeed (in at least one repetition) with probability at least $1 - n^{-C}$.

Proof: Recall that $1 - p \le e^{-p}$ for all real p. Let $t = \lceil \frac{C}{p} \log n \rceil$ and observe

$$\begin{aligned} \mathbf{P} [t \text{ runs all fail}] &\leq (1-p)^t \\ &\leq e^{-pt} \\ &\leq n^{-C}, \end{aligned}$$

thus the probability one of the runs succeeds is at least $1 - n^{-C}$.

Boosting Lemma

Suppose a randomised algorithm succeeds with probability (at least) p. Then for any $C \geq 1$, $\lceil \frac{C}{p} \cdot \log n \rceil$ repetitions are sufficient to succeed (in at least one repetition) with probability at least $1 - n^{-C}$.

Proof: Recall that $1 - p \le e^{-p}$ for all real p. Let $t = \lceil \frac{C}{p} \log n \rceil$ and observe

$$P[t \text{ runs all fail}] \le (1 - p)^t$$

$$\le e^{-pt}$$

$$\le n^{-C},$$

thus the probability one of the runs succeeds is at least $1 - n^{-C}$.

- RANDOMISED-2-SAT

There is a $O(n^2 \log n)$ -step algorithm for 2-SAT which succeeds w.h.p.

Randomised Algorithms

Lecture 6: Linear Programming: Introduction

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Boosting Success Probabilities (Last Lecture)

Introduction

A Simple Example of a Linear Program

Formulating Problems as Linear Programs

Standard and Slack Forms

Boosting Lemma

Suppose a randomised algorithm succeeds with probability (at least) p. Then for any $C \ge 1$, $\lceil \frac{C}{p} \cdot \log n \rceil$ repetitions are sufficient to succeed (in at least one repetition) with probability at least $1 - n^{-C}$.

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Proof: Recall that $1 - p \le e^{-p}$ for all real p. Let $t = \lceil \frac{C}{p} \log n \rceil$ and observe

$$\begin{aligned} \mathbf{P} [t \text{ runs all fail}] &\leq (1-p)^t \\ &\leq e^{-pt} \\ &\leq n^{-C}, \end{aligned}$$

thus the probability one of the runs succeeds is at least $1 - n^{-C}$.

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- RANDOMISED-2-SAT

There is a $O(n^2 \log n)$ -step algorithm for 2-SAT which succeeds w.h.p.

Outline

Boosting Success Probabilities (Last Lecture)

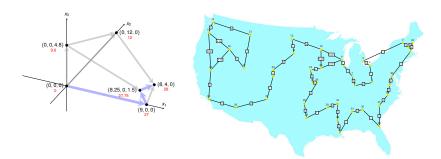
Introduction

A Simple Example of a Linear Program

Formulating Problems as Linear Programs

Standard and Slack Forms

Introduction



- linear programming is a powerful tool in optimisation
- inspired more sophisticated techniques such as quadratic optimisation, convex optimisation, integer programming and semi-definite programming
- we will later use the connection between linear and integer programming to tackle several problems (Vertex-Cover, Set-Cover, TSP, satisfiability)

Outline

Boosting Success Probabilities (Last Lecture)

Introduction

A Simple Example of a Linear Program

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Standard and Slack Forms

What are Linear Programs?

Linear Programming (informal definition)

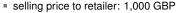
- maximise or minimise an objective, given limited resources (competing constraints)
- constraints are specified as (in)equalities
- objective function and constraints are linear

Laptop

- Laptop
 - selling price to retailer: 1,000 GBP

- Laptop
 - selling price to retailer: 1,000 GBP
 - glass: 4 units

Laptop

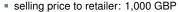


glass: 4 units

copper: 2 units



Laptop



glass: 4 units

copper: 2 units

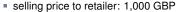
rare-earth elements: 1 unit







Laptop



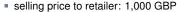
glass: 4 units

copper: 2 units

rare-earth elements: 1 unit

Smartphone

Laptop



glass: 4 units

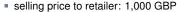
copper: 2 units

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Smartphone

selling price to retailer: 1,000 GBP

Laptop



- glass: 4 units
- copper: 2 units
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- selling price to retailer: 1,000 GBP
- glass: 1 unit

Laptop

selling price to retailer: 1,000 GBP

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rare-earth elements: 1 unit





selling price to retailer: 1,000 GBP

glass: 1 unit

copper: 1 unit

Laptop

selling price to retailer: 1,000 GBP

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copper: 2 units

rare-earth elements: 1 unit











selling price to retailer: 1,000 GBP

glass: 1 unit copper: 1 unit

rare-earth elements: 2 units

- Laptop
 - selling price to retailer: 1,000 GBP
 - glass: 4 units
 - copper: 2 units
 - rare-earth elements: 1 unit









- Smartphone
 - selling price to retailer: 1,000 GBP
 - glass: 1 unit copper: 1 unit
 - rare-earth elements: 2 units
- You have a daily supply of:

- Laptop
 - selling price to retailer: 1,000 GBP
 - glass: 4 units copper: 2 units
 - rare-earth elements: 1 unit
- Smartphone
 - selling price to retailer: 1,000 GBP
 - glass: 1 unit copper: 1 unit
 - rare-earth elements: 2 units
- You have a daily supply of:
 - glass: 20 units

















Laptop

selling price to retailer: 1,000 GBP

glass: 4 unitscopper: 2 units

rare-earth elements: 1 unit



Smartphone

selling price to retailer: 1,000 GBP

glass: 1 unit copper: 1 unit

rare-earth elements: 2 units



de de de de de de de de de de

You have a daily supply of:glass: 20 unitscopper: 10 units

Laptop

- selling price to retailer: 1,000 GBP
- glass: 4 unitscopper: 2 units
- rare-earth elements: 1 unit



Smartphone

- selling price to retailer: 1,000 GBP
- glass: 1 unit copper: 1 unit
- rare-earth elements: 2 units



- glass: 20 unitscopper: 10 units
- rare-earth elements: 14 units



506 506 506 506

Laptop

- selling price to retailer: 1,000 GBP
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Smartphone

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- rare-earth elements: 14 units
- (and enough of everything else...)



506 506 506 506

Laptop

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- (and enough of everything else...)



506 506 506 506

How to maximise your daily earnings?

Linear Program for the Production Problem ———

Linear Program for the Production Problem =

The solution of this linear program yields the optimal production schedule.

Linear Program for the Production Problem

The solution of this linear program yields the optimal production schedule.

Formal Definition of Linear Program —————

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Given a₁, a₂,..., a_n and a set of variables x₁, x₂,..., x_n, a linear function f is defined by

$$f(x_1, x_2, ..., x_n) = a_1x_1 + a_2x_2 + \cdots + a_nx_n.$$

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- Linear Equality: $f(x_1, x_2, ..., x_n) = b$
- Linear Inequality: $f(x_1, x_2, ..., x_n) \le b$

Linear Program for the Production Problem ———

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Formal Definition of Linear Program ——

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Linear Program for the Production Problem ———

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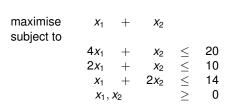
Formal Definition of Linear Program ——

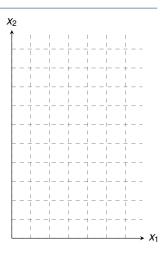
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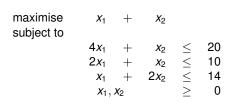
$$f(x_1, x_2, ..., x_n) = a_1x_1 + a_2x_2 + ... + a_nx_n.$$

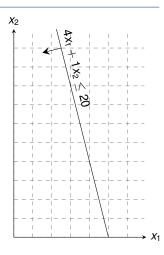
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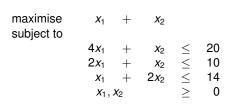
 Linear Proces
- Linear-Progamming Problem: either minimise or maximise a linear function subject to a set of linear constraints

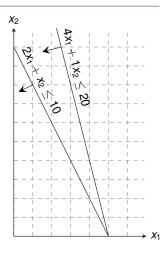


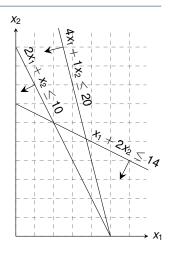


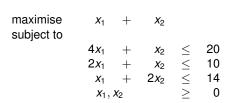


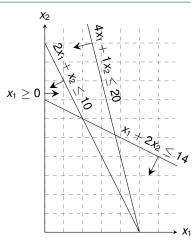


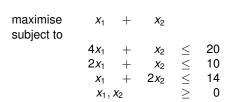


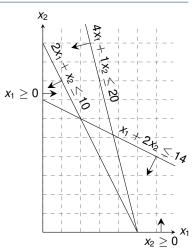


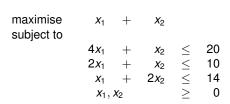


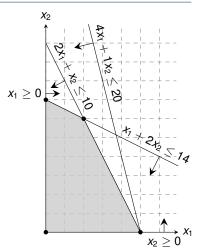


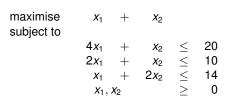


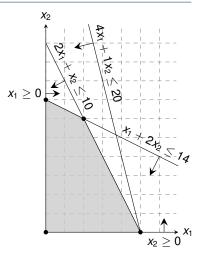


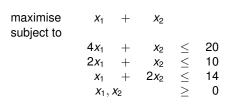


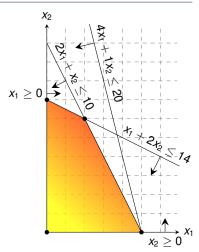


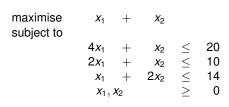


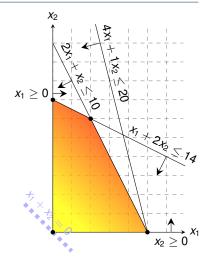


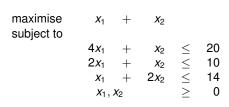


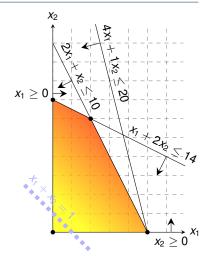


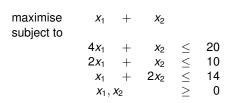


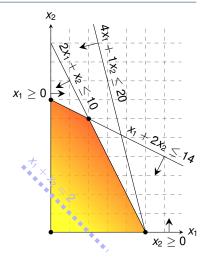


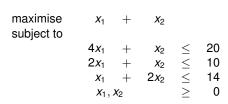


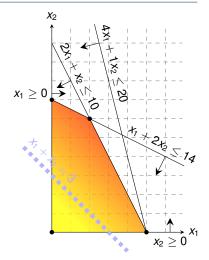


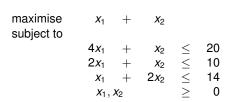


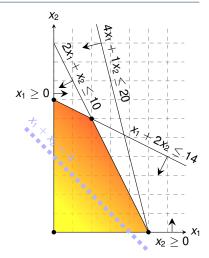


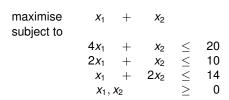


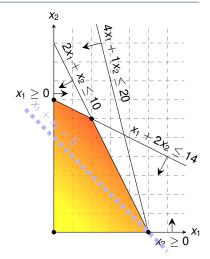


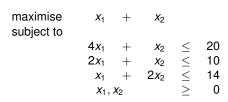


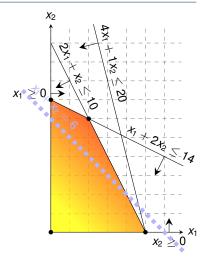


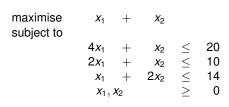


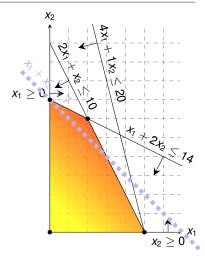


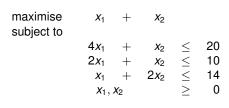


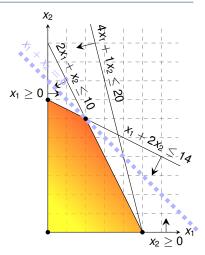


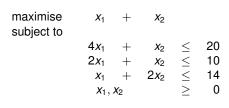


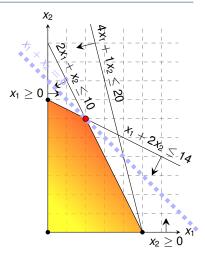


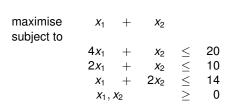




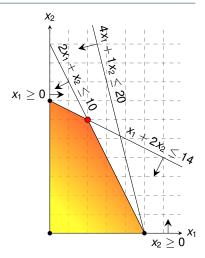






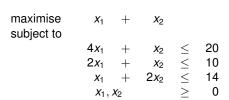


Graphical Procedure: Move the line $x_1 + x_2 = z$ as far up as possible.

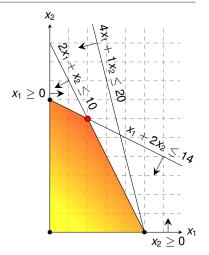




Question: Which aspect did we ignore in the formulation of the linear program?



Graphical Procedure: Move the line $x_1 + x_2 = z$ as far up as possible.



While the same approach also works for higher-dimensions, we need to take a more systematic and algebraic procedure.

Outline

Boosting Success Probabilities (Last Lecture)

Introduction

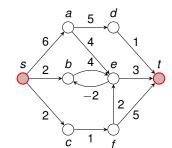
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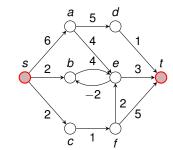
Single-Pair Shortest Path Problem

■ Given: directed graph G = (V, E) with edge weights $w : E \to \mathbb{R}$, pair of vertices $s, t \in V$



- Single-Pair Shortest Path Problem

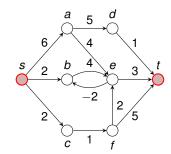
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- Goal: Find a path of minimum weight from s to t in G



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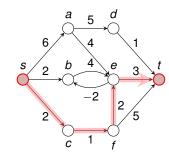
$$p = (v_0 = s, v_1, \dots, v_k = t)$$
 such that $w(p) = \sum_{i=1}^k w(v_{k-1}, v_k)$ is minimised.



Single-Pair Shortest Path Problem

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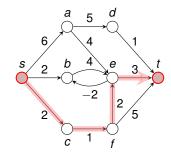
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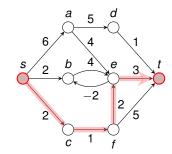


Exercise: Translate the SPSP problem into a linear program!

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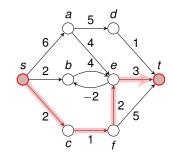
Shortest Paths as LP -

subject to

Single-Pair Shortest Path Problem =

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Shortest Paths as LP -

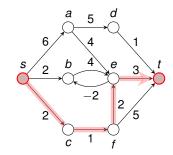
subject to

$$egin{array}{lcl} d_v & \leq & d_u & + & w(u,v) & ext{for each edge } (u,v) \in E, \ d_s & = & 0. \end{array}$$

Single-Pair Shortest Path Problem -

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Shortest Paths as I P =

$$d_t$$

$$d_v \le d_u + w(u,v)$$
 for each edge $(u,v) \in E$, $d_s = 0$.

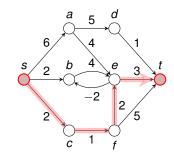
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Shortest Paths as I P =

maximise subject to d_t

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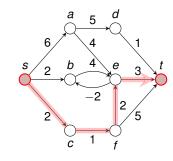
this is a maxim-

isation problem!

Single-Pair Shortest Path Problem

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 such that $w(p) = \sum_{i=1}^k w(v_{k-1}, v_k)$ is minimised.



Shortest Paths as LP maximise d₊ subject to

Recall: When Bellman-Ford terminates. all these inequalities are satisfied.

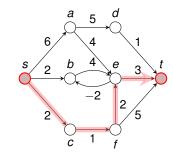
 $\leq d_u + w(u,v)$ for each edge $(u, v) \in E$,

this is a maximisation problem!

Single-Pair Shortest Path Problem

- Given: directed graph G = (V, E) with edge weights $w : E \to \mathbb{R}$, pair of vertices $s, t \in V$
- Goal: Find a path of minimum weight from s to t in G

$$p = (v_0 = s, v_1, \dots, v_k = t)$$
 such that $w(p) = \sum_{i=1}^k w(v_{k-1}, v_k)$ is minimised.



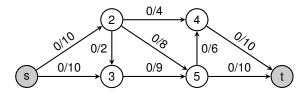
Shortest Paths as LP Recall: When Bellman-Ford terminates, all these inequalities are satisfied. Solution \overline{d} satisfies $\overline{d}_v = \min_{u: (u,v) \in \mathcal{E}} \{\overline{d}_u + w(u,v)\}$

Maximum Flow Problem

• Given: directed graph G=(V,E) with edge capacities $c:E\to\mathbb{R}^+$ (recall c(u,v)=0 if $(u,v)\not\in E$), pair of vertices $s,t\in V$

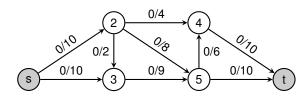
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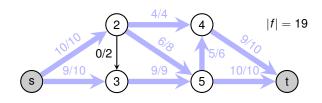
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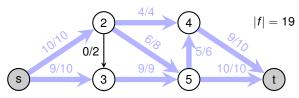
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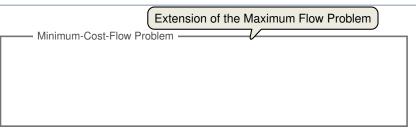


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Maximum Flow as LP



Extension of the Maximum Flow Problem

Minimum-Cost-Flow Problem -

■ Given: directed graph G = (V, E) with capacities $c : E \to \mathbb{R}^+$, pair of vertices $s, t \in V$, cost function $a : E \to \mathbb{R}^+$, flow demand of d units

Extension of the Maximum Flow Problem

Minimum-Cost-Flow Problem -

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- Goal: Find a flow $f: V \times V \to \mathbb{R}$ from s to t with |f| = d while minimising the total cost $\sum_{(u,v) \in F} a(u,v) f_{uv}$ incurred by the flow.

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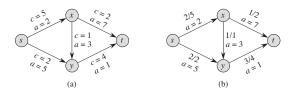


Figure 29.3 (a) An example of a minimum-cost-flow problem. We denote the capacities by c and the costs by a. Vertex s is the source and vertex t is the sink, and we wish to send 4 units of flow from s to t. (b) A solution to the minimum-cost flow problem in which 4 units of flow are sent from s to t. For each edge, the flow and capacity are written as flow/capacity.

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Minimum Cost Flow as a LP

Minimum Cost Flow as LP ——

minimise
$$\sum_{(u,v)\in E} a(u,v) f_{uv}$$
 subject to
$$f_{uv} \leq c(u,v) \quad \text{for } u,v\in V,$$

$$\sum_{v\in V} f_{vu} - \sum_{v\in V} f_{uv} = 0 \quad \text{for } u\in V\setminus \{s,t\},$$

$$\sum_{v\in V} f_{sv} - \sum_{v\in V} f_{vs} = d,$$

$$f_{uv} \geq 0 \quad \text{for } u,v\in V.$$

Minimum Cost Flow as a LP

Minimum Cost Flow as LP

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Real power of Linear Programming comes from the ability to solve **new problems**!

Outline

Boosting Success Probabilities (Last Lecture)

Introduction

A Simple Example of a Linear Program

Formulating Problems as Linear Programs

Standard and Slack Forms

Standard and Slack Forms

Standard Form

maximise
$$\sum_{j=1}^{n} c_j x_j$$
 subject to

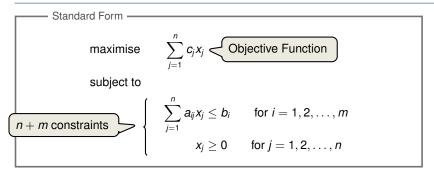
$$\sum_{j=1}^{n} a_{ij} x_{j} \leq b_{i} \quad \text{for } i = 1, 2, \dots, m$$
$$x_{j} \geq 0 \quad \text{for } j = 1, 2, \dots, n$$

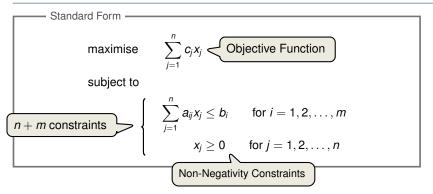
Standard Form

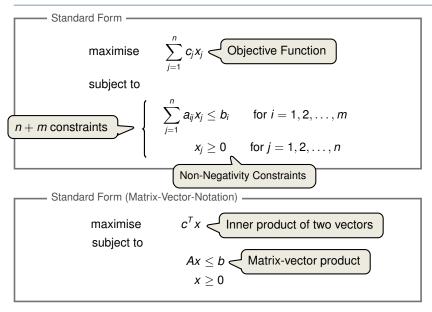
maximise
$$\sum_{j=1}^{n} c_j x_j$$
 Objective Function

subject to

$$\sum_{j=1}^{n} a_{ij} x_{j} \le b_{i} \quad \text{for } i = 1, 2, \dots, n$$
$$x_{j} \ge 0 \quad \text{for } j = 1, 2, \dots, n$$







Converting Linear Programs into Standard Form

Reasons for a LP not being in standard form:

- 1. The objective might be a minimisation rather than maximisation.
- 2. There might be variables without nonnegativity constraints.
- 3. There might be equality constraints.
- 4. There might be inequality constraints (with \geq instead of \leq).

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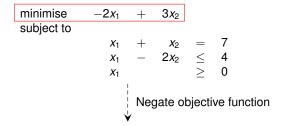
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Equivalence: a correspondence (not necessarily a bijection) between solutions.

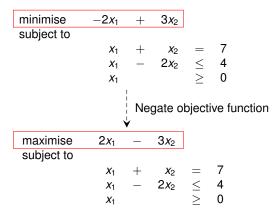
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$$2x_1 - 3x_2$$
 subject to $x_1 + x_2 = 7$ $x_1 - 2x_2 \le 4$ $x_1 \ge 0$

Reasons for a LP not being in standard form:

maximise

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 $2x_{1}$

 $3x_2$

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3. There might be equality constraints.

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maximise subject to

$$2x_1 - 3x_2' + 3x_2''$$

$$\begin{array}{c|cccc} x_1 & + & x_2' & - & x_2'' & = & 7 \\ x_1 & - & 2x_2' & + & 2x_2'' & \leq & 4 \\ x_1, x_2', x_2'' & & \geq & 0 \end{array}$$

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Replace each equality by two inequalities.

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Reasons for a LP not being in standard form:

4. There might be inequality constraints (with \geq instead of \leq).

$$2x_1 - 3x_2' + 3x_2''$$

$$x_1 + x_2' - x_2'' \le 7$$
 $x_1 + x_2' - x_2'' \ge 7$
 $x_1 - 2x_2' + 2x_2'' \le 4$
 $x_1, x_2', x_2'' \ge 0$

Negate respective inequalities.

Reasons for a LP not being in standard form:

4. There might be inequality constraints (with \geq instead of \leq).

$$2x_1 - 3x_2 + 3x_3$$

It is always possible to convert a linear program into standard form.

Goal: Convert standard form into slack form, where all constraints except for the non-negativity constraints are equalities.

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For the simplex algorithm, it is more convenient to work with equality constraints.

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Introducing Slack Variables -

• Let $\sum_{i=1}^{n} a_{ij} x_{ij} \le b_{i}$ be an inequality constraint

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Introducing Slack Variables -

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Introducing Slack Variables -

- Let $\sum_{i=1}^{n} a_{ij}x_{i} \leq b_{i}$ be an inequality constraint
- Introduce a slack variable s by

$$s = b_i - \sum_{j=1}^n a_{ij} x_j$$

Goal: Convert standard form into slack form, where all constraints except for the non-negativity constraints are equalities.

For the simplex algorithm, it is more convenient to work with equality constraints.

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- Let $\sum_{i=1}^{n} a_{ij}x_{i} \leq b_{i}$ be an inequality constraint
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$$s = b_i - \sum_{j=1}^n a_{ij} x_j$$

$$s \geq 0$$
.

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For the simplex algorithm, it is more convenient to work with equality constraints.

- Introducing Slack Variables -

- Let $\sum_{i=1}^{n} a_{ij}x_{i} \leq b_{i}$ be an inequality constraint
- Introduce a slack variable s by

s measures the slack between the two sides of the inequality.

$$s = b_i - \sum_{j=1}^n a_{ij} x_j$$

$$s \ge 0.$$

Goal: Convert standard form into slack form, where all constraints except for the non-negativity constraints are equalities.

For the simplex algorithm, it is more convenient to work with equality constraints.

Introducing Slack Variables

- Let $\sum_{i=1}^{n} a_{ij}x_{i} \leq b_{i}$ be an inequality constraint
- Introduce a slack variable s by

s measures the slack between the two sides of the inequality.

$$s = b_i - \sum_{j=1}^n a_{ij} x_j$$

$$s > 0.$$

• Denote slack variable of the *i*-th inequality by x_{n+i}

subject to

subject to

subject to

maximise subject to

$$2x_1 - 3x_2 + 3x_3$$

$$x_4 = 7 - x_1 - x_2 + x_3$$

$$x_5 = -7 + x_1 + x_2 - x_3$$

$$x_6 = 4 - x_1 + 2x_2 - 2x_3$$

$$x_1, x_2, x_3, x_4, x_5, x_6 \ge 0$$

Use variable z to denote objective function and omit the nonnegativity constraints.

Z	=			$2x_{1}$	_	$3x_{2}$	+	3 <i>x</i> ₃
<i>X</i> ₄	=	7	_	<i>X</i> ₁				
<i>X</i> ₅	=	-7	+	<i>X</i> ₁	+	<i>X</i> ₂	_	<i>X</i> ₃
<i>X</i> ₆	=	4	_	<i>X</i> ₁	+	$2x_{2}$	_	$2x_{3}$

This is called slack form.

$$z = 2x_1 - 3x_2 + 3x_3$$

 $x_4 = 7 - x_1 - x_2 + x_3$
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Basic Variables: $B = \{4, 5, 6\}$

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Non-Basic Variables: $N = \{1, 2, 3\}$

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Slack Form (Formal Definition) ——

Slack form is given by a tuple (N, B, A, b, c, v) so that

$$z = v + \sum_{j \in N} c_j x_j$$

 $x_i = b_i - \sum_{j \in N} a_{ij} x_j$ for $i \in B$,

and all variables are non-negative.

Basic Variables: $B = \{4, 5, 6\}$

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and all variables are non-negative.

Variables/Coefficients on the right hand side are indexed by B and N.

$$z = 28 - \frac{x_3}{6} - \frac{x_5}{6} - \frac{2x_6}{3}$$

$$x_1 = 8 + \frac{x_3}{6} + \frac{x_5}{6} - \frac{x_6}{3}$$

$$x_2 = 4 - \frac{8x_3}{3} - \frac{2x_5}{3} + \frac{x_6}{3}$$

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$$A = \begin{pmatrix} a_{13} & a_{15} & a_{16} \\ a_{23} & a_{25} & a_{26} \\ a_{43} & a_{45} & a_{46} \end{pmatrix} = \begin{pmatrix} -1/6 & -1/6 & 1/3 \\ 8/3 & 2/3 & -1/3 \\ 1/2 & -1/2 & 0 \end{pmatrix}$$

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Slack Form Notation

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$$B = \{1, 2, 4\}, N = \{3, 5, 6\}$$

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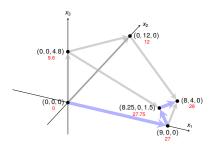
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v = 28



Randomised Algorithms

Lecture 7: Linear Programming: Simplex Algorithm

Thomas Sauerwald (tms41@cam.ac.uk)





Outline

Simplex Algorithm by Example

Details of the Simplex Algorithm

Finding an Initial Solution

Appendix: Cycling and Termination (non-examinable)

Simplex Algorithm: Introduction

Simplex Algorithm ———

- classical method for solving linear programs (Dantzig, 1947)
- usually fast in practice although worst-case runtime not polynomial
- iterative procedure somewhat similar to Gaussian elimination

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Basic Idea:

- Each iteration corresponds to a "basic solution" of the slack form
- All non-basic variables are 0, and the basic variables are determined from the equality constraints
- Each iteration converts one slack form into an equivalent one while the objective value will not decrease
- Conversion ("pivoting") is achieved by switching the roles of one basic and one non-basic variable

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- Each iteration converts one slack form into an equivalent one while the objective value will not decrease In that sense, it is a greedy algorithm.
- Conversion ("pivoting") is achieved by switching the roles of one basic and one non-basic variable

Extended Example: Conversion into Slack Form

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Extended Example: Conversion into Slack Form

$$z = 3x_1 + x_2 + 2x_3$$

 $x_4 = 30 - x_1 - x_2 - 3x_3$
 $x_5 = 24 - 2x_1 - 2x_2 - 5x_3$
 $x_6 = 36 - 4x_1 - x_2 - 2x_3$

$$z = 3x_1 + x_2 + 2x_3$$

 $x_4 = 30 - x_1 - x_2 - 3x_3$
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Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (0, 0, 0, 30, 24, 36)$

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (0, 0, 0, 30, 24, 36)$

This basic solution is **feasible**

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$
Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (0, 0, 0, 30, 24, 36)$
This basic solution is **feasible**
Objective value is 0.

Increasing the value of x_1 would increase the objective value.

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (0, 0, 0, 30, 24, 36)$

This basic solution is **feasible**

Objective value is 0.

Increasing the value of x_1 would increase the objective value.

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

The third constraint is the tightest and limits how much we can increase x_1 .

Increasing the value of x_1 would increase the objective value.

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

The third constraint is the tightest and limits how much we can increase x_1 .

Switch roles of x_1 and x_6 :

Increasing the value of x_1 would increase the objective value.

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

The third constraint is the tightest and limits how much we can increase x_1 .

Switch roles of x_1 and x_6 :

Solving for x₁ yields:

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$
.

Increasing the value of x_1 would increase the objective value.

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

The third constraint is the tightest and limits how much we can increase x_1 .

Switch roles of x_1 and x_6 :

Solving for x₁ yields:

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$
.

• Substitute this into x_1 in the other three equations

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x_6}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_6}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_6}{2}$$

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x_4}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_4}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_4}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_5}{2}$$

Basic solution: $(\overline{x_1},\overline{x_2},\dots,\overline{x_6})=(9,0,0,21,6,0)$ with objective value 27

Increasing the value of x_3 would increase the objective value.

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_4}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_4}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_4}{4}$$

Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (9, 0, 0, 21, 6, 0)$ with objective value 27

Increasing the value of x_3 would increase the objective value.

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x_6}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_6}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_6}{2}$$

The third constraint is the tightest and limits how much we can increase x_3 .

Increasing the value of x_3 would increase the objective value.

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x_6}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_6}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_6}{2}$$

The third constraint is the tightest and limits how much we can increase x_3 .

Switch roles of x_3 and x_5 :

Increasing the value of x_3 would increase the objective value.

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x_6}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_6}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_6}{2}$$

The third constraint is the tightest and limits how much we can increase x_3 .

Switch roles of x_3 and x_5 :

Solving for x₃ yields:

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} - \frac{x_6}{8}$$
.

Increasing the value of x_3 would increase the objective value.

$$z = 27 + \frac{x_2}{4} + \frac{x_3}{2} - \frac{3x_6}{4}$$

$$x_1 = 9 - \frac{x_2}{4} - \frac{x_3}{2} - \frac{x_6}{4}$$

$$x_4 = 21 - \frac{3x_2}{4} - \frac{5x_3}{2} + \frac{x_6}{4}$$

$$x_5 = 6 - \frac{3x_2}{2} - 4x_3 + \frac{x_6}{2}$$

The third constraint is the tightest and limits how much we can increase x_3 .

Switch roles of x_3 and x_5 :

Solving for x₃ yields:

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} - \frac{x_6}{8}$$
.

• Substitute this into x_3 in the other three equations

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_5}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_5}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_6}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

Basic solution: $(\overline{x_1},\overline{x_2},\ldots,\overline{x_6})=(\frac{33}{4},0,\frac{3}{2},\frac{69}{4},0,0)$ with objective value $\frac{111}{4}=27.75$

Increasing the value of x_2 would increase the objective value.

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_5}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (\frac{33}{4}, 0, \frac{3}{2}, \frac{69}{4}, 0, 0)$ with objective value $\frac{111}{4} = 27.75$

Increasing the value of x_2 would increase the objective value.

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_6}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

The second constraint is the tightest and limits how much we can increase x_2 .

Increasing the value of x_2 would increase the objective value.

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_5}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_6}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

The second constraint is the tightest and limits how much we can increase x_2 .

Switch roles of x_2 and x_3 :

Increasing the value of x_2 would increase the objective value.

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_6}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

The second constraint is the tightest and limits how much we can increase x_2 .

Switch roles of x_2 and x_3 :

Solving for x₂ yields:

$$x_2 = 4 - \frac{8x_3}{3} - \frac{2x_5}{3} + \frac{x_6}{3}$$
.

Increasing the value of x_2 would increase the objective value.

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_6}{8} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{16}$$

The second constraint is the tightest and limits how much we can increase x_2 .

Switch roles of x_2 and x_3 :

Solving for x₂ yields:

$$x_2 = 4 - \frac{8x_3}{3} - \frac{2x_5}{3} + \frac{x_6}{3}$$
.

• Substitute this into x_2 in the other three equations

$$z = 28 - \frac{x_3}{6} - \frac{x_5}{6} - \frac{2x_6}{3}$$

$$x_1 = 8 + \frac{x_3}{6} + \frac{x_5}{6} - \frac{x_6}{3}$$

$$x_2 = 4 - \frac{8x_3}{3} - \frac{2x_5}{3} + \frac{x_6}{3}$$

$$x_4 = 18 - \frac{x_3}{2} + \frac{x_5}{2}$$

$$z = 28 - \frac{x_3}{6} - \frac{x_5}{6} - \frac{2x_5}{3}$$

$$x_1 = 8 + \frac{x_3}{6} + \frac{x_5}{6} - \frac{x_6}{3}$$

$$x_2 = 4 - \frac{8x_3}{3} - \frac{2x_5}{3} + \frac{x_6}{3}$$

$$x_4 = 18 - \frac{x_3}{2} + \frac{x_5}{2}$$

Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (8, 4, 0, 18, 0, 0)$ with objective value 28

All coefficients are negative, and hence this basic solution is **optimal**!

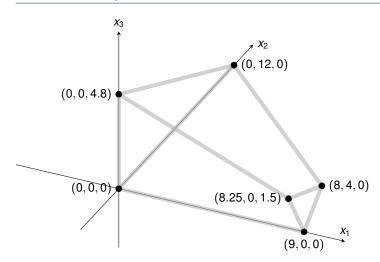
$$z = 28 - \frac{x_3}{6} - \frac{x_5}{6} - \frac{2x_6}{3}$$

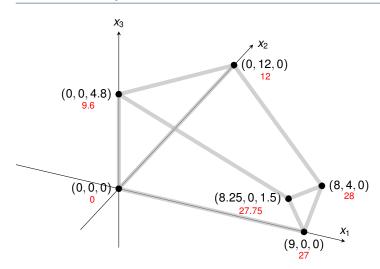
$$x_1 = 8 + \frac{x_3}{6} + \frac{x_5}{6} - \frac{x_6}{3}$$

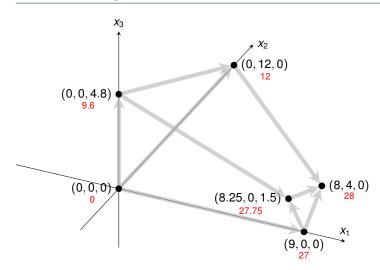
$$x_2 = 4 - \frac{8x_3}{3} - \frac{2x_5}{3} + \frac{x_6}{3}$$

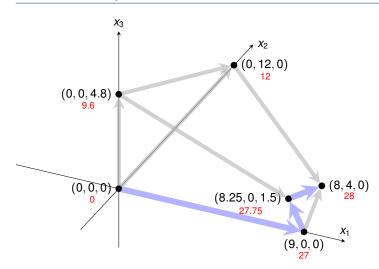
$$x_4 = 18 - \frac{x_3}{3} + \frac{x_5}{2}$$

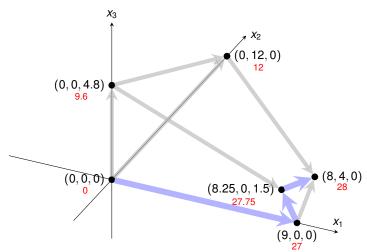
Basic solution: $(\overline{x_1}, \overline{x_2}, \dots, \overline{x_6}) = (8, 4, 0, 18, 0, 0)$ with objective value 28













Exercise: [Ex. 6/7.6] How many basic solutions (including non-feasible ones) are there?

$$z$$
 = $3x_1 + x_2 + 2x_3$
 x_4 = 30 - x_1 - x_2 - $3x_3$
 x_5 = 24 - $2x_1$ - $2x_2$ - $5x_3$
 x_6 = 36 - $4x_1$ - x_2 - $2x_3$

$$z = 3x_1 + x_2 + 2x_3$$

 $x_4 = 30 - x_1 - x_2 - 3x_3$
 $x_5 = 24 - 2x_1 - 2x_2 - 5x_3$
 $x_6 = 36 - 4x_1 - x_2 - 2x_3$

Switch roles of x_1 and x_6 _ - - - -

Switch roles of x_1 and x_{6} ----

$$z = \frac{111}{4} + \frac{x_2}{16} - \frac{x_5}{8} - \frac{11x_6}{16}$$

$$x_1 = \frac{33}{4} - \frac{x_2}{16} + \frac{x_5}{4} - \frac{5x_6}{16}$$

$$x_3 = \frac{3}{2} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$x_4 = \frac{69}{4} + \frac{3x_2}{16} + \frac{5x_5}{8} - \frac{x_6}{8}$$

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$
Switch roles of x_3 and x_5

$$z = \frac{48}{5} + \frac{11x_1}{5} + \frac{x_2}{5} - \frac{2x_5}{5}$$

$$x_4 = \frac{78}{5} + \frac{x_1}{5} + \frac{x_2}{5} + \frac{3x_5}{5}$$

$$x_3 = \frac{24}{5} - \frac{2x_1}{5} - \frac{2x_2}{5} - \frac{x_5}{5}$$

$$x_6 = \frac{132}{5} - \frac{16x_1}{5} - \frac{x_2}{5} + \frac{2x_3}{5}$$
Switch roles of x_1 and x_6 .

Switch roles of x_2 and x_3

$$\frac{1}{3} - \frac{x_2}{16} + \frac{x_5}{8} - \frac{11x_6}{16}$$

$$\frac{3}{4} - \frac{x_2}{16} + \frac{x_5}{8} - \frac{5x_5}{16}$$

$$\frac{3}{4} - \frac{3x_2}{8} - \frac{x_5}{4} + \frac{x_6}{8}$$

$$\frac{3}{4} + \frac{3x_2}{4} + \frac{5x_5}{8} - \frac{x_6}{4}$$

 X_1

 X_4

Extended Example: Alternative Runs (2/2)

$$z = 3x_1 + x_2 + 2x_3$$

$$x_4 = 30 - x_1 - x_2 - 3x_3$$

$$x_5 = 24 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

$$y = 3x_1 + x_2 + 2x_3$$

$$x_6 = 36 - 2x_1 - 2x_2 - 5x_3$$

$$x_6 = 36 - 4x_1 - x_2 - 2x_3$$

$$y = 3x_1 + x_2 + 2x_2 + 2x_3$$

$$x_6 = 36 - 2x_1 - 2x_2 - 5x_3$$

$$x_1 + x_2 + 2x_3 + 3x_5$$

$$x_2 = 48 + \frac{11x_1}{5} + \frac{x_2}{5} - \frac{2x_5}{5}$$

$$x_4 = \frac{78}{5} + \frac{x_1}{5} + \frac{x_2}{5} + \frac{3x_5}{5}$$

$$x_3 = 24 + \frac{24}{5} - \frac{2x_1}{5} - \frac{2x_2}{5} - \frac{x_5}{5}$$

$$x_6 = \frac{132}{5} - \frac{16x_1}{5} - \frac{x_2}{5} + \frac{2x_3}{5}$$
Switch roles of x_1 and x_2

$$x_1 = x_2 + x_3 + x_4 + x_5 + x_5$$

 X_1

 X_4

Outline

Simplex Algorithm by Example

Details of the Simplex Algorithm

Finding an Initial Solution

Appendix: Cycling and Termination (non-examinable)

```
PIVOT(N, B, A, b, c, v, l, e)
      // Compute the coefficients of the equation for new basic variable x_e.
 2. let \widehat{A} be a new m \times n matrix
 3 \hat{b}_e = b_l/a_{le}
 4 for each j \in N - \{e\}
      \hat{a}_{ei} = a_{li}/a_{le}
 6 \hat{a}_{el} = 1/a_{le}
      // Compute the coefficients of the remaining constraints.
 8 for each i \in B - \{l\}
      \hat{b}_i = b_i - a_{i\alpha}\hat{b}_{\alpha}
10 for each j \in N - \{e\}
          \hat{a}_{ii} = a_{ii} - a_{ie}\hat{a}_{ei}
\hat{a}_{il} = -a_{ie}\hat{a}_{el}
13 // Compute the objective function.
14 \quad \hat{v} = v + c_a \hat{b}_a
15 for each j \in N - \{e\}
16
      \hat{c}_i = c_i - c_e \hat{a}_{ei}
17 \quad \hat{c}_l = -c_e \hat{a}_{el}
18 // Compute new sets of basic and nonbasic variables.
19 \hat{N} = N - \{e\} \cup \{l\}
20 \hat{B} = B - \{l\} \cup \{e\}
21 return (\hat{N}, \hat{B}, \hat{A}, \hat{b}, \hat{c}, \hat{v})
```

```
PIVOT(N, B, A, b, c, v, l, e)
      // Compute the coefficients of the equation for new basic variable x_e.
      let \widehat{A} be a new m \times n matrix
 \hat{b}_e = b_l/a_{le}
                                                                                   Rewrite "tight" equation
 4 for each j \in N - \{e\}
                                                                                   for enterring variable x_e.
       \hat{a}_{ei} = a_{li}/a_{le}
 6 \hat{a}_{el} = 1/a_{le}
      // Compute the coefficients of the remaining constraints.
     for each i \in B - \{l\}
       \hat{b}_i = b_i - a_{ia}\hat{b}_a
     for each j \in N - \{e\}
             \hat{a}_{ii} = a_{ii} - a_{ie}\hat{a}_{ei}
     \hat{a}_{il} = -a_{ie}\hat{a}_{el}
     // Compute the objective function.
14 \quad \hat{v} = v + c_a \hat{b}_a
15 for each j \in N - \{e\}
16
      \hat{c}_i = c_i - c_e \hat{a}_{ei}
      \hat{c}_l = -c_e \hat{a}_{el}
18 // Compute new sets of basic and nonbasic variables.
19 \hat{N} = N - \{e\} \cup \{l\}
20 \hat{B} = B - \{l\} \cup \{e\}
21 return (\hat{N}, \hat{B}, \hat{A}, \hat{b}, \hat{c}, \hat{v})
```

```
PIVOT(N, B, A, b, c, v, l, e)
      // Compute the coefficients of the equation for new basic variable x_e.
      let \widehat{A} be a new m \times n matrix
 \hat{b}_e = b_l/a_{le}
                                                                                    Rewrite "tight" equation
 4 for each j \in N - \{e\}
       \hat{a}_{ei} = a_{li}/a_{le}
                                                                                    for enterring variable x_e.
 6 \hat{a}_{el} = 1/a_{le}
      // Compute the coefficients of the remaining constraints.
      for each i \in B - \{l\}
       \hat{b}_i = b_i - a_{i\alpha}\hat{b}_{\alpha}
                                                                                     Substituting x<sub>e</sub> into
     for each j \in N - \{e\}
                                                                                       other equations.
              \hat{a}_{ii} = a_{ii} - a_{ie}\hat{a}_{ei}
      \hat{a}_{il} = -a_{ie}\hat{a}_{el}
      // Compute the objective function.
14 \quad \hat{v} = v + c_a \hat{b}_a
15 for each j \in N - \{e\}
16
           \hat{c}_i = c_i - c_e \hat{a}_{ei}
      \hat{c}_l = -c_e \hat{a}_{el}
18 // Compute new sets of basic and nonbasic variables.
19 \hat{N} = N - \{e\} \cup \{l\}
20 \hat{B} = B - \{l\} \cup \{e\}
21 return (\hat{N}, \hat{B}, \hat{A}, \hat{b}, \hat{c}, \hat{v})
```

```
PIVOT(N, B, A, b, c, v, l, e)
      // Compute the coefficients of the equation for new basic variable x_e.
      let \widehat{A} be a new m \times n matrix
 \hat{b}_e = b_l/a_{le}
                                                                                   Rewrite "tight" equation
 4 for each j \in N - \{e\}
       \hat{a}_{ei} = a_{li}/a_{le}
                                                                                  for enterring variable x_e.
 6 \hat{a}_{el} = 1/a_{le}
      // Compute the coefficients of the remaining constraints.
      for each i \in B - \{l\}
       \hat{b}_i = b_i - a_{i\alpha}\hat{b}_{\alpha}
                                                                                   Substituting x<sub>e</sub> into
      for each j \in N - \{e\}
                                                                                     other equations.
               \hat{a}_{ii} = a_{ii} - a_{ie}\hat{a}_{ei}
      \hat{a}_{il} = -a_{ie}\hat{a}_{el}
      // Compute the objective function.
14 \quad \hat{v} = v + c_a \hat{b}_a
                                                                                   Substituting xe into
15 for each j \in N - \{e\}
16
           \hat{c}_i = c_i - c_e \hat{a}_{ei}
                                                                                    objective function.
      \hat{c}_l = -c_e \hat{a}_{el}
18 // Compute new sets of basic and nonbasic variables.
19 \hat{N} = N - \{e\} \cup \{l\}
20 \hat{B} = B - \{l\} \cup \{e\}
21 return (\hat{N}, \hat{B}, \hat{A}, \hat{b}, \hat{c}, \hat{v})
```

```
PIVOT(N, B, A, b, c, v, l, e)
      // Compute the coefficients of the equation for new basic variable x_e.
     let \widehat{A} be a new m \times n matrix
 \hat{b}_e = b_l/a_{le}
                                                                                Rewrite "tight" equation
 4 for each j \in N - \{e\}
       \hat{a}_{ei} = a_{li}/a_{le}
                                                                                for enterring variable x_e.
 6 \hat{a}_{el} = 1/a_{le}
      // Compute the coefficients of the remaining constraints.
     for each i \in B - \{l\}
       \hat{b}_i = b_i - a_{ia}\hat{b}_a
                                                                                Substituting x<sub>e</sub> into
      for each j \in N - \{e\}
                                                                                  other equations.
               \hat{a}_{ii} = a_{ii} - a_{ie}\hat{a}_{ei}
     \hat{a}_{il} = -a_{ie}\hat{a}_{el}
     // Compute the objective function.
14 \quad \hat{v} = v + c_a \hat{b}_a
                                                                                Substituting xe into
15 for each j \in N - \{e\}
16
          \hat{c}_i = c_i - c_e \hat{a}_{ei}
                                                                                 objective function.
     \hat{c}_l = -c_e \hat{a}_{el}
    // Compute new sets of basic and nonbasic variables.
19 \hat{N} = N - \{e\} \cup \{l\}
                                                                                 Update non-basic
20 \hat{B} = B - \{l\} \cup \{e\}
                                                                                and basic variables
21 return (\hat{N}, \hat{B}, \hat{A}, \hat{b}, \hat{c}, \hat{v})
```

```
PIVOT(N, B, A, b, c, v, l, e)
      // Compute the coefficients of the equation for new basic variable x_e.
     let \widehat{A} be a new m \times n matrix
 \hat{b}_e = b_l/a_{le}
                                                                                Rewrite "tight" equation
    for each j \in N - \{e\} Need that a_{le} \neq 0!
           \hat{a}_{ei} = a_{li}/a_{le}
                                                                               for enterring variable x_e.
 6 \hat{a}_{el} = 1/a_{le}
      // Compute the coefficients of the remaining constraints.
     for each i \in B - \{l\}
       \hat{b}_i = b_i - a_{ia}\hat{b}_a
                                                                                Substituting x<sub>e</sub> into
      for each j \in N - \{e\}
                                                                                  other equations.
               \hat{a}_{ii} = a_{ii} - a_{ie}\hat{a}_{ei}
     \hat{a}_{il} = -a_{ie}\hat{a}_{el}
     // Compute the objective function.
14 \quad \hat{v} = v + c_a \hat{b}_a
                                                                                Substituting xe into
15 for each j \in N - \{e\}
16
      \hat{c}_i = c_i - c_e \hat{a}_{ei}
                                                                                 objective function.
     \hat{c}_l = -c_e \hat{a}_{el}
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                                                                                 Update non-basic
20 \hat{B} = B - \{l\} \cup \{e\}
                                                                                and basic variables
21 return (\hat{N}, \hat{B}, \hat{A}, \hat{b}, \hat{c}, \hat{v})
```

Lemma 29.1

Consider a call to PIVOT(N, B, A, b, c, v, I, e) in which $a_{le} \neq 0$. Let the values returned from the call be $(\widehat{N}, \widehat{B}, \widehat{A}, \widehat{b}, \widehat{c}, \widehat{v})$, and let \overline{x} denote the basic solution after the call. Then

Lemma 29.1

Consider a call to PIVOT(N, B, A, b, c, v, l, e) in which $a_{le} \neq 0$. Let the values returned from the call be $(\widehat{N}, \widehat{B}, \widehat{A}, \widehat{b}, \widehat{c}, \widehat{v})$, and let \overline{x} denote the basic solution after the call. Then

- 1. $\overline{x}_i = 0$ for each $j \in \widehat{N}$.
- 2. $\overline{x}_e = b_l/a_{le}$.
- 3. $\overline{x}_i = b_i a_{ie}\widehat{b}_e$ for each $i \in \widehat{B} \setminus \{e\}$.

Lemma 29.1

Consider a call to PIVOT(N,B,A,b,c,v,l,e) in which $a_{le}\neq 0$. Let the values returned from the call be $(\widehat{N},\widehat{B},\widehat{A},\widehat{b},\widehat{c},\widehat{v})$, and let \overline{x} denote the basic solution after the call. Then

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Proof:

Lemma 29.1

Consider a call to PIVOT(N, B, A, b, c, v, l, e) in which $a_{le} \neq 0$. Let the values returned from the call be $(\widehat{N}, \widehat{B}, \widehat{A}, \widehat{b}, \widehat{c}, \widehat{v})$, and let \overline{x} denote the basic solution after the call. Then

- 1. $\overline{x}_j = 0$ for each $j \in \widehat{N}$.
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Proof:

- 1. holds since the basic solution always sets all non-basic variables to zero.
- 2. When we set each non-basic variable to 0 in a constraint

$$x_i = \widehat{b}_i - \sum_{j \in \widehat{N}} \widehat{a}_{ij} x_j,$$

we have $\overline{x}_i = \hat{b}_i$ for each $i \in \hat{B}$. Hence $\overline{x}_e = \hat{b}_e = b_l/a_{le}$.

3. After substituting into the other constraints, we have

$$\overline{x}_i = \widehat{b}_i = b_i - a_{ie}\widehat{b}_e.$$

Lemma 29.1

Consider a call to PIVOT(N, B, A, b, c, v, l, e) in which $a_{le} \neq 0$. Let the values returned from the call be $(\widehat{N}, \widehat{B}, \widehat{A}, \widehat{b}, \widehat{c}, \widehat{v})$, and let \overline{x} denote the basic solution after the call. Then

- 1. $\overline{x}_j = 0$ for each $j \in \widehat{N}$.
- 2. $\overline{x}_e = b_l/a_{le}$.
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Proof:

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3. After substituting into the other constraints, we have

$$\overline{x}_i = \widehat{b}_i = b_i - a_{ie}\widehat{b}_e$$
.

Formalizing the Simplex Algorithm: Questions

Questions:

- How do we determine whether a linear program is feasible?
- What do we do if the linear program is feasible, but the initial basic solution is not feasible?
- How do we determine whether a linear program is unbounded?
- How do we choose the entering and leaving variables?

Formalizing the Simplex Algorithm: Questions

Questions:

- How do we determine whether a linear program is feasible?
- What do we do if the linear program is feasible, but the initial basic solution is not feasible?
- How do we determine whether a linear program is unbounded?
- How do we choose the entering and leaving variables?

Example before was a particularly nice one!

```
SIMPLEX(A, b, c)
     (N, B, A, b, c, v) = \text{INITIALIZE-SIMPLEX}(A, b, c)
     let \Delta be a new vector of length m
     while some index j \in N has c_i > 0
           choose an index e \in N for which c_e > 0
          for each index i \in B
                if a_{ie} > 0
                     \Delta_i = b_i/a_{ie}
                else \Delta_i = \infty
 9
          choose an index l \in B that minimizes \Delta_i
10
          if \Delta_I == \infty
11
                return "unbounded"
12.
          else (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, e)
     for i = 1 to n
14
          if i \in B
15
               \bar{x}_i = b_i
          else \bar{x}_i = 0
16
     return (\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)
```

```
SIMPLEX(A, b, c)
                                                                            Returns a slack form with a
     (N, B, A, b, c, v) = \text{INITIALIZE-SIMPLEX}(A, b, c)
                                                                        feasible basic solution (if it exists)
     let \Delta be a new vector of length m
     while some index j \in N has c_i > 0
           choose an index e \in N for which c_e > 0
          for each index i \in B
                if a_{ie} > 0
                     \Delta_i = b_i/a_{ie}
                else \Delta_i = \infty
          choose an index l \in B that minimizes \Delta_i
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           else (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, e)
     for i = 1 to n
          if i \in B
14
15
               \bar{x}_i = b_i
          else \bar{x}_i = 0
16
     return (\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)
```

```
SIMPLEX(A, b, c)
                                                                           Returns a slack form with a
     (N, B, A, b, c, v) = INITIALIZE-SIMPLEX(A, b, c)
                                                                       feasible basic solution (if it exists)
    let \Delta be a new vector of length m
    while some index j \in N has c_i > 0
          choose an index e \in N for which c_e > 0
          for each index i \in B
               if a_{ie} > 0
                     \Delta_i = b_i/a_{ie}
               else \Delta_i = \infty
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               return "unbounded"
          else (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, e)
     for i = 1 to n
14
          if i \in B
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               \bar{x}_i = b_i
          else \bar{x}_i = 0
16
     return (\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)
```

```
SIMPLEX(A, b, c)
                                                                          Returns a slack form with a
     (N, B, A, b, c, v) = \text{INITIALIZE-SIMPLEX}(A, b, c)
                                                                      feasible basic solution (if it exists)
    let \Delta be a new vector of length m
    while some index j \in N has c_i > 0
                                                                              Main Loop:
          choose an index e \in N for which c_e > 0
          for each index i \in B
               if a_{ie} > 0
                    \Delta_i = b_i/a_{ie}
               else \Delta_i = \infty
          choose an index l \in B that minimizes \Delta_i
10
          if \Delta_I == \infty
11
               return "unbounded"
          else (N, B, A, b, c, \nu) = \text{PIVOT}(N, B, A, b, c, \nu, l, e)
     for i = 1 to n
14
          if i \in B
15
               \bar{x}_i = b_i
          else \bar{x}_i = 0
16
```

return $(\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)$

```
SIMPLEX(A, b, c)
                                                                            Returns a slack form with a
     (N, B, A, b, c, v) = \text{INITIALIZE-SIMPLEX}(A, b, c)
                                                                        feasible basic solution (if it exists)
     let \Delta be a new vector of length m
    while some index j \in N has c_i > 0
           choose an index e \in N for which c_e > 0
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                if a_{ie} > 0
                     \Delta_i = b_i/a_{ie}
                else \Delta_i = \infty
          choose an index l \in B that minimizes \Delta_i
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           else (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, e)
     for i = 1 to n
14
          if i \in B
15
               \bar{x}_i = b_i
          else \bar{x}_i = 0
16
     return (\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)
```

Main Loop:

- terminates if all coefficients in objective function are non-positive
- Line 4 picks enterring variable x_{α} with positive coefficient
- Lines 6 9 pick the tightest constraint, associated with x1
- Line 11 returns "unbounded" if there are no constraints
- Line 12 calls PIVOT, switching roles of x_i and x_p

```
SIMPLEX(A, b, c)
                                                                          Returns a slack form with a
     (N, B, A, b, c, v) = \text{INITIALIZE-SIMPLEX}(A, b, c)
                                                                      feasible basic solution (if it exists)
    let \Delta be a new vector of length m
    while some index j \in N has c_i > 0
                                                                             Main Loop:
          choose an index e \in N for which c_e > 0

    terminates if all coefficients in

          for each index i \in B
                                                                                  objective function are
               if a_{ie} > 0
                                                                                  non-positive
                    \Delta_i = b_i/a_{ie}
                                                                                  Line 4 picks enterring variable
               else \Delta_i = \infty
                                                                                   x_{\alpha} with positive coefficient
          choose an index l \in B that minimizes \Delta_i
                                                                                ■ Lines 6 — 9 pick the tightest
10
          if \Delta_I == \infty
                                                                                  constraint, associated with x1
11
               return "unbounded"
12
          else (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, e)
                                                                                Line 11 returns "unbounded" if
     for i = 1 to n
                                                                                  there are no constraints
14
          if i \in B
                                                                                Line 12 calls PIVOT, switching
15
               \bar{x}_i = b_i
                                                                                   roles of x_i and x_p
          else \bar{x}_i = 0
16
     return (\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)
```

Return corresponding solution.

```
SIMPLEX(A, b, c)
                                                                            Returns a slack form with a
     (N, B, A, b, c, v) = \text{INITIALIZE-SIMPLEX}(A, b, c)
                                                                       feasible basic solution (if it exists)
    let \Delta be a new vector of length m
    while some index j \in N has c_i > 0
           choose an index e \in N for which c_e > 0
          for each index i \in B
                if a_{ie} > 0
                    \Delta_i = b_i/a_{ie}
                else \Delta_i = \infty
          choose an index l \in B that minimizes \Delta_i
10
          if \Delta_I == \infty
11
                return "unbounded"
          else (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, e)
     for i = 1 to n
          if i \in B
14
15
              \bar{x}_i = b_i
          else \bar{x}_i = 0
16
     return (\bar{x}_1, \bar{x}_2, \dots, \bar{x}_n)
```

Lemma 29.2

Suppose the call to INITIALIZE-SIMPLEX in line 1 returns a slack form for which the basic solution is feasible. Then if SIMPLEX returns a solution, it is a feasible solution. If SIMPLEX returns "unbounded", the linear program is unbounded.

```
SIMPLEX (A,b,c)

1 (N,B,A,b,c,\nu) = INITIALIZE-SIMPLEX (A,b,c)

2 \underbrace{\det \Delta \text{ be a new vector of length } m}_{}

3 \underbrace{\text{while some index } j \in N \text{ has } c_j > 0}_{}

4 \underbrace{\text{choose an index } e \in N \text{ for which } c_e > 0}_{}

5 \underbrace{\text{for each index } i \in B}_{}

6 \underbrace{\text{if } a_{ie} > 0}_{}

7 \underbrace{\Delta_i = b_i/a_{ie}}_{}

8 \underbrace{\text{else } \Delta_i = \infty}_{}

9 \underbrace{\text{choose an index } l \in B \text{ that minimizes } \Delta_i}_{}

10 \underbrace{\text{if } \Delta_l = \infty}_{}

11 \underbrace{\text{return "unbounded"}}_{}
```

Proof is based on the following three-part loop invariant:

Lemma 29 2 =

Suppose the call to INITIALIZE-SIMPLEX in line 1 returns a slack form for which the basic solution is feasible. Then if SIMPLEX returns a solution, it is a feasible solution. If SIMPLEX returns "unbounded", the linear program is unbounded.

```
SIMPLEX (A,b,c)

1 (N,B,A,b,c,\nu) = INITIALIZE-SIMPLEX (A,b,c)

2 |\det\Delta be a new vector of length m

3 while some index j \in N has c_j > 0

4 choose an index e \in N for which c_e > 0

5 for each index i \in B

6 if a_{ie} > 0

7 \Delta_i = b_i/a_{ie}

8 else \Delta_i = \infty

9 choose an index l \in B that minimizes \Delta_i

10 if \Delta_l = \infty

11 return "unbounded"
```

Proof is based on the following three-part loop invariant:

- 1. the slack form is always equivalent to the one returned by INITIALIZE-SIMPLEX,
- 2. for each $i \in B$, we have $b_i \ge 0$,
- 3. the basic solution associated with the (current) slack form is feasible.

Lemma 29 2 =

Suppose the call to INITIALIZE-SIMPLEX in line 1 returns a slack form for which the basic solution is feasible. Then if SIMPLEX returns a solution, it is a feasible solution. If SIMPLEX returns "unbounded", the linear program is unbounded.

Outline

Simplex Algorithm by Example

Details of the Simplex Algorithm

Finding an Initial Solution

Appendix: Cycling and Termination (non-examinable)

Finding an Initial Solution

maximise subject to

Finding an Initial Solution

maximise subject to

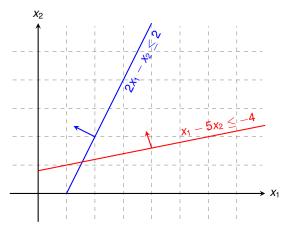
Finding an Initial Solution

maximise
$$2x_1 - x_2$$
 subject to
$$2x_1 - x_2 \leq 2 \\ x_1 - 5x_2 \leq -4 \\ x_1, x_2 \geq 0$$
 Conversion into slack form
$$z = 2x_1 - x_2 \\ x_3 = 2 - 2x_1 + x_2 \\ x_4 = -4 - x_1 + 5x_2$$
 Basic solution $(x_1, x_2, x_3, x_4) = (0, 0, 2, -4)$ is not feasible!

Geometric Illustration

maximise subject to

$$2x_1 - x_2$$

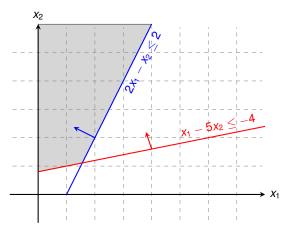


Geometric Illustration

maximise subject to

$$2x_1 - x_2$$

$$\begin{array}{ccccc} 2x_1 & - & x_2 & \leq & 2 \\ x_1 & - & 5x_2 & \leq & -4 \\ x_1, x_2 & \geq & 0 \end{array}$$



Geometric Illustration

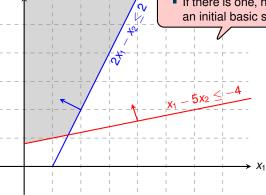
maximise subject to

$$2x_1 - x_2$$

 χ_2

Questions:

- How to determine whether there is any feasible solution?
- If there is one, how to determine an initial basic solution?



$$\sum_{j=1}^n c_j x_j$$

$$\begin{array}{ccc} \sum_{j=1}^n a_{ij} x_j & \leq & b_i & \text{ for } i=1,2,\ldots,m, \\ x_j & \geq & 0 & \text{ for } j=1,2,\ldots,n \end{array}$$

$$\sum_{j=1}^{n} c_j x_j$$

$$\begin{array}{cccc} \sum_{j=1}^n a_{ij} x_j & \leq & b_i & \text{for } i=1,2,\ldots,m, \\ x_j & \geq & 0 & \text{for } j=1,2,\ldots,n \\ & & \downarrow & \text{Formulating an Auxiliary Linear Program} \end{array}$$

maximise
$$\sum_{j=1}^{n} c_{j}x_{j}$$
 subject to
$$\sum_{j=1}^{n} a_{ij}x_{j} \leq b_{i} \quad \text{for } i=1,2,\ldots,m,$$

$$x_{j} \geq 0 \quad \text{for } j=1,2,\ldots,n$$
 Formulating an Auxiliary Linear Program maximise subject to
$$\sum_{j=1}^{n} a_{ij}x_{j} - x_{0} \leq b_{i} \quad \text{for } i=1,2,\ldots,m,$$

$$x_{i} \geq 0 \quad \text{for } j=0,1,\ldots,n$$

maximise
$$\sum_{j=1}^{n} c_{j}x_{j}$$
 subject to
$$\sum_{j=1}^{n} a_{ij}x_{j} \leq b_{i} \quad \text{for } i=1,2,\ldots,m,$$

$$x_{j} \geq 0 \quad \text{for } j=1,2,\ldots,n$$

$$\downarrow$$
 Formulating an Auxiliary Linear Program maximise subject to
$$\sum_{j=1}^{n} a_{ij}x_{j} - x_{0} \leq b_{i} \quad \text{for } i=1,2,\ldots,m,$$

$$x_{j} \geq 0 \quad \text{for } j=0,1,\ldots,n$$

Let L_{aux} be the auxiliary LP of a linear program L in standard form. Then L is feasible if and only if the optimal objective value of L_{aux} is 0.

Lemma 29.11

maximise
$$\sum_{j=1}^{n} c_{j}x_{j}$$
 subject to
$$\sum_{j=1}^{n} a_{ij}x_{j} \leq b_{i} \quad \text{for } i=1,2,\ldots,m,$$

$$x_{j} \geq 0 \quad \text{for } j=1,2,\ldots,n$$
 Formulating an Auxiliary Linear Program maximise subject to
$$\sum_{j=1}^{n} a_{ij}x_{j} - x_{0} \leq b_{i} \quad \text{for } i=1,2,\ldots,m,$$

$$x_{i} \geq 0 \quad \text{for } j=0,1,\ldots,n$$

Let L_{aux} be the auxiliary LP of a linear program L in standard form. Then L is feasible if and only if the optimal objective value of L_{aux} is 0.

Proof. Exercise!

Lemma 29.11

■ Let us illustrate the role of x_0 as "distance from feasibility"	

- Let us illustrate the role of x_0 as "distance from feasibility"
- We'll also see that increasing x_0 enlarges the feasible region

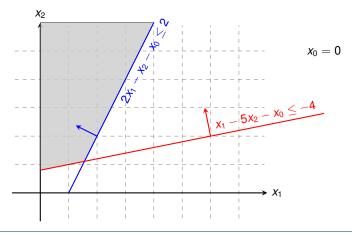
$$-x_0$$

$$2x_1 - x_2 - x_0$$

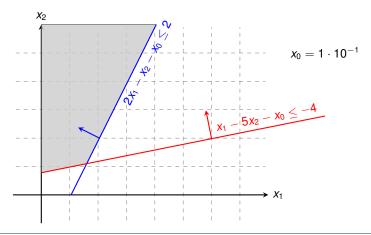
 $x_1 - 5x_2 - x_0$

$$x_1 - x_2 - x_0 \le 2$$

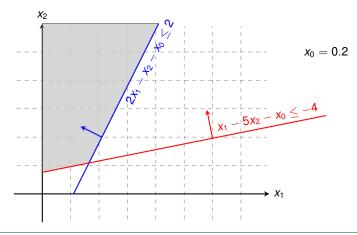
 $x_1 - 5x_2 - x_0 \le -4$
 $x_0, x_1, x_2 \ge 0$



$$-x_0$$

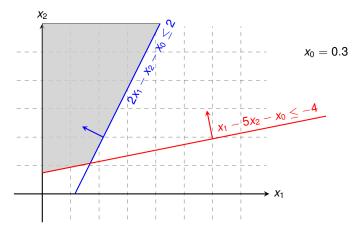


$$-x_0$$



$$-x_0$$

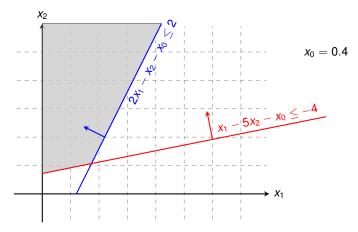
 $2x_{1}$



$$-x_0$$

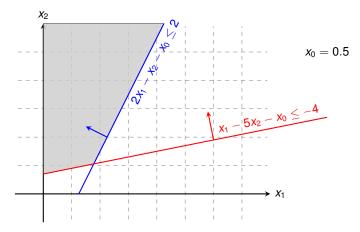
 $2x_{1}$

$$\begin{cases} 2x_1 - x_2 - x_0 \le 2 \\ x_1 - 5x_2 - x_0 \le -4 \\ x_0, x_1, x_2 \ge 0 \end{cases}$$

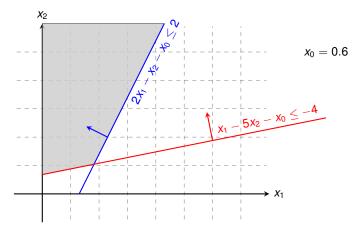


$$-x_0$$

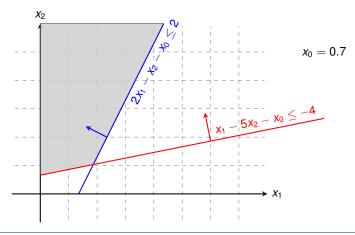
$$x_0 \leq$$



$$-x_0$$



$$-x_0$$



$$-x_0$$

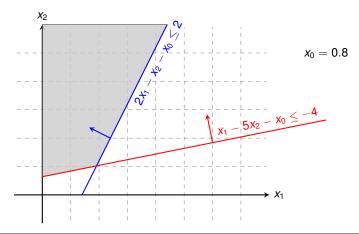
 $2x_{1}$ *X*₁

$$- x_2 - x_0$$

$$\leq$$

$$x_1 - x_2 - x_0 \le 2$$

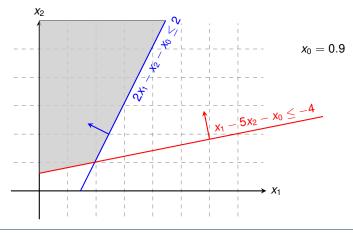
 $x_1 - 5x_2 - x_0 \le -4$
 $x_0, x_1, x_2 \ge 0$



$$-x_0$$

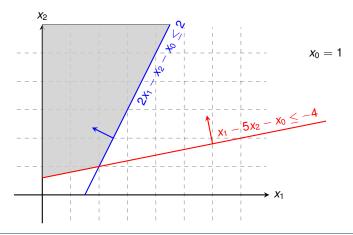
$$x_1 - x_2 - x_0$$

 $x_1 - 5x_2 - x_0$



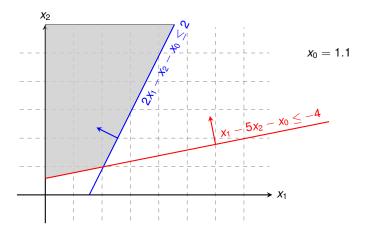
$$-x_0$$

 $2x_{1}$ *X*₁

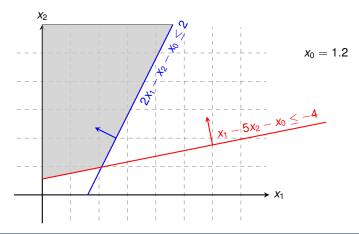


$$-x_0$$

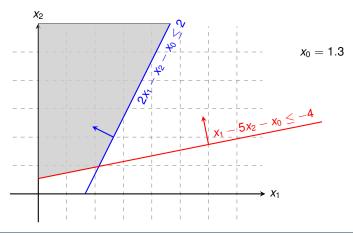
 $2x_{1}$ *X*₁



$$-x_0$$

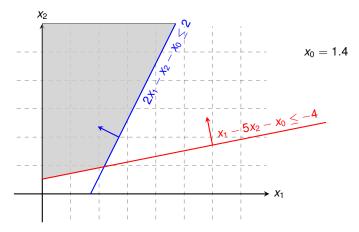


$$-x_0$$

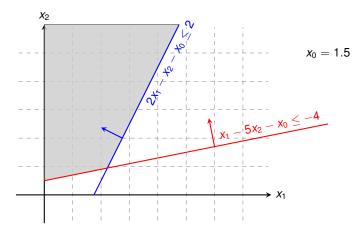


$$-x_0$$

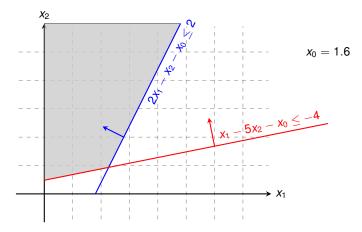
 $2x_{1}$ *X*₁



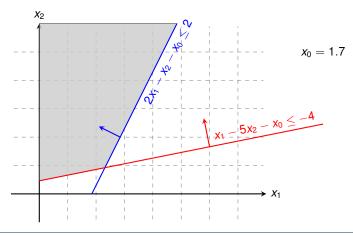
$$-x_0$$



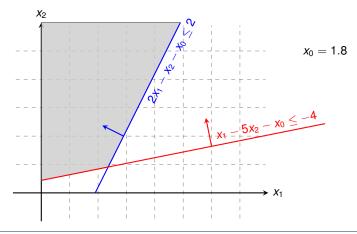
$$-x_0$$



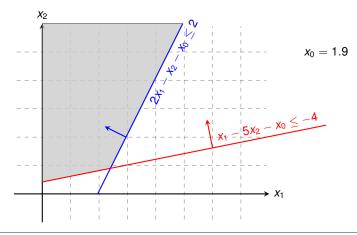
$$-x_0$$



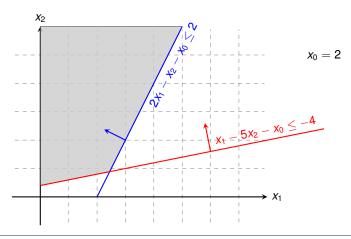
$$-x_0$$



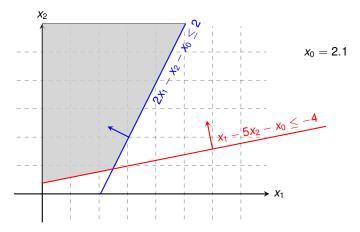
$$-x_0$$



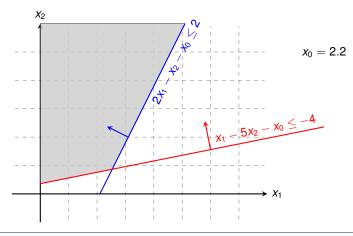
$$-x_0$$



$$-x_0$$

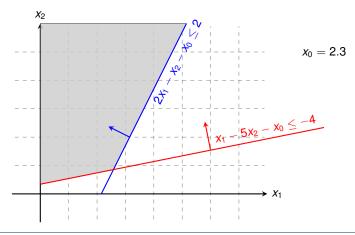


$$-x_0$$

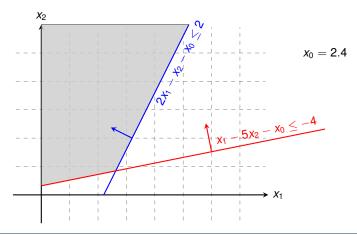


$$-x_0$$

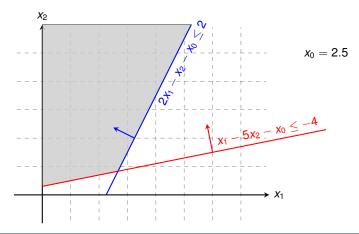
 $2x_{1}$ *X*₁



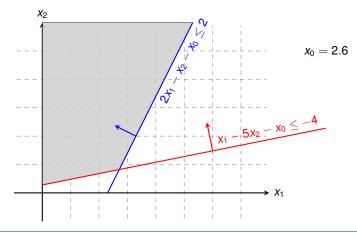
$$-x_0$$



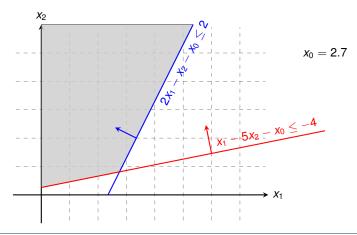
$$-x_0$$



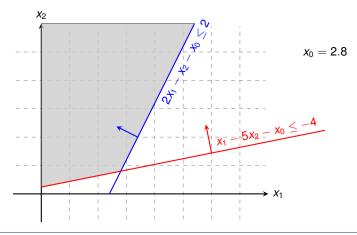
$$-x_0$$



$$-x_0$$

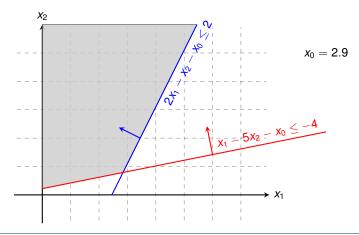


$$-x_0$$

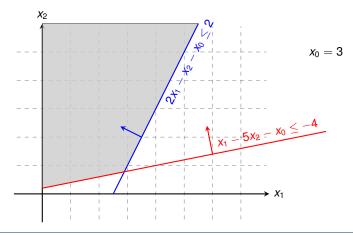


$$-x_0$$

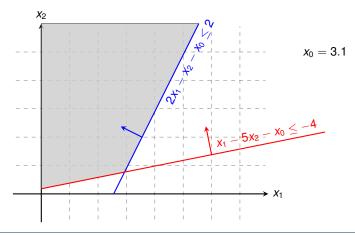
$$\leq$$



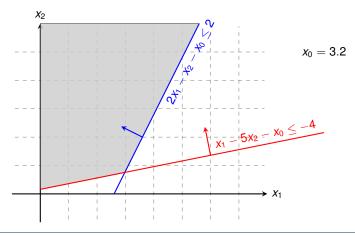
$$-x_0$$



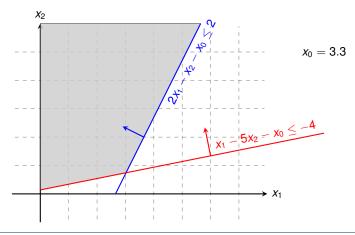
$$-x_0$$



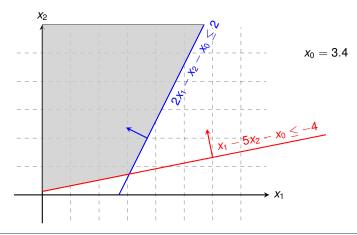
$$-x_0$$



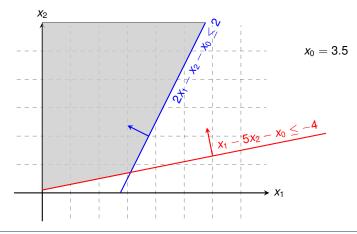
$$-x_0$$



$$-x_0$$

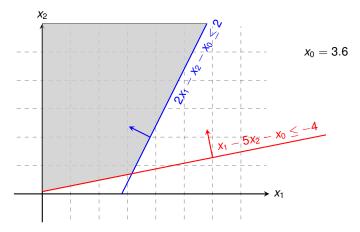


$$-x_0$$

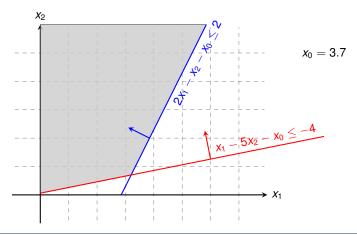


$$-x_0$$

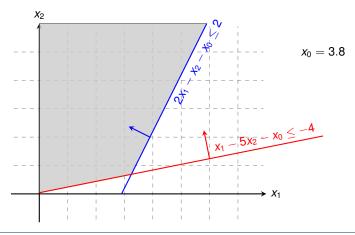
 X_0, X_1, X_2



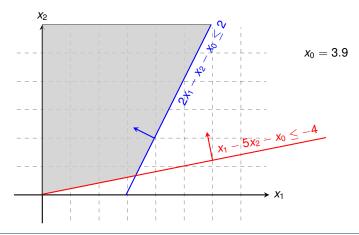
$$-x_0$$



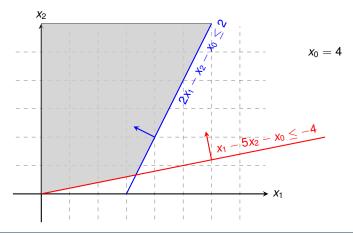
$$-x_0$$



$$-x_0$$



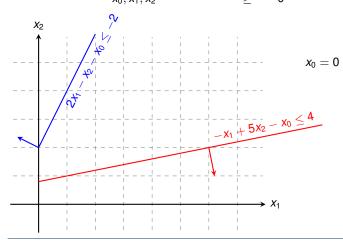
$$-x_0$$



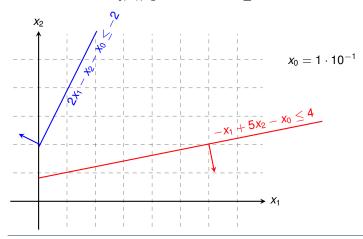
Let us now feasible	modify the ori	ginal linear pr	ogram so that	it is not

- Let us now modify the original linear program so that it is not feasible
- \Rightarrow Hence the auxiliary linear program has only a solution for a sufficiently large $x_0 > 0$!

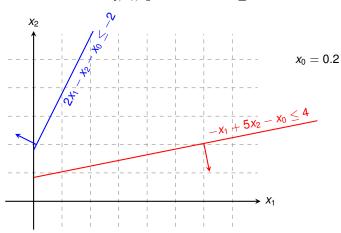
$$-x_0$$



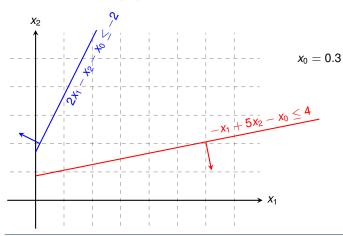
$$-x_0$$



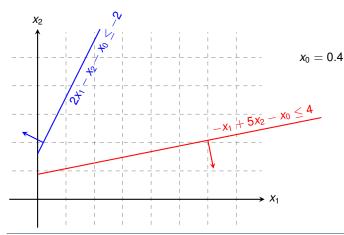
$$-x_0$$



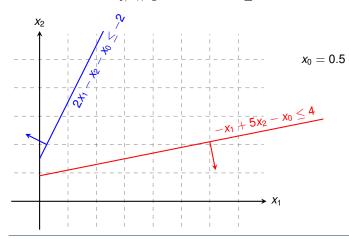
$$-x_0$$



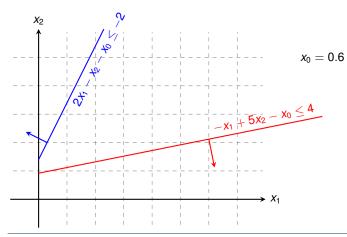
$$-x_0$$



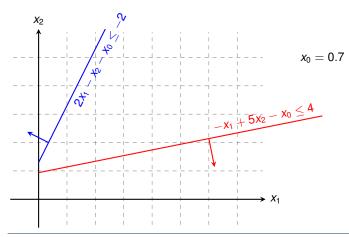
$$-x_0$$



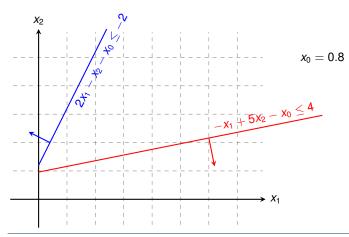
$$-x_0$$



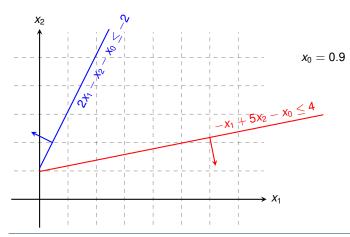
$$-x_0$$



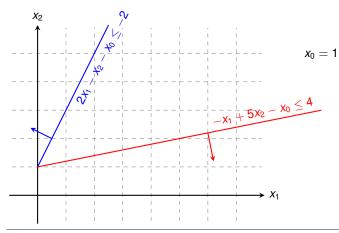
$$-x_0$$



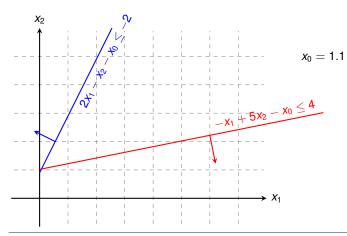
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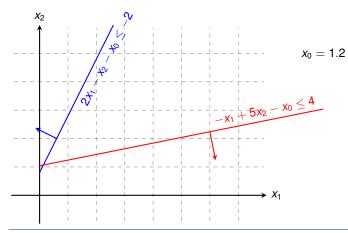
$$-x_0$$



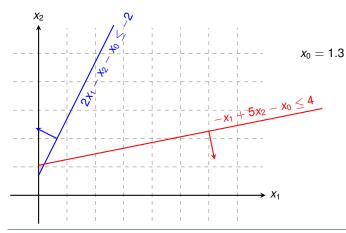
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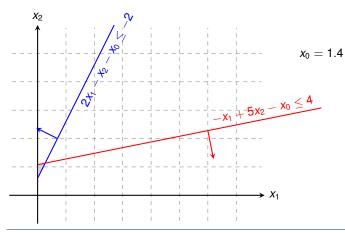
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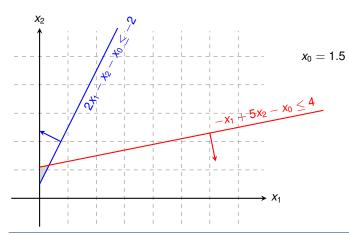
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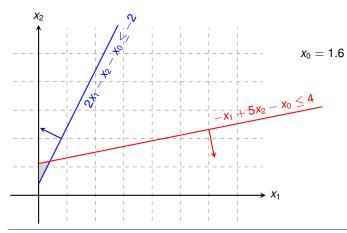
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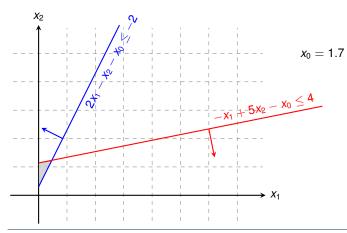
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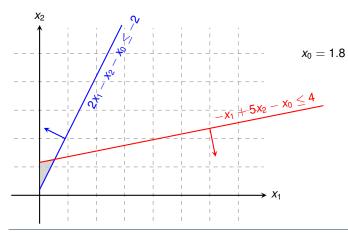
$$-x_0$$



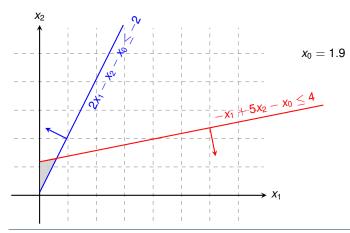
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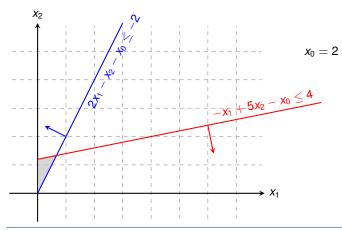
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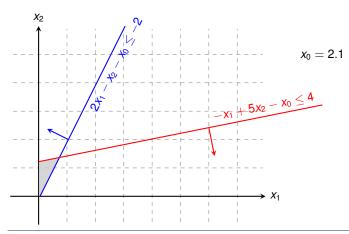
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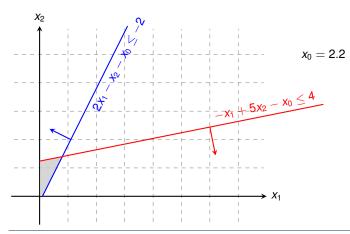
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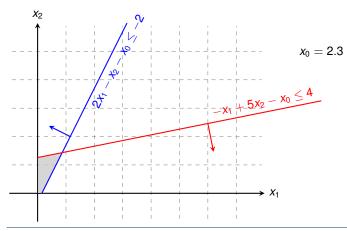
$$-x_0$$



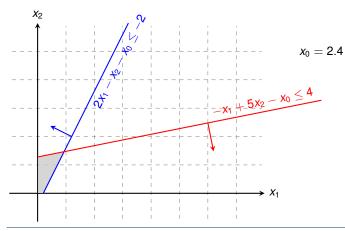
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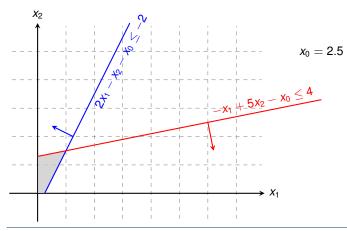
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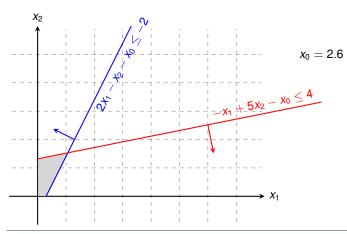
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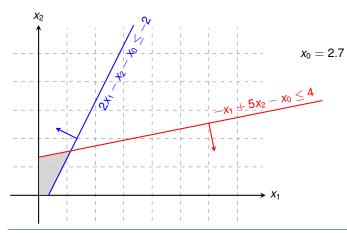
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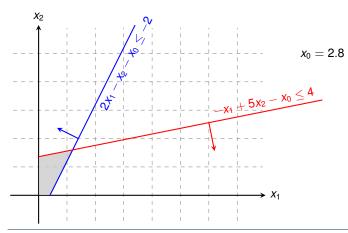
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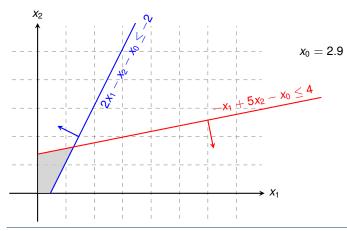
$$-x_0$$



$$-x_0$$

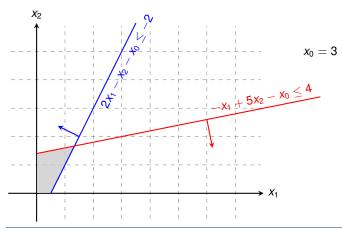


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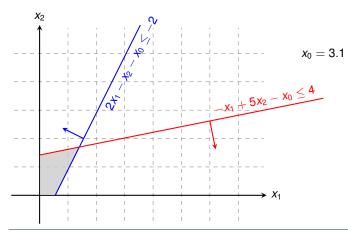


maximise subject to

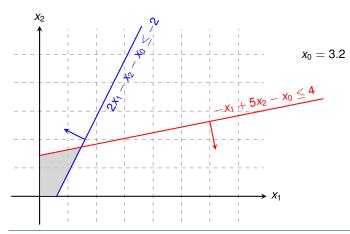
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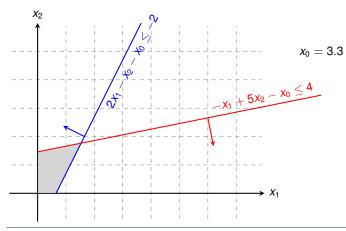
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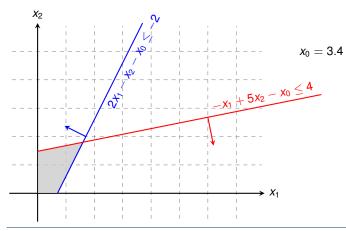
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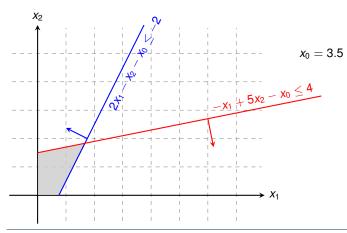
$$-x_0$$



$$-x_0$$

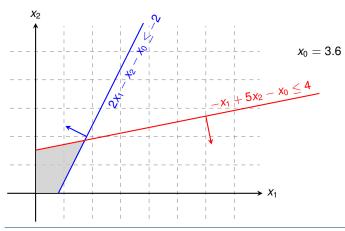


$$-x_0$$

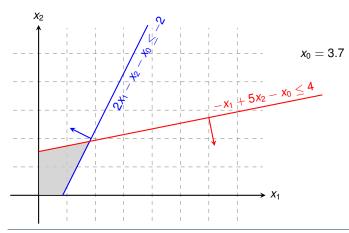


maximise subject to

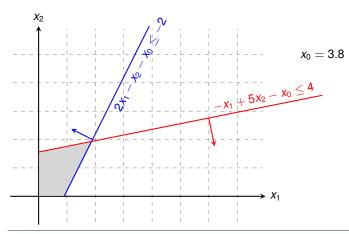
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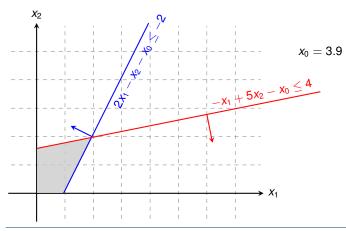
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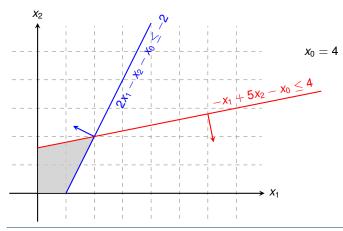
$$-x_0$$



$$-x_0$$



$$-x_0$$



```
INITIALIZE-SIMPLEX (A, b, c)
     let k be the index of the minimum b_i
                                   // is the initial basic solution feasible?
 2 if b_{\nu} > 0
          return (\{1, 2, ..., n\}, \{n + 1, n + 2, ..., n + m\}, A, b, c, 0)
    form L_{\text{aux}} by adding -x_0 to the left-hand side of each constraint
          and setting the objective function to -x_0
 5 let (N, B, A, b, c, \nu) be the resulting slack form for L_{aux}
   l = n + k
     //L_{\text{any}} has n+1 nonbasic variables and m basic variables.
 8 (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, 0)
 9 // The basic solution is now feasible for L_{\text{aux}}.
10 iterate the while loop of lines 3-12 of SIMPLEX until an optimal solution
          to L_{\text{aux}} is found
     if the optimal solution to L_{\text{aux}} sets \bar{x}_0 to 0
12
          if \bar{x}_0 is basic
13
               perform one (degenerate) pivot to make it nonbasic
14
          from the final slack form of L_{\text{aux}}, remove x_0 from the constraints and
               restore the original objective function of L, but replace each basic
               variable in this objective function by the right-hand side of its
               associated constraint
15
          return the modified final slack form
```

else return "infeasible"

INITIALIZE-SIMPLEX (A, b, c)

to L_{aux} is found

if \bar{x}_0 is basic

else return "infeasible"

2 **if** $b_{\nu} > 0$

l = n + k

12

13

14

15

```
Test solution with N = \{1, 2, \dots, n\}, B = \{n + 1, n + 1\}
                                                   \{2,\ldots,n+m\},\ \overline{x}_i=b_i\ \text{for}\ i\in B,\ \overline{x}_i=0\ \text{otherwise}.
     let k be the index of the minimum b_k
                                   // is the initial basic solution feasible?
          return (\{1, 2, ..., n\}, \{n + 1, n + 2, ..., n + m\}, A, b, c, 0)
    form L_{\text{aux}} by adding -x_0 to the left-hand side of each constraint
          and setting the objective function to -x_0
   let (N, B, A, b, c, \nu) be the resulting slack form for L_{aux}
     //L_{\text{any}} has n+1 nonbasic variables and m basic variables.
8 (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, 0)
    // The basic solution is now feasible for L_{\text{aux}}.
10 iterate the while loop of lines 3-12 of SIMPLEX until an optimal solution
     if the optimal solution to L_{\text{aux}} sets \bar{x}_0 to 0
               perform one (degenerate) pivot to make it nonbasic
          from the final slack form of L_{\text{aux}}, remove x_0 from the constraints and
               restore the original objective function of L, but replace each basic
               variable in this objective function by the right-hand side of its
          return the modified final slack form
```

associated constraint

Test solution with $N = \{1, 2, ..., n\}$, $B = \{n + 1, n + 2, ..., n + m\}$, $\overline{x}_i = b_i$ for $i \in B$, $\overline{x}_i = 0$ otherwise.

ℓ will be the leaving variable so

that x_{ℓ} has the most negative value.

Initialize-Simplex (A, b, c)

- 1 let k be the index of the minimum b_i
- 2 if $b_k \ge 0$ // is the initial basic solution feasible?
- 3 **return** $(\{1, 2, ..., n\}, \{n+1, n+2, ..., n+m\}, A, b, c, 0)$
- 4 form L_{aux} by adding $-x_0$ to the left-hand side of each constraint
- and setting the objective function to $-x_0$ 5 let (N, B, A, b, c, v) be the resulting slack form for L_{aux}
- l = n + k
- 7 // L_{aux} has n+1 nonbasic variables and m basic variables.
- 8 (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, 0)
- 9 // The basic solution is now feasible for L_{max} .
- 10 iterate the **while** loop of lines 3–12 of SIMPLEX until an optimal solution to $L_{\rm aux}$ is found
- 11 **if** the optimal solution to L_{aux} sets \bar{x}_0 to 0
- 12 **if** \bar{x}_0 is basic
- 13 perform one (degenerate) pivot to make it nonbasic
- from the final slack form of L_{mix}, remove x₀ from the constraints and restore the original objective function of L, but replace each basic variable in this objective function by the right-hand side of its associated constraint
- 15 return the modified final slack form
- 16 else return "infeasible"

```
Test solution with N = \{1, 2, \dots, n\}, B = \{n + 1, n + 1\}
INITIALIZE-SIMPLEX (A, b, c)
                                                   \{2,\ldots,n+m\},\ \overline{x}_i=b_i\ \text{for}\ i\in B,\ \overline{x}_i=0\ \text{otherwise}.
     let k be the index of the minimum b_k
                                   // is the initial basic solution feasible?
 2 if b_{\nu} > 0
          return (\{1, 2, ..., n\}, \{n + 1, n + 2, ..., n + m\}, A, b, c, 0)
     form L_{\text{aux}} by adding -x_0 to the left-hand side of each constraint
          and setting the objective function to -x_0
                                                                                ℓ will be the leaving variable so
    let (N, B, A, b, c, \nu) be the resulting slack form for L_{aux}
    l = n + k
                                                                            that x_{\ell} has the most negative value.
     //L_{aux} has n+1 nonbasic variables and m basic variables.
   (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, 0)
                                                                Pivot step with x_{\ell} leaving and x_0 entering.
     // The basic solution is now feasible for L_{\text{aux}}.
    iterate the while loop of lines 3-12 of SIMPLEX until an optimal solution
          to L_{\text{aux}} is found
     if the optimal solution to L_{\text{aux}} sets \bar{x}_0 to 0
12
          if \bar{x}_0 is basic
13
               perform one (degenerate) pivot to make it nonbasic
14
          from the final slack form of L_{\text{aux}}, remove x_0 from the constraints and
               restore the original objective function of L, but replace each basic
               variable in this objective function by the right-hand side of its
               associated constraint
15
          return the modified final slack form
     else return "infeasible"
```

```
Test solution with N = \{1, 2, \dots, n\}, B = \{n + 1, n + 1\}
INITIALIZE-SIMPLEX (A, b, c)
                                                  \{2,\ldots,n+m\},\ \overline{x}_i=b_i\ \text{for}\ i\in B,\ \overline{x}_i=0\ \text{otherwise}.
     let k be the index of the minimum b_k
                                  // is the initial basic solution feasible?
 2 if b_{\nu} > 0
          return (\{1, 2, ..., n\}, \{n + 1, n + 2, ..., n + m\}, A, b, c, 0)
     form L_{\text{aux}} by adding -x_0 to the left-hand side of each constraint
          and setting the objective function to -x_0
                                                                               ℓ will be the leaving variable so
    let (N, B, A, b, c, \nu) be the resulting slack form for L_{aux}
    l = n + k
                                                                           that x_{\ell} has the most negative value.
     //L_{aux} has n+1 nonbasic variables and m basic variables.
   (N, B, A, b, c, v) = PIVOT(N, B, A, b, c, v, l, 0)
                                                               Pivot step with x_{\ell} leaving and x_0 entering.
     // The basic solution is now feasible for L_{\text{aux}}.
    iterate the while loop of lines 3-12 of SIMPLEX until an optimal solution
          to L_{\text{aux}} is found
                                                                             This pivot step does not change
     if the optimal solution to L_{\text{aux}} sets \bar{x}_0 to 0
12
          if \bar{x}_0 is basic
                                                                                the value of any variable.
               perform one (degenerate) pivot to make it nonbasic
13
14
          from the final slack form of L_{\text{aux}}, remove x_0 from the constraints and
               restore the original objective function of L, but replace each basic
               variable in this objective function by the right-hand side of its
               associated constraint
15
          return the modified final slack form
     else return "infeasible"
```

maximise
$$2x_1 - x_2$$
 subject to $2x_1 - x_2 \le 2$ $x_1 - 5x_2 \le -4$ $x_1, x_2 \ge 0$

maximise subject to

$$2x_1 - x_2$$

Formulating the auxiliary linear program (as basic solution would not be feasible!)

maximise subject to

$$2x_1 - x_2$$

Formulating the auxiliary linear program (as basic solution would not be feasible!)

maximise subject to

$$x_0$$

$$2x_1 - x_2$$

$$\begin{array}{ccccc} 2x_1 & - & x_2 & \leq & 2 \\ x_1 & - & 5x_2 & \leq & -4 \\ & x_1, x_2 & \geq & 0 \end{array}$$

Formulating the auxiliary linear program (as basic solution would not be feasible!)

maximise subject to

$$X_0$$

Converting into slack form

maximise subject to
$$2x_1 - x_2 \leq 2$$

$$x_1 - 5x_2 \leq -4$$

$$x_1, x_2 \geq 0$$
Formulating the auxiliary linear program (as basic solution would not be feasible!)

maximise subject to
$$2x_1 - x_2 - x_0$$

$$2x_1 - x_2 - x_0 \leq 2$$

$$x_1 - 5x_2 - x_0 \leq -4$$

$$x_1, x_2, x_0 \geq 0$$
Converting into slack form
$$z = x_3 = 2 - 2x_1 + x_2 + x_0$$

 X_4

5*x*₂

 X_0

maximise subject to
$$2x_1 - x_2 \leq 2$$

$$x_1 - 5x_2 \leq -4$$

$$x_1, x_2 \geq 0$$
Formulating the auxiliary linear program (as basic solution would not be feasible!)

maximise
$$-x_0$$
subject to
$$2x_1 - x_2 - x_0 \leq 2$$

$$x_1 - 5x_2 - x_0 \leq 2$$

$$x_1 - 5x_2 - x_0 \leq -4$$

$$x_1, x_2, x_0 \geq 0$$
Basic solution (0,0,0,2,-4) not feasible!
$$Converting into slack form$$

$$z = x_3 = 2 - 2x_1 + x_2 + x_0$$

$$x_4 = -4 - x_1 + 5x_2 + x_0$$

$$z = x_3 = 2 - 2x_1 + x_2 + x_0$$

 $x_4 = -4 - x_1 + 5x_2 + x_0$
Pivot with x_0 entering and x_4 leaving

Basic solution (4, 0, 0, 6, 0) is feasible!

Example of Initialize-SIMPLEX (2/3)

$$\begin{array}{rclcrcr}
z & = & - & x_0 \\
x_2 & = & \frac{4}{5} & - & \frac{x_0}{5} & + & \frac{x_1}{5} & + & \frac{x_5}{5} \\
x_3 & = & \frac{14}{5} & + & \frac{4x_0}{5} & - & \frac{9x_1}{5} & + & \frac{x_5}{5}
\end{array}$$

$$z = -x_0$$

$$x_2 = \frac{4}{5} - \frac{x_0}{5} + \frac{x_1}{5} + \frac{x_4}{5}$$

$$x_3 = \frac{14}{5} + \frac{4x_0}{5} - \frac{9x_1}{5} + \frac{x_4}{5}$$

$$\text{Set } x_0 = 0 \text{ and express objective function}$$

$$\text{by non-basic variables}$$

$$\begin{array}{rclcrcr}
z & = & - & x_0 \\
x_2 & = & \frac{4}{5} & - & \frac{x_0}{5} & + & \frac{x_1}{5} & + & \frac{x_2}{5} \\
x_3 & = & \frac{14}{5} & + & \frac{4x_0}{5} & - & \frac{9x_1}{5} & + & \frac{x_2}{5}
\end{array}$$

$$2x_1 - x_2 = 2x_1 - \left(\frac{4}{5} - \frac{x_0}{5} + \frac{x_1}{5} + \frac{x_4}{5}\right)$$

Set $x_0 = 0$ and express objective function by non-basic variables

$$z = -\frac{4}{5} + \frac{9x_1}{5} - \frac{x_2}{5}$$

$$x_2 = \frac{4}{5} + \frac{x_1}{5} + \frac{x_2}{5}$$

$$x_3 = \frac{14}{5} - \frac{9x_1}{5} + \frac{x_2}{5}$$

$$\begin{array}{rclcrcr}
z & = & - & x_0 \\
x_2 & = & \frac{4}{5} & - & \frac{x_0}{5} & + & \frac{x_1}{5} & + & \frac{x_2}{5} \\
x_3 & = & \frac{14}{5} & + & \frac{4x_0}{5} & - & \frac{9x_1}{5} & + & \frac{x_2}{5}
\end{array}$$

$$2x_1 - x_2 = 2x_1 - \left(\frac{4}{5} - \frac{x_0}{5} + \frac{x_1}{5} + \frac{x_4}{5}\right)$$

Set $x_0 = 0$ and express objective function by non-basic variables

$$z = -\frac{4}{5} + \frac{9x_1}{5} - \frac{x_4}{5}$$

$$x_2 = \frac{4}{5} + \frac{x_1}{5} + \frac{x_4}{5}$$

$$x_3 = \frac{14}{5} - \frac{9x_1}{5} + \frac{x_4}{5}$$

Basic solution $(0, \frac{4}{5}, \frac{14}{5}, 0)$, which is feasible!

Example of Initialize-Simplex (3/3)

$$\begin{array}{rclcrcr}
z & = & - & x_0 \\
x_2 & = & \frac{4}{5} & - & \frac{x_0}{5} & + & \frac{x_1}{5} & + & \frac{x_2}{5} \\
x_3 & = & \frac{14}{5} & + & \frac{4x_0}{5} & - & \frac{9x_1}{5} & + & \frac{x_2}{5}
\end{array}$$

$$2x_1 - x_2 = 2x_1 - \left(\frac{4}{5} - \frac{x_0}{5} + \frac{x_1}{5} + \frac{x_4}{5}\right)$$

Set $x_0 = 0$ and express objective function by non-basic variables

$$z = -\frac{4}{5} + \frac{9x_1}{5} - \frac{x_4}{5}$$

$$x_2 = \frac{4}{5} + \frac{x_1}{5} + \frac{x_4}{5}$$

$$x_3 = \frac{14}{5} - \frac{9x_1}{5} + \frac{x_4}{5}$$

Basic solution $(0, \frac{4}{5}, \frac{14}{5}, 0)$, which is feasible!

Lemma 29.12

If a linear program L has no feasible solution, then INITIALIZE-SIMPLEX returns "infeasible". Otherwise, it returns a valid slack form for which the basic solution is feasible.

Fundamental Theorem of Linear Programming

Theorem 29.13 (Fundamental Theorem of Linear Programming)

For any linear program *L*, given in standard form, either:

- 1. L is infeasible \Rightarrow SIMPLEX returns "infeasible".
- 2. *L* is unbounded \Rightarrow SIMPLEX returns "unbounded".
- 3. L has an optimal solution with a finite objective value
 - ⇒ SIMPLEX returns an optimal solution with a finite objective value.

Small Technicality: need to equip SIMPLEX with an "anti-cycling strategy" (see extra slides)

Fundamental Theorem of Linear Programming

Theorem 29.13 (Fundamental Theorem of Linear Programming)

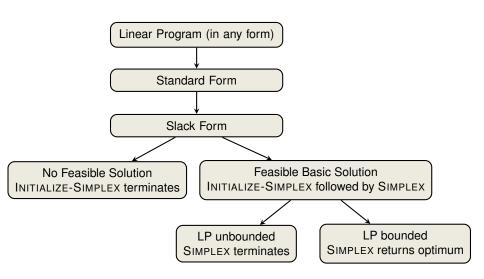
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- 3. L has an optimal solution with a finite objective value
 - ⇒ SIMPLEX returns an optimal solution with a finite objective value.

Small Technicality: need to equip SIMPLEX with an "anti-cycling strategy" (see extra slides)

Proof requires the concept of duality, which is not covered in this course (for details see CLRS3, Chapter 29.4)

Workflow for Solving Linear Programs



Linear Programming and Simplex: Summary and Outlook Linear Programming

extremely versatile tool for modelling problems of all kinds

Linear Programming ———

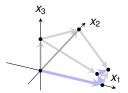
- extremely versatile tool for modelling problems of all kinds
- basis of Integer Programming, to be discussed in later lectures

Linear Programming ——

- extremely versatile tool for modelling problems of all kinds
- basis of Integer Programming, to be discussed in later lectures

Simplex Algorithm

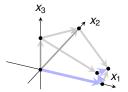
• In practice: usually terminates in polynomial time, i.e., O(m+n)



- Linear Programming _
- extremely versatile tool for modelling problems of all kinds
- basis of Integer Programming, to be discussed in later lectures

Simplex Algorithm

- In practice: usually terminates in polynomial time, i.e., O(m + n)
- In theory: even with anti-cycling may need exponential time



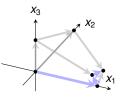
Linear Programming

- extremely versatile tool for modelling problems of all kinds
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Simplex Algorithm

- In practice: usually terminates in polynomial time, i.e., O(m+n)
- In theory: even with anti-cycling may need exponential time

Research Problem: Is there a pivoting rule which makes SIMPLEX a polynomial-time algorithm?

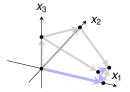


Linear Programming -

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Simplex Algorithm

- In practice: usually terminates in polynomial time, i.e., O(m+n)
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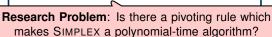
Polynomial-Time Algorithms -

Linear Programming -

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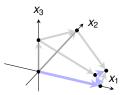
Simplex Algorithm

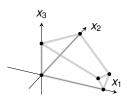
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- In theory: even with anti-cycling may need exponential time



Polynomial-Time Algorithms -

 Interior-Point Methods: traverses the interior of the feasible set of solutions (not just vertices!)



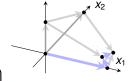


Linear Programming -

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Simplex Algorithm

- In practice: usually terminates in polynomial time, i.e., O(m+n)
- In theory: even with anti-cycling may need exponential time

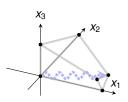


Xз

Research Problem: Is there a pivoting rule which makes SIMPLEX a polynomial-time algorithm?

Polynomial-Time Algorithms -

 Interior-Point Methods: traverses the interior of the feasible set of solutions (not just vertices!)



Outlook: Alternatives to Worst Case Analysis (non-examinable)

1.2 Famous Failures and the Need for Alternatives

For many problems a bit beyond the scope of an undergraduate course, the downside of worst-case analysis rears its ugly head. This section reviews four famous examples in which worst-case analysis gives misleading or useless advice about how to solve a problem. These examples motivate the alternatives to worst-case analysis that are surveyed in Section 1.4 and described in detail in later chapters of the book.

1.2.1 The Simplex Method for Linear Programming

Perhaps the <u>most famous failure of worst-case analysis concerns linear programming</u>, the problem of optimizing a linear function subject to linear constraints (Figure 1.1). Dantzig proposed in the 1940s an algorithm for solving linear programs called the *simplex method*. The simplex method solves linear programs using greedy local

Source: "Beyond the Worst-Case Analysis of Algorithms" by Tim Roughgarden, 2020

Outline

Simplex Algorithm by Example

Details of the Simplex Algorithm

Finding an Initial Solution

Appendix: Cycling and Termination (non-examinable)

$$z = x_1 + x_2 + x_3$$

 $x_4 = 8 - x_1 - x_2$
 $x_5 = x_2 - x_3$

$$z$$
 = x_1 + x_2 + x_3
 x_4 = 8 - x_1 - x_2
 x_5 = x_2 - x_3
Pivot with x_1 entering and x_4 leaving

$$z = x_1 + x_2 + x_3$$

 $x_4 = 8 - x_1 - x_2$
 $x_5 = x_2 - x_3$
 $\begin{vmatrix} \text{Pivot with } x_1 \text{ entering and } x_4 \text{ leaving} \end{vmatrix}$
 $z = 8 + x_3 - x_4$
 $x_1 = 8 - x_2 - x_3$
 $\begin{vmatrix} \text{Pivot with } x_3 \text{ entering and } x_5 \text{ leaving} \end{vmatrix}$



Exercise: Execute one more step of the Simplex Algorithm on the tableau from the previous slide.

Cycling: SIMPLEX may fail to terminate.

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

Anti-Cycling Strategies

1. Bland's rule: Choose entering variable with smallest index

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

Anti-Cycling Strategies

- 1. Bland's rule: Choose entering variable with smallest index
- 2. Random rule: Choose entering variable uniformly at random

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

Anti-Cycling Strategies

- Bland's rule: Choose entering variable with smallest index
- 2. Random rule: Choose entering variable uniformly at random
- 3. Perturbation: Perturb the input slightly so that it is impossible to have two solutions with the same objective value

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

Anti-Cycling Strategies

- Bland's rule: Choose entering variable with smallest index
- 2. Random rule: Choose entering variable uniformly at random
- 3. Perturbation: Perturb the input slightly so that it is impossible to have two solutions with the same objective value

Replace each b_i by $\hat{b}_i = b_i + \epsilon_i$, where $\epsilon_i \gg \epsilon_{i+1}$ are all small.

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

Anti-Cycling Strategies

- 1. Bland's rule: Choose entering variable with smallest index
- 2. Random rule: Choose entering variable uniformly at random
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Lemma 29.7

Assuming Initialize-SIMPLEX returns a slack form for which the basic solution is feasible, SIMPLEX either reports that the program is unbounded or returns a feasible solution in at most $\binom{n+m}{m}$ iterations.

It is theoretically possible, but very rare in practice.

Cycling: SIMPLEX may fail to terminate.

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Lemma 29.7

Assuming Initialize-SIMPLEX returns a slack form for which the basic solution is feasible, SIMPLEX either reports that the program is unbounded or returns a feasible solution in at most $\binom{n+m}{m}$ iterations.

Every set *B* of basic variables uniquely determines a slack form, and there are at most $\binom{n+m}{m}$ unique slack forms.

Randomised Algorithms

Lecture 8: Solving a TSP Instance using Linear Programming

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Introduction

Examples of TSP Instances

Demonstration

The Traveling Salesman Problem (TSP)

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

The Traveling Salesman Problem (TSP)

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

Formal Definition			

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

Formal Definition

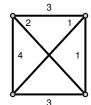
■ Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

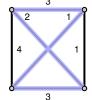
Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.



Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

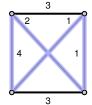
- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of *G* with minimum cost.



$$3+2+1+3=9$$

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
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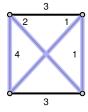
$$2+4+1+1=8$$

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

Formal Definition

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

Solution space consists of at most n! possible tours!



$$2+4+1+1=8$$

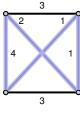
Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

Formal Definition

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

Solution space consists of at most n! possible tours!

Actually the right number is (n-1)!/2



$$2+4+1+1=8$$

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

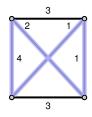
Formal Definition

Special Instances

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

Solution space consists of at most n! possible tours!

Actually the right number is (n-1)!/2



$$2+4+1+1=8$$

8. Solving TSP via Linear Programming © T. Sauerwald

Given a set of cities along with the cost of travel between them, find the cheapest route visiting all cities and returning to your starting point.

Formal Definition

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

Solution space consists of at most n! possible tours!

 $\frac{4}{3}$ 2+4+1+1=8

3

Actually the right number is
$$(n-1)!/2$$

Special Instances

Metric TSP: costs satisfy triangle inequality:

$$\forall u, v, w \in V$$
: $c(u, w) \leq c(u, v) + c(v, w)$.

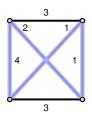
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Formal Definition

- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

Solution space consists of at most *n*! possible tours!

Actually the right number is (n-1)!/2



$$2+4+1+1=8$$

Special Instances

■ Metric TSP: costs satisfy triangle inequality: < NP hard (Ex. 35.2-2)

Even this version is

$$\forall u, v, w \in V$$
: $c(u, w) \leq c(u, v) + c(v, w)$.

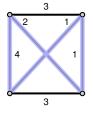
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- Given: A complete undirected graph G = (V, E) with nonnegative integer cost c(u, v) for each edge $(u, v) \in E$
- Goal: Find a hamiltonian cycle of G with minimum cost.

Solution space consists of at most n! possible tours!

Actually the right number is (n-1)!/2



$$2+4+1+1=8$$

Special Instances

 Metric TSP: costs satisfy triangle inequality:
 Even this version is NP hard (Ex. 35.2-2)

$$\forall u, v, w \in V$$
: $c(u, w) \leq c(u, v) + c(v, w)$.

 Euclidean TSP: cities are points in the Euclidean space, costs are equal to their (rounded) Euclidean distance

Outline

Introduction

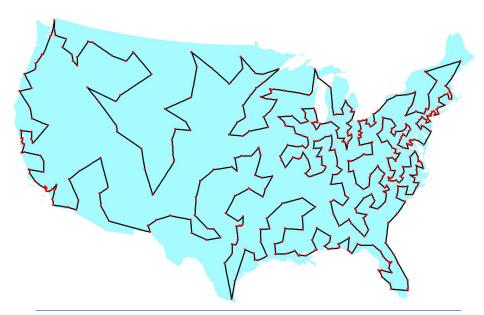
Examples of TSP Instances

Demonstration

33 city contest (1964)



532 cities (1987 [Padberg, Rinaldi])



13,509 cities (1999 [Applegate, Bixby, Chavatal, Cook])



SOLUTION OF A LARGE-SCALE TRAVELING-SALESMAN PROBLEM*

G. DANTZIG, R. FULKERSON, AND S. JOHNSON

The Rand Corporation, Santa Monica, California

(Received August 9, 1954)

It is shown that a certain tour of 49 cities, one in each of the 48 states and Washington, D. C., has the shortest road distance.

THE TRAVELING-SALESMAN PROBLEM might be described as I follows: Find the shortest route (tour) for a salesman starting from a given city, visiting each of a specified group of cities, and then returning to the original point of departure. More generally, given an n by n symmetric matrix $D = (d_{IJ})$, where d_{IJ} represents the 'distance' from I to J, arrange the points in a cyclic order in such a way that the sum of the d_{IJ} between consecutive points is minimal. Since there are only a finite number of possibilities (at most $\frac{1}{2}(n-1)!$) to consider, the problem is to devise a method of picking out the optimal arrangement which is reasonably efficient for fairly large values of n. Although algorithms have been devised for problems of similar nature, e.g., the optimal assignment problem, 3,7,8 little is known about the traveling-salesman problem. We do not claim that this note alters the situation very much; what we shall do is outline a way of approaching the problem that sometimes, at least, enables one to find an optimal path and prove it so. In particular, it will be shown that a certain arrangement of 49 cities, one in each of the 48 states and Washington, D. C., is best, the d_{IJ} used representing road distances as taken from an atlas.

The 42 (49) Cities

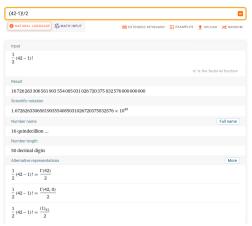
- 1. Manchester, N. H.
- 2. Montpelier, Vt.
- 3. Detroit, Mich.
- 4. Cleveland, Ohio 5. Charleston, W. Va.
- 6. Louisville, Kv.
- 7. Indianapolis, Ind.
- 8. Chicago, Ill.
- Milwaukee, Wis.
- 10. Minneapolis, Minn.
- 11. Pierre, S. D.
- 12. Bismarck, N. D.
- 13. Helena, Mont.
- 14. Seattle, Wash.
- 15. Portland, Ore.
- 16. Boise, Idaho
- 17. Salt Lake City, Utah

- Carson City, Nev. Los Angeles, Calif.
- Phoenix, Ariz.
- Santa Fe, N. M. 22. Denver, Colo.
- 23. Chevenne, Wvo.
- 24. Omaha, Neb.
- Des Moines, Iowa
- 26. Kansas City, Mo. 27. Topeka, Kans.
- 28. Oklahoma City, Okla.
- 29. Dallas, Tex. 30. Little Rock, Ark.
- Memphis, Tenn.
- 32. Jackson, Miss.
- 33. New Orleans, La.

- 34. Birmingham, Ala.
- 35. Atlanta, Ga.
- 36. Jacksonville, Fla.
- 37. Columbia, S. C.
- 38. Raleigh, N. C. 39. Richmond, Va.
- 40. Washington, D. C.
- 41. Boston, Mass.
- 42. Portland, Me.
- A. Baltimore, Md.
- B. Wilmington, Del. C. Philadelphia, Penn.
- D. Newark, N. J.
- E. New York, N. Y.
- F. Hartford, Conn.
- G. Providence, R. I.

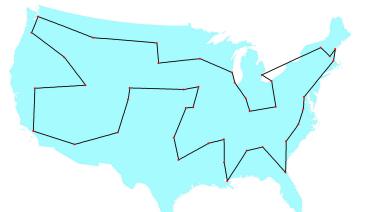
Combinatorial Explosion





Solution of this TSP problem

Dantzig, Fulkerson and Johnson found an optimal tour through 42 cities.



http://www.math.uwaterloo.ca/tsp/history/img/dantzig_big.html

TABLE I

ROAD DISTANCES BETWEEN CITIES IN ADJUSTED UNITS The figures in the table are mileages between the two specified numbered cities, less 11. 50 49 21 15 divided by 17, and rounded to the nearest integer. 61 62 21 48 60 16 17 18 59 60 15 20 26 17 10 62 66 20 25 31 22 15 40 44 50 41 12 108 117 66 71 77 68 61 51 46 13 145 149 104 108 114 106 99 88 84 14 | 181 18; 140 144 150 142 13; 124 120 99 8; 15 187 191 146 150 156 142 137 130 125 105 90 81 16 | 161 170 120 124 130 115 110 104 105 90 142 146 101 104 111 97 91 85 86 75 174 178 133 138 143 129 123 117 118 107 19 185 186 142 143 140 130 126 124 128 118 93 101 72 69 20 164 16; 120 123 124 106 106 10; 110 104 86 97 71 93 82 62 117 122 77 80 83 68 62 66 61 50 48 34 28 42 82 59 23 114 118 73 78 84 69 63 36 72 48 53 41 34 28 29 22 23 35 69 105 102 27 19 21 14 29 40 27 36 66 33 36 30 77 115 110 83 63 97 85 119 115 88 66 98 34 45 48 46 71 96 130 126 98 75 98 85 57 59 38 43 49 60 71 103 141 136 109 90 115 99 51 63 75 106 142 140 112 93 126 108 88 60 43 38 22 26 32 36 76 87 120 155 150 123 100 123 109 86 62 71 60 86 97 126 160 155 128 104 128 113 90 67 76 62 78 89 121 159 155 127 108 136 124 101 81 54 50 31 25 32 41 46 64 83 90 130 164 160 133 114 146 134 111 85 42 44 51 60 66 83 102 110 147 185 179 155 133 159 146 122 98 105 107 79 52 71 93 98 136 172 172 148 126 158 147 124 121 97 99 71 65 36 47 53 73 96 99 137 176 178 151 131 163 159 135 108 102 103 73 67 64 69 34 36 46 51 70 93 97 134 171 176 151 129 161 163 139 118 102 101 71 65 65 70 35 33 40 45 65 87 91 117 166 171 144 125 157 156 139 113 95 97 67 60 62 67

58 63 83 105 109 147 186 188 164 144 176 182 161 134 119 116 86 78 84 88 101 108 88 80 86 61 66 84 111 113 150 186 192 166 147 180 188 167 140 124 119 90 87 90 94 107 114 77 86 92 98 80

8 9 10 11 12 13 14 15 16 17 18 19 20 21 22 23 24 25 26 27 28 29 30 31 32 33 34 35 36 37 38 39 40 41

79 82 62 53 59 66 45 38 45 27 15 6

Hence this is an instance of the Metric TSP, but not Euclidean TSP.

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Idea: Indicator variable x(i, j), i > j, which is one if the tour includes edge $\{i, j\}$ (in either direction)

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minimize
$$\sum_{i=1}^{42} \sum_{j=1}^{i-1} c(i,j) x(i,j)$$
 subject to
$$\sum_{j < i} x(i,j) + \sum_{j > i} x(j,i) = 2 \qquad \text{for each } 1 \leq i \leq 42$$

$$0 \leq x(i,j) \leq 1 \quad \text{for each } 1 \leq j < i \leq 42$$

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- As long as solution of LP has fractional $x(i,j) \in (0,1)$:
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 - Return best of these two solutions

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 - Return best of these two solutions
- If solution of LP integral, return objective value

Bound-Step: If the best known integral solution so far is better than the solution of a LP, no need to explore branch further!

Outline

Introduction

Examples of TSP Instances

Demonstration

In the following, there are a few different runs of the demo.	

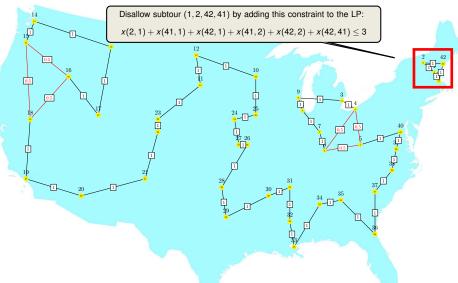
Iteration 1:



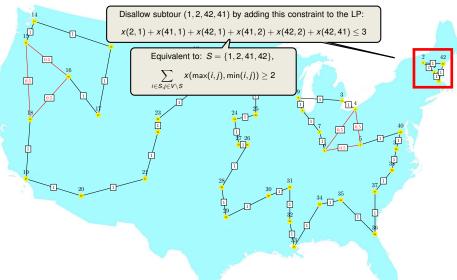
Iteration 1: Eliminate Subtour 1, 2, 41, 42



Iteration 1: Eliminate Subtour 1, 2, 41, 42



Iteration 1: Eliminate Subtour 1, 2, 41, 42

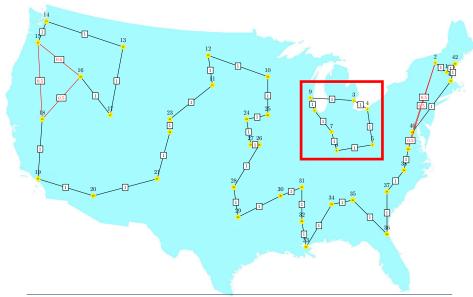


Iteration 2:



Iteration 2: Eliminate Subtour 3 – 9

Objective value: -676.000000, 861 variables, 946 constraints, 1802 iterations



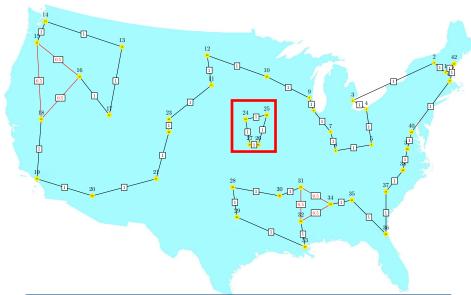
Iteration 3:

Objective value: -681.000000, 861 variables, 947 constraints, 1984 iterations



Iteration 3: Eliminate Subtour 24, 25, 26, 27

Objective value: -681.000000, 861 variables, 947 constraints, 1984 iterations



Iteration 4:

Objective value: -682.500000, 861 variables, 948 constraints, 1492 iterations



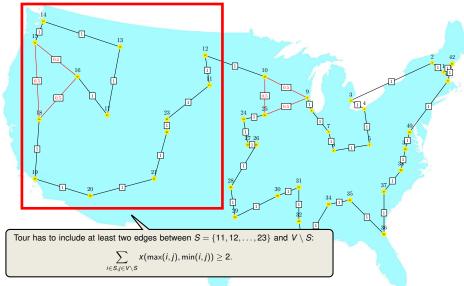
Iteration 4: Eliminate Cut 11 – 23

Objective value: -682.500000, 861 variables, 948 constraints, 1492 iterations



Iteration 4: Eliminate Cut 11 – 23

Objective value: -682.500000, 861 variables, 948 constraints, 1492 iterations



Iteration 5:

Objective value: -686.000000, 861 variables, 949 constraints, 2446 iterations



Iteration 5: Eliminate Subtour 13 – 23

Objective value: -686.000000, 861 variables, 949 constraints, 2446 iterations



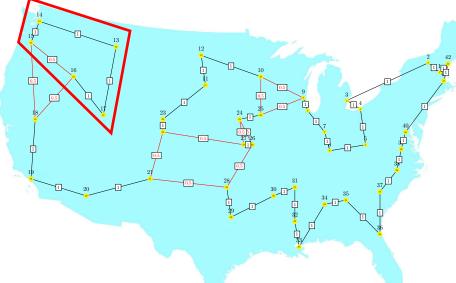
Iteration 6:

Objective value: -694.500000, 861 variables, 950 constraints, 1690 iterations



Iteration 6: Eliminate Cut 13 – 17

Objective value: -694.500000, 861 variables, 950 constraints, 1690 iterations



Iteration 7:

Objective value: -697.000000, 861 variables, 951 constraints, 2212 iterations



Iteration 7: Branch 1a $x_{18,15} = 0$

Objective value: -697.000000, 861 variables, 951 constraints, 2212 iterations



Iteration 8:

Objective value: -698.000000, 861 variables, 952 constraints, 1878 iterations



Iteration 8: Branch 2a $x_{17,13} = 0$

Objective value: -698.000000, 861 variables, 952 constraints, 1878 iterations



Iteration 9:

Objective value: -699.000000, 861 variables, 953 constraints, 2281 iterations



Iteration 9: Branch 2b $x_{17,13} = 1$

Objective value: -699.000000, 861 variables, 953 constraints, 2281 iterations



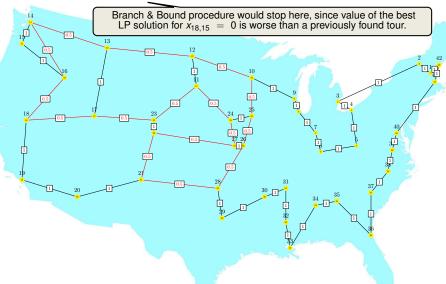
Iteration 10:

Objective value: -700.000000, 861 variables, 954 constraints, 2398 iterations



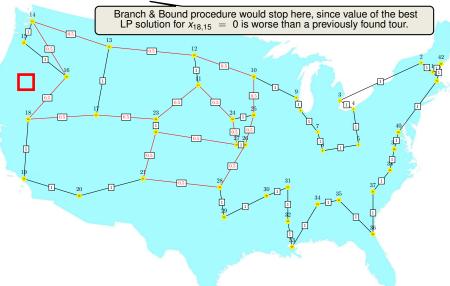
Iteration 10:

Objective value: -700.000000, 861 variables, 954 constraints, 2398 iterations



Iteration 10: Branch 1b $x_{18.15} = 1$

Objective value: -700.000000, 861 variables, 954 constraints, 2398 iterations



Iteration 11:

Objective value: -701.000000, 861 variables, 953 constraints, 2506 iterations



Iteration 11: Branch & Bound terminates

Objective value: -701.000000, 861 variables, 953 constraints, 2506 iterations



1: LP solution 641

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Eliminate Subtour 1, 2, 41, 42

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2: LP solution 676

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Leliminate Subtour 1, 2, 41, 42

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Eliminate Subtour 3 – 9

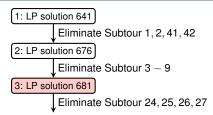
1: LP solution 641

Eliminate Subtour 1, 2, 41, 42

2: LP solution 676

Eliminate Subtour 3 – 9

3: LP solution 681



```
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Eliminate Subtour 1, 2, 41, 42

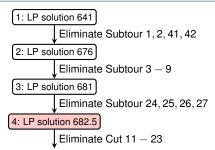
2: LP solution 676

Eliminate Subtour 3 – 9

3: LP solution 681

Eliminate Subtour 24, 25, 26, 27

4: LP solution 682.5
```



```
1: LP solution 641

Eliminate Subtour 1, 2, 41, 42

2: LP solution 676

Eliminate Subtour 3 - 9

3: LP solution 681

Eliminate Subtour 24, 25, 26, 27

4: LP solution 682.5

Eliminate Cut 11 - 23

5: LP solution 686
```

```
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Eliminate Subtour 3 – 9

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Eliminate Cut 11 – 23

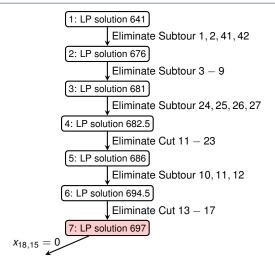
5: LP solution 686

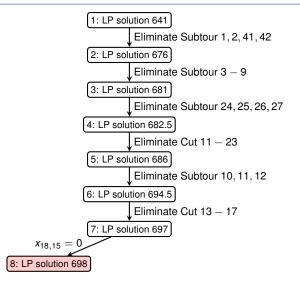
Eliminate Subtour 10, 11, 12
```

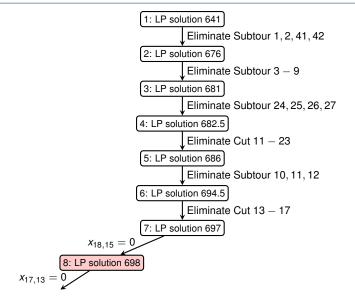
```
1: LP solution 641
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2: LP solution 676
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          Eliminate Subtour 24, 25, 26, 27
4: LP solution 682.5
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5: LP solution 686
          Eliminate Subtour 10, 11, 12
6: LP solution 694.5
```

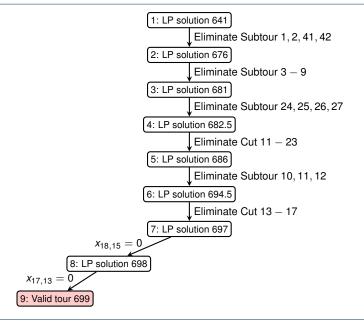
```
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```

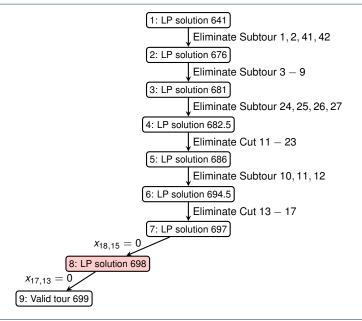
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7: LP solution 697
```

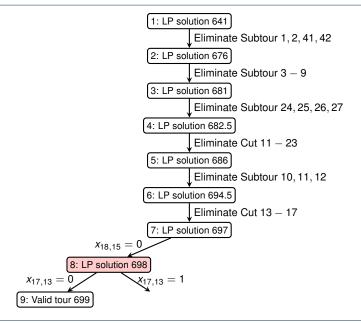


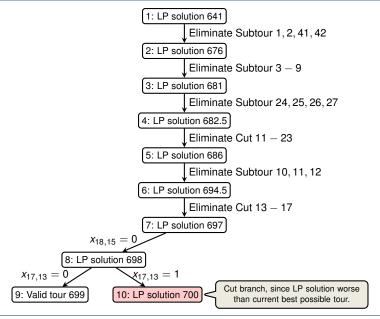


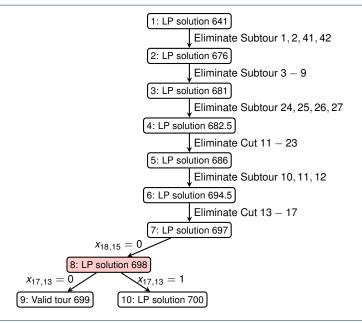


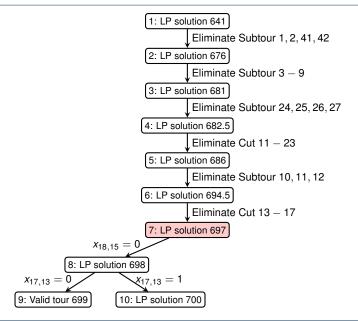


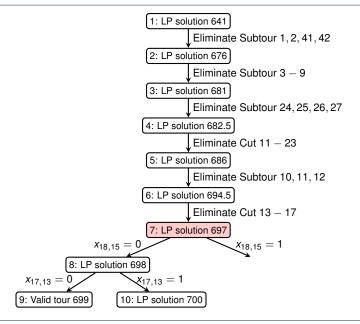


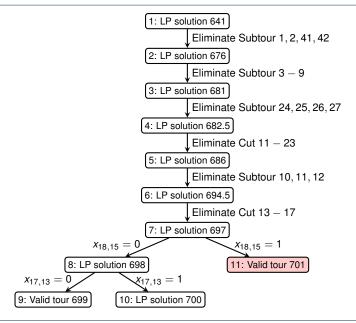


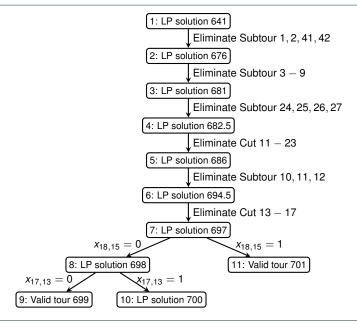








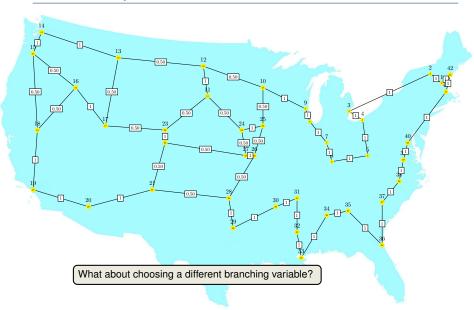




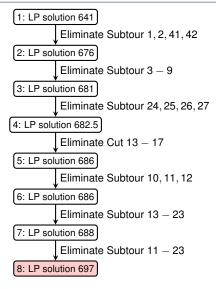
Iteration 7: Objective 697



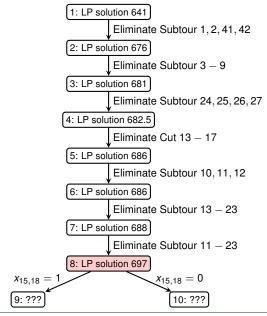
Iteration 7: Objective 697



Solving Progress (Alternative Branch 1)



Solving Progress (Alternative Branch 1)



Alternative Branch 1: $x_{18,15}$, Objective 697



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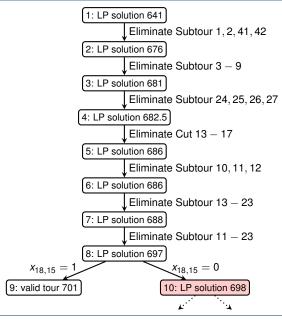
Alternative Branch 1a: $x_{18,15} = 1$, Objective 701 (Valid Tour)



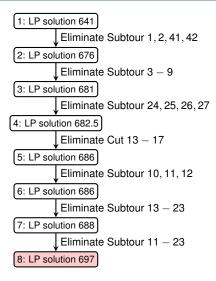
Alternative Branch 1b: $x_{18,15} = 0$, Objective 698



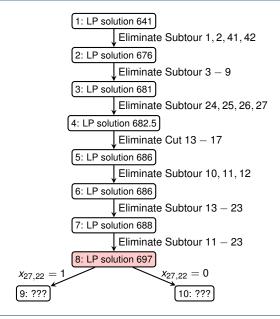
Solving Progress (Alternative Branch 1)



Solving Progress (Alternative Branch 2)



Solving Progress (Alternative Branch 2)



Alternative Branch 2: $x_{27,22}$, Objective 697



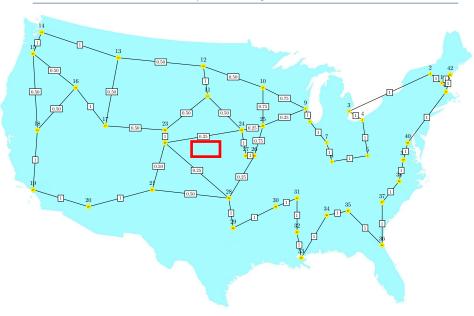
Alternative Branch 2: $x_{27,22}$, Objective 697



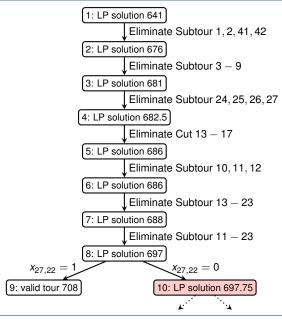
Alternative Branch 2a: $x_{27,22} = 1$, Objective 708 (Valid tour)



Alternative Branch 2b: $x_{27,22} = 0$, Objective 697.75



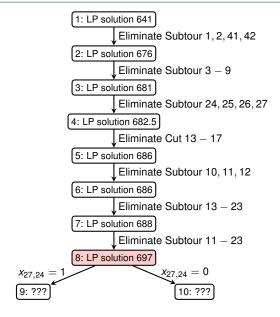
Solving Progress (Alternative Branch 2)



Solving Progress (Alternative Branch 3)

```
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          Eliminate Subtour 3 - 9
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4: LP solution 682.5
          Eliminate Cut 13 - 17
5: LP solution 686
          Eliminate Subtour 10, 11, 12
 6: LP solution 686
          Eliminate Subtour 13 - 23
7: LP solution 688
          Eliminate Subtour 11 - 23
8: LP solution 697
```

Solving Progress (Alternative Branch 3)



Alternative Branch 3: $x_{27,24}$, Objective 697



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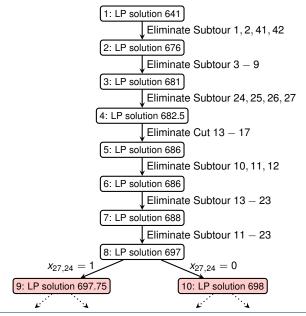
Alternative Branch 3a: $x_{27,24} = 1$, Objective 697.75



Alternative Branch 3b: $x_{27,24} = 0$, Objective 698



Solving Progress (Alternative Branch 3)



Solving Progress (Alternative Branch 3)

```
1: LP solution 641

Eliminate Subtour 1, 2, 41, 42

2: LP solution 676

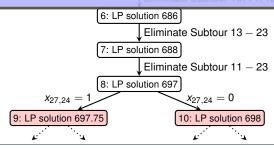
Eliminate Subtour 3 – 9

3: LP solution 681

Eliminate Subtour 24, 25, 26, 27

4: LP solution 682.5
```

Not only do we have to explore (and branch further in) both subtrees, but also the optimal tour is in the subtree with larger LP solution!



Conclusion (1/2)

How can one generate these constraints automatically?

Conclusion (1/2)

How can one generate these constraints automatically?
 Subtour Elimination: Finding Connected Components
 Small Cuts: Finding the Minimum Cut in Weighted Graphs

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 BFS may be more attractive, even though it might need more memory.

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- Why don't we add all possible Subtour Eliminiation constraints to the LP? There are exponentially many of them!
- Should the search tree be explored by BFS or DFS?
 BFS may be more attractive, even though it might need more memory.

CONCLUDING REMARK

It is clear that we have left unanswered practically any question one might pose of a theoretical nature concerning the traveling-salesman problem; however, we hope that the feasibility of attacking problems involving a moderate number of points has been successfully demonstrated, and that perhaps some of the ideas can be used in problems of similar nature.

- Eliminate Subtour 1, 2, 41, 42
- Eliminate Subtour 3 9
- Eliminate Subtour 10, 11, 12
- Eliminate Subtour 11 23
- Eliminate Subtour 13 23
- Eliminate Cut 13 17
- Eliminate Subtour 24, 25, 26, 27

- Eliminate Subtour 1, 2, 41, 42
- Eliminate Subtour 3 9
- Eliminate Subtour 10, 11, 12
- Eliminate Subtour 11 23
- Eliminate Subtour 13 23
- Eliminate Cut 13 17
- Eliminate Subtour 24, 25, 26, 27

THE 49-CITY PROBLEM*

The optimal tour \bar{x} is shown in Fig. 16. The proof that it is optimal is given in Fig. 17. To make the correspondence between the latter and its programming problem clear, we will write down in addition to 42 relations in non-negative variables (2), a set of 25 relations which suffice to prove that D(x) is a minimum for \bar{x} . We distinguish the following subsets of the 42 cities:

```
\begin{array}{lll} S_1 = \{1,\,2,\,41,\,42\} & S_5 = \{13,\,14,\,\cdots,\,23\} \\ S_2 = \{3,\,4,\,\cdots,\,9\} & S_6 = \{13,\,14,\,15,\,16,\,17\} \\ S_3 = \{11,\,2,\,\cdots,\,9,\,29,\,30,\,\cdots,\,42\} & S_7 = \{24,\,25,\,26,\,27\}. \\ S_4 = \{11,\,12,\,\cdots,\,23\} & S_7 = \{24,\,25,\,26,\,27\}. \end{array}
```



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Contact page

Tools

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From Wikipedia, the free encyclopedia

IBM ILOG CPLEX Optimization Studio (often informally referred to simply as CPLEX) is an optimization software package. In 2004, the work on CPLEX earned the first INFORMS Impact Prize.

The CPLEX Optimizer was named for the simplex method as implemented in the C programming language. although today it also supports other types of mathematical optimization and offers interfaces other than just C. It was originally developed by Robert E. Bixby and was offered commercially starting in 1988 by

CPLEX Optimization Inc., which was acquired by ILOG in 1997; ILOG was subsequently acquired by IBM in January 2009.^[1] CPLEX continues to be actively developed under IBM.

The IBM ILOG CPLEX Optimizer solves integer programming problems, very large^[2] linear programming problems using either primal or dual variants of the simplex method or the barrier interior

CPLEX

Developer(s) IRM Stable release 12.6

Development status Active

Type Technical computing License Proprietary

Website ibm.com/software

/products /ibmilogcpleoptistud/₫ Welcome to IBM(R) ILOG(R) CPLEX(R) Interactive Optimizer 12.6.1.0 with Simplex. Mixed Integer & Barrier Optimizers 5725-A06 5725-A29 5724-Y48 5724-Y49 5724-Y54 5724-Y55 5655-Y21 Copyright IBM Corp. 1988, 2014. All Rights Reserved. Type 'help' for a list of available commands. Type 'help' followed by a command name for more information on commands. CPLEX> read tsp.lp Problem 'tsp.lp' read. Read time = 0.00 sec. (0.06 ticks) CPLEX> primopt Tried aggregator 1 time. LP Presolve eliminated 1 rows and 1 columns. Reduced LP has 49 rows. 860 columns. and 2483 nonzeros. Presolve time = 0.00 sec. (0.36 ticks)Iteration log . . . Iteration: 1 Infeasibility = 33,999999 Iteration: 26 Objective 1510,000000 Objective = Iteration: 90 923,000000 Iteration: 155 Objective 711.000000 Primal simplex - Optimal: Objective = 6.9900000000e+02 Solution time = 0.00 sec. Iterations = 168 (25) Deterministic time = 1.16 ticks (288.86 ticks/sec)

CPLEX>

CPLEX> display	solut	ior				
Variable Name			Sol	lution		
x_2_1					900000	
x_42_1					900000	
x_3_2				1.6	000000	
x_4_3				1.6	000000	
x_5_4				1.6	00000	
x_6_5				1.6	00000	
x_7_6				1.6	00000	
x 8 7				1.6	00000	
x 9 8				1.6	999999	
x 10 9				1.6	900000	
x_11_10					000000	
x 12 11					900000	
x_13_12					900000	
x_14_13					900000	
x_15_14					900000	
x 16 15					900000	
					000000	
x_17_16						
x_18_17					900000	
x_19_18					900000	
x_20_19					00000	
x_21_20					900000	
x_22_21					909999	
x_23_22					900000	
x_24_23					900000	
x_25_24				1.6	900000	
x_26_25				1.6	000000	
x_27_26				1.6	00000	
x_28_27				1.6	000000	
x_29_28				1.6	00000	
x_30_29				1.6	00000	
x_31_30				1.6	00000	
x_32_31				1.6	909999	
x 33 32				1.6	900000	
x_34_33					000000	
x_35_34					000000	
x_36_35					900000	
x_37_36					900000	
x_37_36 1.000000 x_38_37 1.000000						
x_38_37 1.000000 x_39_38 1.000000						
x_41_40						
x_42_41						
All other varia	ables	in	the	range	1-861	are

Randomised Algorithms

Lecture 9: Approximation Algorithms: MAX-3-CNF and Vertex-Cover

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Randomised Approximation

MAX-3-CNF

Weighted Vertex Cover

Approximation Ratio —

A randomised algorithm for a problem has approximation ratio $\rho(n)$, if for any input of size n, the expected cost (value) $\mathbf{E}[C]$ of the returned solution and optimal cost C^* satisfy:

$$\max\left(\frac{\mathbf{E}\left[\,C\,
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- Maximisation problem: $\frac{C^*}{\mathbf{E}[C]} \ge 1$ Minimisation problem: $\frac{\mathbf{E}[C]}{C^*} \ge 1$

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Randomised Approximation Schemes

not covered here (non-examinable)

An approximation scheme is an approximation algorithm, which given any input and $\epsilon > 0$, is a $(1 + \epsilon)$ -approximation algorithm.

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Randomised Approximation Schemes

An approximation scheme is an approximation algorithm, which given any input and $\epsilon > 0$, is a $(1 + \epsilon)$ -approximation algorithm.

- It is a polynomial-time approximation scheme (PTAS) if for any fixed $\epsilon > 0$, the runtime is polynomial in n. For example, $O(n^{2/\epsilon})$.
- It is a fully polynomial-time approximation scheme (FPTAS) if the runtime is polynomial in both $1/\epsilon$ and n. For example, $O((1/\epsilon)^2 \cdot n^3)$.

Outline

Randomised Approximation

MAX-3-CNF

Weighted Vertex Cover

• Given: 3-CNF formula, e.g.: $(x_1 \lor x_3 \lor \overline{x_4}) \land (x_2 \lor \overline{x_3} \lor \overline{x_5}) \land \cdots$

MAX-3-CNF Satisfiability —

- Given: 3-CNF formula, e.g.: $(x_1 \lor x_3 \lor \overline{x_4}) \land (x_2 \lor \overline{x_3} \lor \overline{x_5}) \land \cdots$
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Relaxation of the satisfiability problem. Want to compute how "close" the formula to being satisfiable is.

Assume that no literal (including its negation) appears more than once in the same clause.

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Example:

$$(x_1 \vee x_3 \vee \overline{x_4}) \wedge (x_1 \vee \overline{x_3} \vee \overline{x_5}) \wedge (x_2 \vee \overline{x_4} \vee x_5) \wedge (\overline{x_1} \vee x_2 \vee \overline{x_3})$$

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$$x_1 = 1, x_2 = 0, x_3 = 1, x_4 = 0 \text{ and } x_5 = 1 \text{ satisfies 3 (out of 4 clauses)}$$

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$$x_1 = 1, x_2 = 0, x_3 = 1, x_4 = 0 \text{ and } x_5 = 1 \text{ satisfies 3 (out of 4 clauses)}$$

Idea: What about assigning each variable uniformly and independently at random?

Theorem 35.6 -

Given an instance of MAX-3-CNF with n variables x_1, x_2, \ldots, x_n and m clauses, the randomised algorithm that sets each variable independently at random is a randomised 8/7-approximation algorithm.

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$$P[clause i is not satisfied] = $\frac{1}{2} \cdot \frac{1}{2} \cdot \frac{1}{2} = \frac{1}{8}$
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(Linearity of Expectations)

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⇒ P[clause *i* is satisfied] = $1 - \frac{1}{8} = \frac{7}{8}$

⇒ E[Y_i] = P[Y_i = 1] · 1 = $\frac{7}{8}$.

$$\mathbf{E}[Y] = \mathbf{E}\left[\sum_{i=1}^{m} Y_{i}\right] = \sum_{i=1}^{m} \mathbf{E}[Y_{i}]$$
[Linearity of Expectations]

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Linearity of Expectations

Theorem 35.6

Given an instance of MAX-3-CNF with n variables x_1, x_2, \ldots, x_n and m clauses, the randomised algorithm that sets each variable independently at random is a randomised 8/7-approximation algorithm.

Proof:

• For every clause i = 1, 2, ..., m, define a random variable:

$$Y_i = \mathbf{1}\{\text{clause } i \text{ is satisfied}\}$$

Since each literal (including its negation) appears at most once in clause i,

$$\mathbf{P}[\text{clause } i \text{ is not satisfied}] = \frac{1}{2} \cdot \frac{1}{2} \cdot \frac{1}{2} = \frac{1}{8}$$
⇒
$$\mathbf{P}[\text{clause } i \text{ is satisfied}] = 1 - \frac{1}{8} = \frac{7}{8}$$
⇒
$$\mathbf{E}[Y_i] = \mathbf{P}[Y_i = 1] \cdot 1 = \frac{7}{8}.$$

• Let $Y := \sum_{i=1}^{m} Y_i$ be the number of satisfied clauses. Then,

$$\mathbf{E}[Y] = \mathbf{E}\left[\sum_{i=1}^{m} Y_i\right] = \sum_{i=1}^{m} \mathbf{E}[Y_i] = \sum_{i=1}^{m} \frac{7}{8} = \frac{7}{8} \cdot m. \quad \Box$$
Linearity of Expectations
maximum number of satisfiable clauses is m.

- Theorem 35.6 -

Given an instance of MAX-3-CNF with n variables x_1, x_2, \ldots, x_n and m clauses, the randomised algorithm that sets each variable independently at random is a polynomial-time randomised 8/7-approximation algorithm.

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Given an instance of MAX-3-CNF with n variables x_1, x_2, \ldots, x_n and m clauses, the randomised algorithm that sets each variable independently at random is a polynomial-time randomised 8/7-approximation algorithm.

Corollary

For any instance of MAX-3-CNF, there exists an assignment which satisfies at least $\frac{7}{8}$ of all clauses.

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There is $\omega \in \Omega$ such that $Y(\omega) \geq \mathbf{E}[Y]$

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Probabilistic Method: powerful tool to show existence of a non-obvious property.

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Corollary -

Any instance of MAX-3-CNF with at most 7 clauses is satisfiable.

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Follows from the previous Corollary.

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Algorithm: Assign x_1 so that the conditional expectation is maximised and recurse.

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GREEDY-3-CNF(ϕ , n, m)

- 1: **for** j = 1, 2, ..., n
- 2: Compute **E**[$Y \mid x_1 = v_1 \dots, x_{i-1} = v_{i-1}, x_i = 1$]
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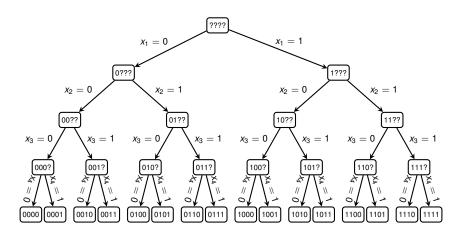
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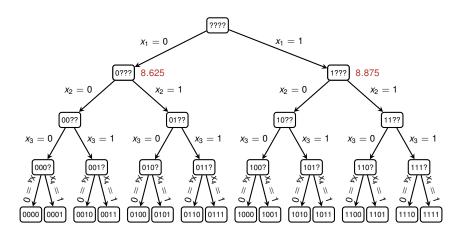
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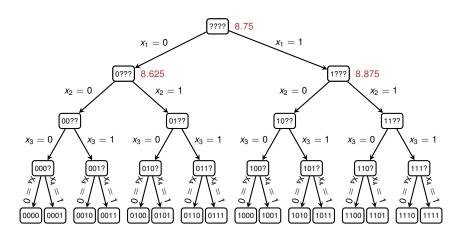
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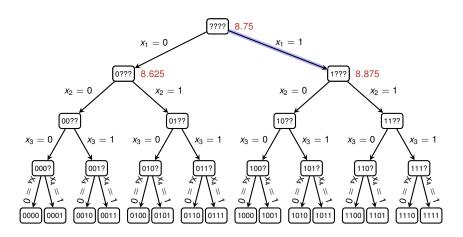
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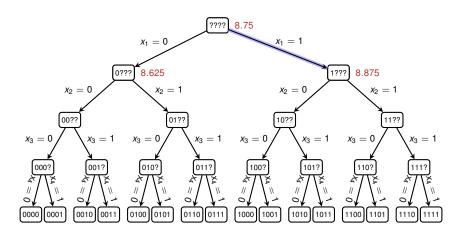
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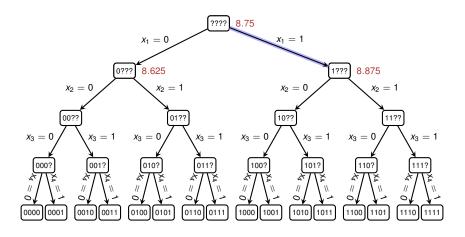
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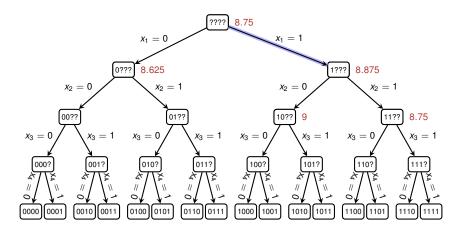
 $(\underline{x} \vee \underline{x} \vee \underline{x}) \wedge (\underline{x} \vee \underline{x}) \wedge (\underline{x} \vee \underline{x} \vee \underline{x}) \wedge (\underline{x} \vee \underline{x})$



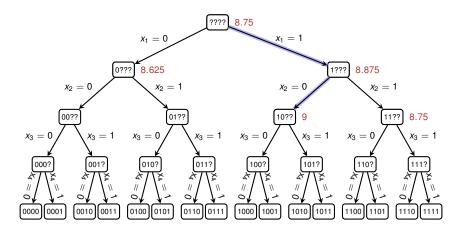
$$1 \wedge 1 \wedge 1 \wedge (\overline{x_3} \vee x_4) \wedge 1 \wedge (\overline{x_2} \vee \overline{x_3}) \wedge (x_2 \vee x_3) \wedge (\overline{x_2} \vee x_3) \wedge 1 \wedge (x_2 \vee \overline{x_3} \vee \overline{x_4})$$



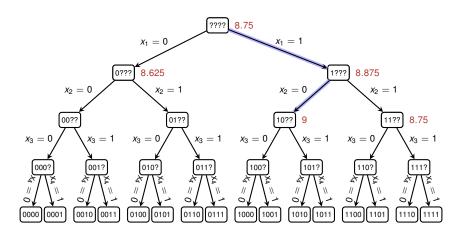
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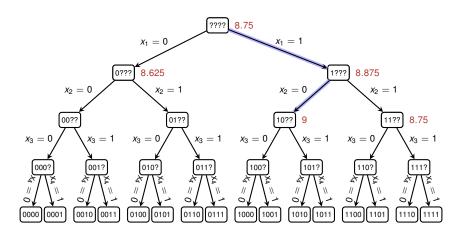
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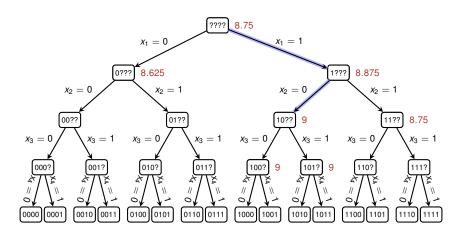
 $1 \wedge 1 \wedge 1 \wedge (\overline{x_3} \vee x_4) \wedge 1 \wedge (\overline{x_2} \vee \overline{x_3}) \wedge (\cancel{x_2} \vee x_3) \wedge (\overline{x_2} \vee x_3) \wedge 1 \wedge (\cancel{x_2} \vee \overline{x_3} \vee \overline{x_4})$



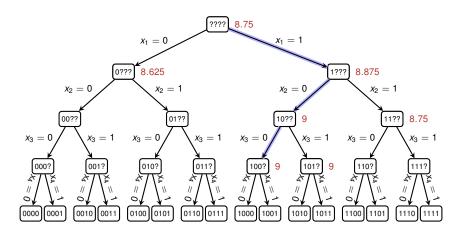
$$1 \wedge 1 \wedge 1 \wedge (\overline{x_3} \vee x_4) \wedge 1 \wedge 1 \wedge (x_3) \wedge 1 \wedge 1 \wedge (\overline{x_3} \vee \overline{x_4})$$



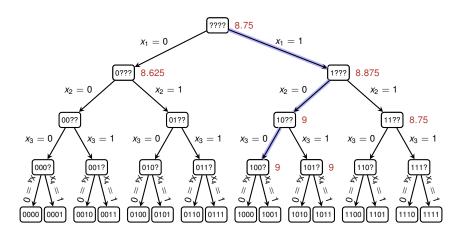
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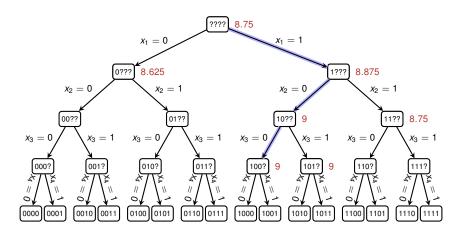


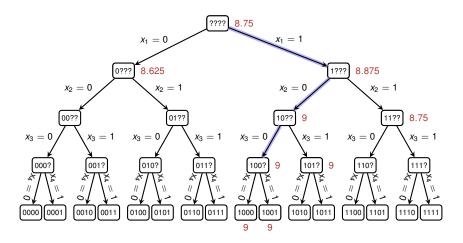
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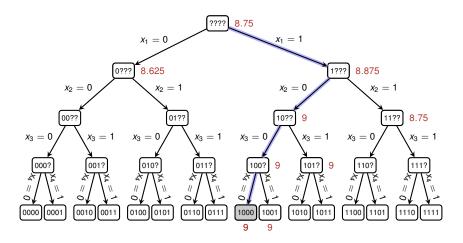


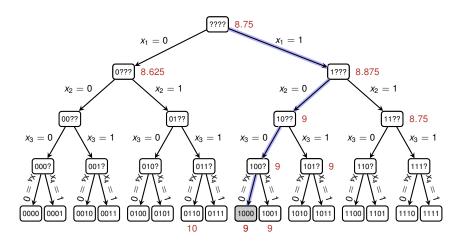
 $1 \land 1 \land 1 \land (\overline{x_2} \lor x_4) \land 1 \land 1 \land (x_3) \land 1 \land 1 \land (\overline{x_3} \lor \overline{x_4})$

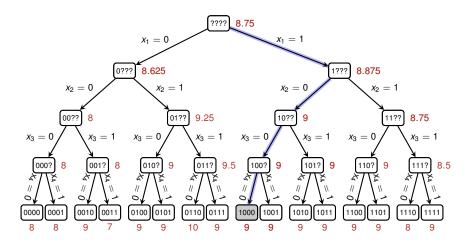


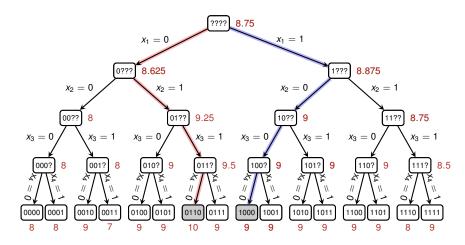






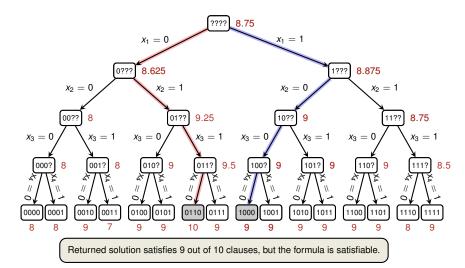






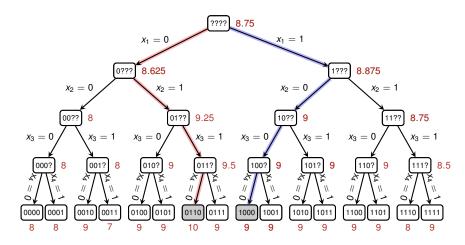
Run of GREEDY-3-CNF(φ , n, m)

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$$\mathbf{E}\left[\;Y\mid x_{1}=v_{1},\ldots,x_{j-1}=v_{j-1},x_{j}=v_{j}\;\right]\geq\mathbf{E}\left[\;Y\mid x_{1}=v_{1},\ldots,x_{j-1}=v_{j-1}\;\right]$$

This algorithm is deterministic.

Theorem

GREEDY-3-CNF(ϕ , n, m) is a polynomial-time 8/7-approximation.

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 - In iteration j = 1, 2, ..., n, $Y = Y(\phi)$ averages over 2^{n-j+1} assignments
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Theorem 35.6 ——

Given an instance of MAX-3-CNF with n variables x_1, x_2, \ldots, x_n and m clauses, the randomised algorithm that sets each variable independently at random is a randomised 8/7-approximation algorithm.

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Theorem (Hastad'97)

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Essentially there is nothing smarter than just guessing!

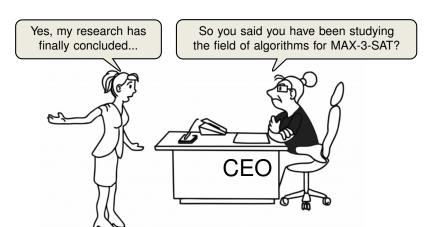


Source of Image: Stefan Szeider, TU Vienna

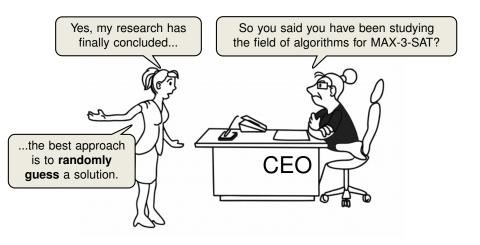
So you said you have been studying the field of algorithms for MAX-3-SAT?



Source of Image: Stefan Szeider, TU Vienna



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Outline

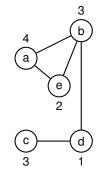
Randomised Approximation

MAX-3-CNF

Weighted Vertex Cover

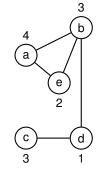
Vertex Cover Problem

- Given: Undirected, vertex-weighted graph G = (V, E)
- Goal: Find a minimum-weight subset $V' \subseteq V$ such that if $\{u, v\} \in E(G)$, then $u \in V'$ or $v \in V'$.



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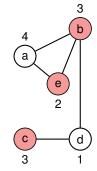




Question: How can we deal with graphs that have negative weights?

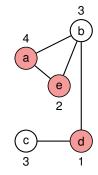
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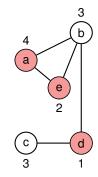
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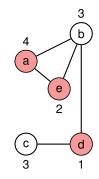
This is (still) an NP-hard problem.



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Applications:

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Applications:

 Every edge forms a task, and every vertex represents a person/machine which can execute that task

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This is (still) an NP-hard problem.

Applications:

- Every edge forms a task, and every vertex represents a person/machine which can execute that task
- Weight of a vertex could be salary of a person
- Perform all tasks with the minimal amount of resources

```
APPROX-VERTEX-COVER (G)

1 C = \emptyset

2 E' = G.E

3 while E' \neq \emptyset

4 let (u, v) be an arbitrary edge of E'

5 C = C \cup \{u, v\}

7 return C
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This algorithm is a 2-approximation for unweighted graphs!

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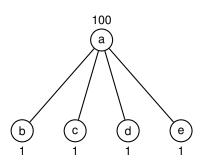
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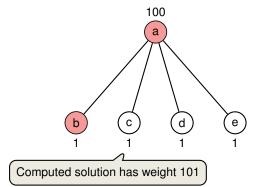
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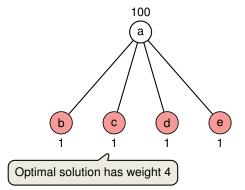
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minimize $\sum_{v \in V} w(v)x(v)$ subject to $x(u) + x(v) \geq 1 \qquad \text{for each } (u,v) \in E$ $x(v) \in \{0,1\} \qquad \text{for each } v \in V$

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Linear Program ———

Idea: Round the solution of an associated linear program.

```
0-1 Integer Program —
              \sum_{v\in V}w(v)x(v)
minimize
               x(u) + x(v) > 1 for each (u, v) \in E
subject to
                       x(v) \in \{0,1\} for each v \in V
                    optimum is a lower bound on the optimal
                       weight of a minimum weight-cover.
 Linear Program
               \sum w(v)x(v)
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Rounding Rule: if $x(v) \ge 1/2$ then round up, otherwise round down.

Linear Program

The Algorithm

```
APPROX-MIN-WEIGHT-VC(G,w)

1 C=\emptyset

2 compute \bar{x}, an optimal solution to the linear program

3 for each \nu \in V

4 if \bar{x}(\nu) \geq 1/2

5 C=C \cup \{\nu\}

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Theorem 35.7 -

 $\label{lem:approx-Min-Weight-VC} \mbox{ is a polynomial-time 2-approximation algorithm for the minimum-weight vertex-cover problem.}$

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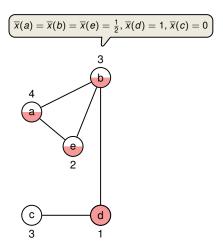
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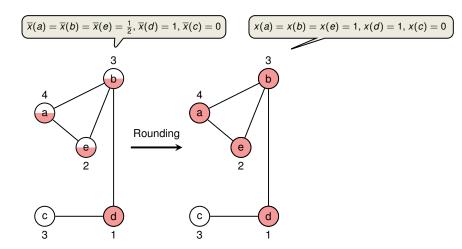
is polynomial-time because we can solve the linear program in polynomial time

Example of APPROX-MIN-WEIGHT-VC



fractional solution of LP with weight = 5.5

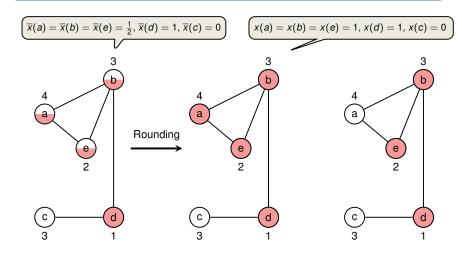
Example of APPROX-MIN-WEIGHT-VC



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rounded solution of LP with weight = 10

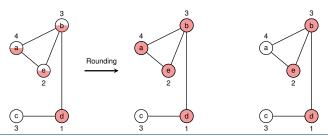
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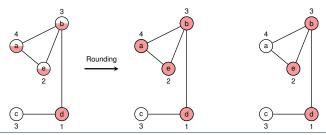
rounded solution of LP with weight = 10

optimal solution with weight = 6

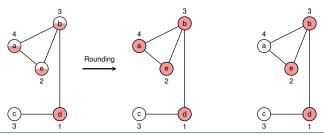


Proof (Approximation Ratio is 2 and Correctness):

• Let C^* be an optimal solution to the minimum-weight vertex cover problem

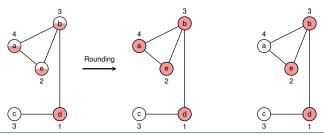


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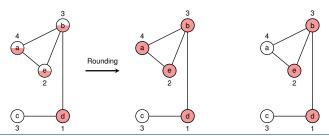


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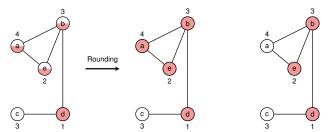
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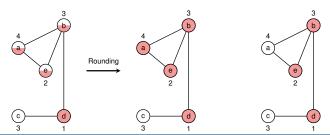
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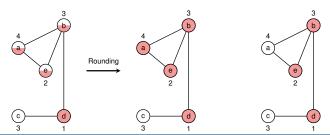
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 - \Rightarrow at least one of $\overline{x}(u)$ and $\overline{x}(v)$ is at least 1/2



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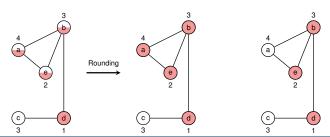
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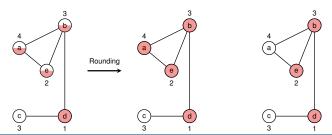
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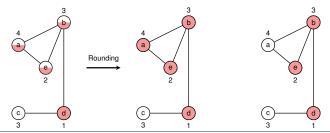


- Let C* be an optimal solution to the minimum-weight vertex cover problem
- Let z^* be the value of an optimal solution to the linear program, so

$$z^* \leq w(C^*)$$

- **Step 1:** The computed set *C* covers all vertices:
 - Consider any edge $(u, v) \in E$ which imposes the constraint $x(u) + x(v) \ge 1$ \Rightarrow at least one of $\overline{x}(u)$ and $\overline{x}(v)$ is at least $1/2 \Rightarrow C$ covers edge (u, v)
- Step 2: The computed set C satisfies $w(C) \le 2z^*$:

$$w(C^*) \geq z^*$$

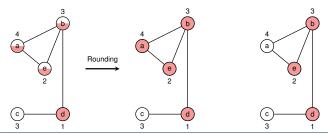


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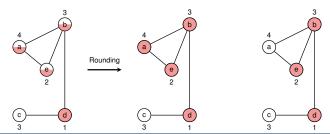


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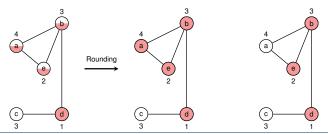


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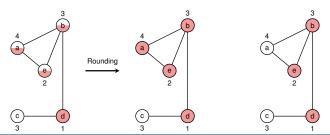


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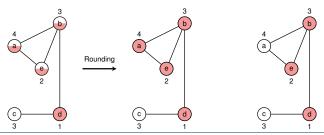


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Randomised Algorithms

Lecture 10: Approximation Algorithms: Set-Cover and MAX-CNF

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025



Outline

Weighted Set Cover

MAX-CNF

Set Cover Problem —

- Given: set X, |X| = n, a family of subsets \mathcal{F} , and cost function $c : \mathcal{F} \to \mathbb{R}^+$
- ullet Goal: Find a minimum-cost subset $\mathcal{C} \subseteq \mathcal{F}$

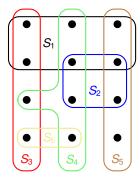
s.t.
$$X = \bigcup_{S \in \mathcal{C}} S$$
.

Set Cover Problem -

- Given: set X, |X| = n, a family of subsets \mathcal{F} , and cost function $c : \mathcal{F} \to \mathbb{R}^+$
- Goal: Find a minimum-cost subset $C \subseteq \mathcal{F}$

Sum over the costs of all sets in $\mathcal C$

s.t.
$$X = \bigcup_{S \in \mathcal{C}} S$$
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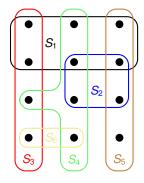


Set Cover Problem -

- Given: set X, |X| = n, a family of subsets \mathcal{F} , and cost function $c : \mathcal{F} \to \mathbb{R}^+$
- lacktriangle Goal: Find a minimum-cost subset $\mathcal{C} \subseteq \mathcal{F}$

Sum over the costs of all sets in \mathcal{C}

s.t.
$$X = \bigcup_{S \in \mathcal{C}} S$$
.



 S_1 S_2 S_3 S_4 S_5 S_6 c: 2 3 3 5 1 2

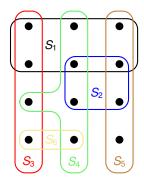
s.t.

Set Cover Problem —

- Given: set X, |X| = n, a family of subsets \mathcal{F} , and cost function $c : \mathcal{F} \to \mathbb{R}^+$
- Goal: Find a minimum-cost subset $C \subseteq \mathcal{F}$

Sum over the costs of all sets in C

$$X = \bigcup_{S \in \mathcal{C}} S.$$



 S_1 S_2 S_3 S_4 S_5 S_6 c: 2 3 3 5 1 2

Remarks:

- generalisation of the weighted Vertex-Cover problem
- models resource allocation problems

Set Cover Problem —

s.t.

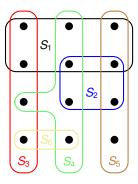
- Given: set X, |X| = n, a family of subsets \mathcal{F} , and cost function $c : \mathcal{F} \to \mathbb{R}^+$
- lacktriangle Goal: Find a minimum-cost subset $\mathcal{C} \subseteq \mathcal{F}$

Sum over the costs of all sets in \mathcal{C}

$$X = \bigcup_{S \in \mathcal{C}} S$$
.



Question: How can we reduce the Vertex-Cover problem to the Set-Cover problem?



Remarks:

- generalisation of the weighted Vertex-Cover problem
- models resource allocation problems

Setting up an Integer Program



Question: Try to formulate the integer program and linear program of the weighted SET-COVER problem (solution on next slide!)

Setting up an Integer Program

- 0-1 Integer Program ———

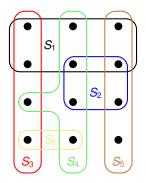
minimize
$$\sum_{S\in\mathcal{F}}c(S)y(S)$$
 subject to
$$\sum_{S\in\mathcal{F}:\ x\in\mathcal{S}}y(S)\ \geq\ 1\qquad \text{for each }x\in X$$

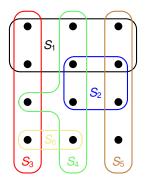
$$y(S)\ \in\ \{0,1\}\qquad \text{for each }S\in\mathcal{F}$$

Setting up an Integer Program

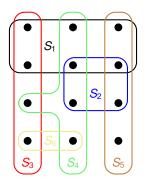
Linear Program
$$\sum_{S\in\mathcal{F}}c(S)y(S)$$
 subject to
$$\sum_{S\in\mathcal{F}\colon x\in S}y(S)\ \geq\ 1\qquad \text{for each }x\in X$$

$$y(S)\ \in\ [0,1]\qquad \text{for each }S\in\mathcal{F}$$

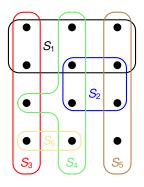




	S_1	S_2	S_3	S_4	S_5	S_6
c :	2	3	3	5	1	2
$\overline{y}(.)$:	1/2	1/2	1/2	1/2	1	1/2

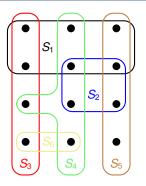


	S_1	S_2	S_3	S_4	S_5	S_6	
c :	2	3	3	5	1	2	
$\overline{y}(.)$:	1/2	1/2	1/2	1/2	1	1/2	Cost equals 8.5





The strategy employed for Vertex-Cover would take all 6 sets!



$$S_1$$
 S_2 S_3 S_4 S_5 S_6 $c:$ 2 3 3 5 1 2 $\overline{y}(.)$: 1/2 1/2 1/2 1/2 1/2 1 1/2

The strategy employed for Vertex-Cover would take all 6 sets!

Even worse: If all \overline{y} 's were below 1/2, we would not even return a valid cover!

Cost equals 8.5

	S ₁	S_2	S_3	S_4	S ₅	S ₆ 2 1/2
C :	2	3	3	5	1	2
$\overline{y}(.)$:	1/2	1/2	1/2	1/2	1	1/2

	S_1	S_2	S ₃	S_4	S ₅	S_6
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Idea: Interpret the \overline{y} -values as probabilities for picking the respective set.

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Randomised Rounding —

- Let $C \subseteq \mathcal{F}$ be a random set with each set S being included independently with probability $\overline{y}(S)$.
- More precisely, if \(\overline{y} \) denotes the optimal solution of the LP, then we compute an integral solution \(y \) by:

$$y(S) = \begin{cases} 1 & \text{with probability } \overline{y}(S) \\ 0 & \text{otherwise.} \end{cases}$$
 for all $S \in \mathcal{F}$.

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 for all $S \in \mathcal{F}$.

• Therefore, $\mathbf{E}[y(S)] = \overline{y}(S)$.

	S_1	S_2	S_3	S_4	S_5	S_6	
C :	2	3	3	5	1	2	
$\overline{y}(.)$:	1/2	1/2	1/2	1/2	1	1/2	

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- Lemma -

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Idea: Interpret the \overline{y} -values as probabilities for picking the respective set.

Lemma -

The expected cost satisfies

$$\mathbf{E}[c(\mathcal{C})] = \sum_{S \in \mathcal{F}} c(S) \cdot \overline{y}(S).$$

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Lemma -

The expected cost satisfies

$$\mathbf{E}[c(\mathcal{C})] = \sum_{S \in \mathcal{F}} c(S) \cdot \overline{y}(S).$$

• The probability that an element $x \in X$ is covered satisfies

$$\mathbf{P}\left[x\in\bigcup_{S\in\mathcal{C}}S\right]\geq 1-\frac{1}{e}.$$

Lemma

Let $\mathcal{C} \subseteq \mathcal{F}$ be a random subset with each set \mathcal{S} being included independently with probability $\overline{y}(\mathcal{S})$.

- The expected cost satisfies $\mathbf{E}[c(C)] = \sum_{S \in \mathcal{F}} c(S) \cdot \overline{y}(S)$.
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Proof:

$$\mathbf{E}[c(\mathcal{C})] = \mathbf{E}\left[\sum_{S \in \mathcal{C}} c(S)\right]$$

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Proof:

$$\mathbf{E}[c(\mathcal{C})] = \mathbf{E}\left[\sum_{S \in \mathcal{C}} c(S)\right] = \mathbf{E}\left[\sum_{S \in \mathcal{F}} \mathbf{1}_{S \in \mathcal{C}} \cdot c(S)\right]$$

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$$= \sum_{S \in \mathcal{F}} \mathbf{P}[S \in \mathcal{C}] \cdot c(S)$$

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Proof:

• Step 1: The expected cost of the random set $\mathcal{C} \checkmark$

$$\begin{split} \mathbf{E}\left[c(\mathcal{C})\right] &= \mathbf{E}\left[\sum_{S \in \mathcal{C}} c(S)\right] = \mathbf{E}\left[\sum_{S \in \mathcal{F}} \mathbf{1}_{S \in \mathcal{C}} \cdot c(S)\right] \\ &= \sum_{S \in \mathcal{F}} \mathbf{P}\left[S \in \mathcal{C}\right] \cdot c(S) = \sum_{S \in \mathcal{F}} \overline{y}(S) \cdot c(S). \end{split}$$

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Step 2: The probability for an element to be (not) covered

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• Step 2: The probability for an element to be (not) covered

$$P[x \notin \cup_{S \in C} S]$$

Lemma

Let $\mathcal{C} \subseteq \mathcal{F}$ be a random subset with each set S being included independently with probability $\overline{y}(S)$.

- The expected cost satisfies $\mathbf{E}[c(\mathcal{C})] = \sum_{S \in \mathcal{F}} c(S) \cdot \overline{y}(S)$.
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Proof:

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Step 2: The probability for an element to be (not) covered

$$\mathbf{P}[x \notin \cup_{S \in \mathcal{C}} S] = \prod_{S \in \mathcal{F}: x \in S} \mathbf{P}[S \notin \mathcal{C}]$$

Lemma

Let $\mathcal{C} \subseteq \mathcal{F}$ be a random subset with each set S being included independently with probability $\overline{y}(S)$.

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Step 2: The probability for an element to be (not) covered

$$\mathbf{P}[x \notin \cup_{S \in \mathcal{C}} S] = \prod_{S \in \mathcal{F}: x \in S} \mathbf{P}[S \notin \mathcal{C}] = \prod_{S \in \mathcal{F}: x \in S} (1 - \overline{y}(S))$$

Lemma

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$$1+x \leq e^x$$
 for any $x \in \mathbb{R}$

Lemma

Let $\mathcal{C} \subseteq \mathcal{F}$ be a random subset with each set \mathcal{S} being included independently with probability $\overline{y}(\mathcal{S})$.

- The expected cost satisfies $\mathbf{E}[c(\mathcal{C})] = \sum_{S \in \mathcal{F}} c(S) \cdot \overline{y}(S)$.
- The probability that x is covered satisfies $P[x \in \bigcup_{S \in C} S] \ge 1 \frac{1}{e}$.

Proof:

• Step 1: The expected cost of the random set $\mathcal{C} \checkmark$

$$\begin{split} \mathbf{E}\left[\,c(\mathcal{C})\,\right] &= \mathbf{E}\left[\,\sum_{S\in\mathcal{C}} c(S)\,\right] \\ &= \sum_{S\in\mathcal{F}} \mathbf{1}_{S\in\mathcal{C}} \cdot c(S)\,\right] \\ &= \sum_{S\in\mathcal{F}} \mathbf{P}\left[\,S\in\mathcal{C}\,\right] \cdot c(S) = \sum_{S\in\mathcal{F}} \overline{y}(S) \cdot c(S). \end{split}$$

Step 2: The probability for an element to be (not) covered

$$\mathbf{P}[x \notin \cup_{S \in \mathcal{C}} S] = \prod_{S \in \mathcal{F}: \ x \in S} \mathbf{P}[S \notin \mathcal{C}] = \prod_{S \in \mathcal{F}: \ x \in S} (1 - \overline{y}(S))$$

$$\leq \prod_{S \in \mathcal{F}: \ x \in S} e^{-\overline{y}(S)}$$

$$1 + x \leq e^{x} \text{ for any } x \in \mathbb{R}$$

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WEIGHTED SET COVER-LP(X, \mathcal{F}, c)

1: compute \overline{y}, an optimal solution to the linear program

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clearly runs in polynomial-time!

Theorem

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 - Linearity $\Rightarrow \mathbf{E}[c(\mathcal{C})] \leq 2\ln(n) \cdot \sum_{S \in \mathcal{F}} c(S) \cdot \overline{y}(S) \leq 2\ln(n) \cdot c(\mathcal{C}^*)$

Theorem

- With probability at least $1 \frac{1}{n}$, the returned set C is a valid cover of X.
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Proof:

- Step 1: The probability that C is a cover √
 - By previous Lemma, an element $x \in X$ is covered in one of the $2 \ln n$ iterations with probability at least $1 \frac{1}{e}$, so that

$$\mathbf{P}[x \not\in \cup_{S \in \mathcal{C}} S] \le \left(\frac{1}{e}\right)^{2 \ln n} = \frac{1}{n^2}.$$

$$\mathbf{P}[X = \cup_{S \in \mathcal{C}} S] = 1 - \mathbf{P}\left[\bigcup_{x \in X} \{x \notin \cup_{S \in \mathcal{C}} S\}\right]$$

$$\boxed{\mathbf{P}[A \cup B] \leq \mathbf{P}[A] + \mathbf{P}[B]} \geq 1 - \sum_{x \in X} \mathbf{P}[x \notin \bigcup_{S \in \mathcal{C}} S] \geq 1 - n \cdot \frac{1}{n^2} = 1 - \frac{1}{n}.$$

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[Exercise Question (9/10).10] gives a different perspective on the amplification procedure through non-linear randomised rounding.

Outline

Weighted Set Cover

MAX-CNF

MAX-CNF

Recall:

MAX-3-CNF Satisfiability ————

- Given: 3-CNF formula, e.g.: $(x_1 \lor x_3 \lor \overline{x_4}) \land (x_2 \lor \overline{x_3} \lor \overline{x_5}) \land \cdots$
- Goal: Find an assignment of the variables that satisfies as many clauses as possible.

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Why study this generalised problem?

- Allowing arbitrary clause lengths makes the problem more interesting (we will see that simply guessing is not the best!)
- a nice concluding example where we can practice previously learned approaches

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P[clause *i* is satisfied] =
$$1 - 2^{-\ell} := \alpha_{\ell}$$
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In particular, the guessing algorithm is a randomised 2-approximation.

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First statement as in the proof of Theorem 35.6. For clause i not to be satisfied, all ℓ occurring variables must be set to a specific value.

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Proof:

- First statement as in the proof of Theorem 35.6. For clause i not to be satisfied, all ℓ occurring variables must be set to a specific value.
- As before, let $Y := \sum_{i=1}^{m} Y_i$ be the number of satisfied clauses. Then,

$$\mathbf{E}[Y] = \mathbf{E}\left[\sum_{i=1}^{m} Y_i\right] = \sum_{i=1}^{m} \mathbf{E}[Y_i] \ge \sum_{i=1}^{m} \frac{1}{2} = \frac{1}{2} \cdot m. \quad \Box$$

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maximize $\sum_{i=1}^{m} z_{i}$ subject to $\sum_{i=1}^{m} v_{i} + \sum_{i=1}^{m} (1 - v_{i}) > z_{i} \quad \text{for each } i = 1, 2$

subject to
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- In the corresponding LP each ∈ {0, 1} is replaced by ∈ [0, 1]
- Let $(\overline{y}, \overline{z})$ be the optimal solution of the LP
- Obtain an integer solution y through randomised rounding of \overline{y}

Lemma

For any clause i of length ℓ ,

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• Assume w.l.o.g. all literals in clause i appear non-negated (otherwise replace every occurrence of x_i by $\overline{x_i}$ in the whole formula)

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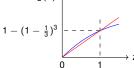
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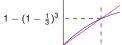
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Randomised Rounding yields a $1/(1-1/e)\approx 1.5820$ randomised approximation algorithm for MAX-CNF.

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- Then the expected number of satisfied clauses is:

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$$\boxed{\text{Since } (1 - 1/x)^x \le 1/e}$$

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Algorithm sets each variable x_i to TRUE with prob. $\frac{1}{2} \cdot \frac{1}{2} + \frac{1}{2} \cdot \overline{y}_i$. Note, however, that variables are **not** independently assigned!

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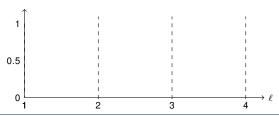
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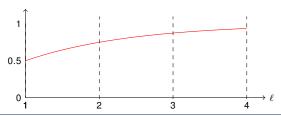
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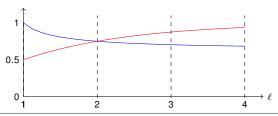
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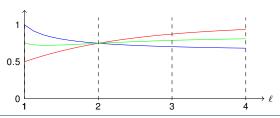
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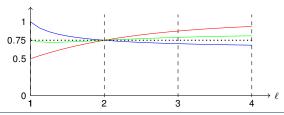
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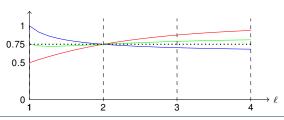
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MAX-CNF Conclusion

Summary

- Since $\alpha_2 = \beta_2 = 3/4$, we cannot achieve a better approximation ratio than 4/3 by combining Algorithm 1 & 2 in a different way
- The 4/3-approximation algorithm can be easily derandomised
 - Idea: use the conditional expectation trick for both Algorithm 1 & 2 and output the better solution
- The 4/3-approximation algorithm applies unchanged to a weighted version of MAX-CNF, where each clause has a non-negative weight
- Even MAX-2-CNF (every clause has length 2) is NP-hard!

Randomised Algorithms

Lecture 11: Spectral Graph Theory

Thomas Sauerwald (tms41@cam.ac.uk)

Lent 2025

Outline

Introduction to (Spectral) Graph Theory and Clustering

Matrices, Spectrum and Structure

A Simplified Clustering Problem



Source: Wikipedia

Seven Bridges at Königsberg 1737



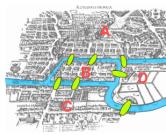
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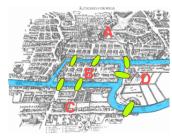
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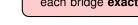


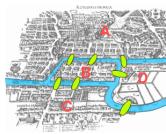




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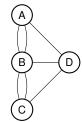
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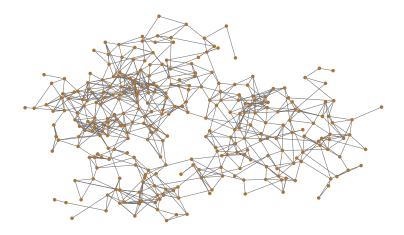


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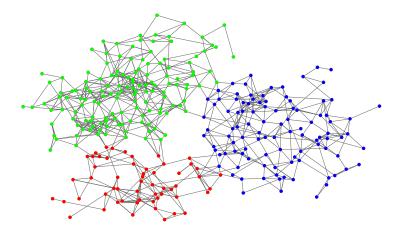
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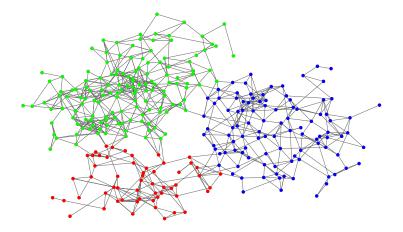
Graphs Nowadays: Clustering



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Graphs Nowadays: Clustering



Goal: Use spectrum of graphs (unstructured data) to extract clustering (communitites) or other structural information.

Graph Clustering (applications)

- Applications of Graph Clustering
 - Community detection
 - Group webpages according to their topics
 - Find proteins performing the same function within a cell
 - Image segmentation
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Graphs and Matrices

Graphs



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Graphs and Matrices

Graphs



- Connectivity
- Bipartiteness
- Number of triangles
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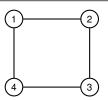
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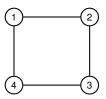
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Properties of A:

- The sum of elements in each row/column i equals the degree of the corresponding vertex i, deg(i)
- Since G is undirected, A is symmetric

Eigenvalues and Eigenvectors ——

Let $\mathbf{M} \in \mathbb{R}^{n \times n}$, $\lambda \in \mathbb{C}$ is an eigenvalue of \mathbf{M} if and only if there exists $x \in \mathbb{C}^n \setminus \{\mathbf{0}\}$ such that

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= orthogonal and normalised

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Remark: For symmetric matrices we have algebraic multiplicity = geometric multiplicity (otherwise >)

Example 1

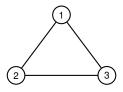


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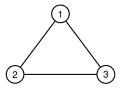
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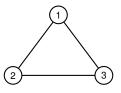


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Question: What are the Eigenvalues and Eigenvectors?



$$\mathbf{A} = \begin{pmatrix} 0 & 1 & 1 \\ 1 & 0 & 1 \\ 1 & 1 & 0 \end{pmatrix}$$

Solution:

- The three eigenvalues are $\lambda_1 = \lambda_2 = -1, \lambda_3 = 2$.
- The three eigenvectors are (for example):

$$f_1 = \begin{pmatrix} 1 \\ 0 \\ -1 \end{pmatrix}, \quad f_2 = \begin{pmatrix} -\frac{1}{2} \\ 1 \\ -\frac{1}{2} \end{pmatrix}, \quad f_3 = \begin{pmatrix} 1 \\ 1 \\ 1 \end{pmatrix}.$$

Laplacian Matrix —

Let G = (V, E) be a *d*-regular undirected graph. The (normalised) Laplacian matrix of G is the n by n matrix \mathbf{L} defined as

$$\mathbf{L} = \mathbf{I} - \frac{1}{d}\mathbf{A},$$

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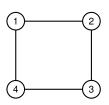
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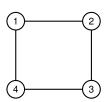


$$\boldsymbol{L} = \begin{pmatrix} 1 & -1/2 & 0 & -1/2 \\ -1/2 & 1 & -1/2 & 0 \\ 0 & -1/2 & 1 & -1/2 \\ -1/2 & 0 & -1/2 & 1 \end{pmatrix}$$

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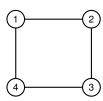
Properties of L:

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Properties of L:

- The sum of elements in each row/column equals zero
- L is symmetric

Relating Spectrum of Adjacency Matrix and Laplacian Matrix

- Correspondence between Adjacency and Laplacian Matrix -

A and L have the same set of eigenvectors.



Exercise: Prove this correspondence. Hint: Use that $\mathbf{L} = \mathbf{I} - \frac{1}{d}\mathbf{A}$. [Exercise 11/12.1]

Eigenvalues and eigenvectors -

Let $\mathbf{M} \in \mathbb{R}^{n \times n}$, $\lambda \in \mathbb{C}$ is an eigenvalue of \mathbf{M} if and only if there exists $x \in \mathbb{C}^n \setminus \{\mathbf{0}\}$ such that

$$\mathbf{M}\mathbf{x} = \lambda \mathbf{x}$$
.

We call x an eigenvector of **M** corresponding to the eigenvalue λ .

- Graph Spectrum -

Let **L** be the Laplacian matrix of a d-regular graph G with n vertices. Then, **L** has n real eigenvalues $\lambda_1 \leq \cdots \leq \lambda_n$ and n corresponding orthonormal eigenvectors f_1, \ldots, f_n . These eigenvalues associated with their multiplicities constitute the spectrum of G.

Useful Facts of Graph Spectrum

Lemma

Let **L** be the Laplacian matrix of an undirected, regular graph G = (V, E) with eigenvalues $\lambda_1 \leq \cdots \leq \lambda_n$.

- 1. $\lambda_1 = 0$ with eigenvector **1**
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The proof of these properties is based on a powerful characterisation of eigenvalues/vectors!

Courant-Fischer Min-Max Formula (non-examinable)

Let **M** be an *n* by *n* symmetric matrix with eigenvalues $\lambda_1 \leq \cdots \leq \lambda_n$. Then,

$$\lambda_k = \min_{\mathcal{S}: \ \dim(\mathcal{S}) = k} \max_{x \in \mathcal{S}, x \neq \mathbf{0}} \frac{x^T \mathbf{M} x}{x^T x},$$

where S is a subspace of \mathbb{R}^n . The eigenvectors corresponding to $\lambda_1,\ldots,\lambda_k$ minimise such expression.

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minimised by f2

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Quadratic Forms of the Laplacian

Lemma -

Let **L** be the Laplacian matrix of a *d*-regular graph G = (V, E) with n vertices. For any $x \in \mathbb{R}^n$,

$$x^T \mathbf{L} x = \sum_{\{u,v\} \in E} \frac{(x_u - x_v)^2}{d}.$$

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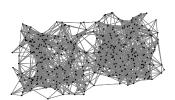
Proof:

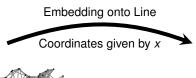
$$x^{T} \mathbf{L} x = x^{T} \left(\mathbf{I} - \frac{1}{d} \mathbf{A} \right) x = x^{T} x - \frac{1}{d} x^{T} \mathbf{A} x$$

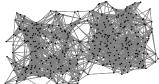
$$= \sum_{u \in V} x_{u}^{2} - \frac{2}{d} \sum_{\{u,v\} \in E} x_{u} x_{v}$$

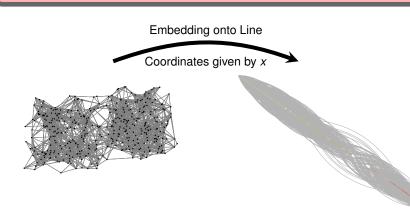
$$= \frac{1}{d} \sum_{\{u,v\} \in E} (x_{u}^{2} + x_{v}^{2} - 2x_{u} x_{v})$$

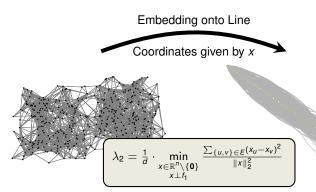
$$= \sum_{\{u,v\} \in E} \frac{(x_{u} - x_{v})^{2}}{d}.$$



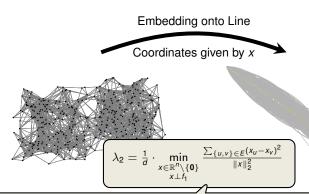








Question: How can we visualize a complicated object like an unknown graph with many vertices in low-dimensional space?



The coordinates in the vector **x** indicate how similar/dissimilar vertices are. Edges between dissimilar vertices are penalised quadratically.

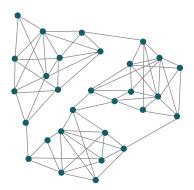
Outline

Introduction to (Spectral) Graph Theory and Clustering

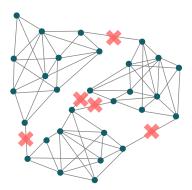
Matrices, Spectrum and Structure

A Simplified Clustering Problem

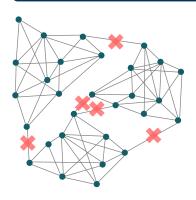
Partition the graph into **connected components** so that any pair of vertices in the same component is connected, but vertices in different components are not.

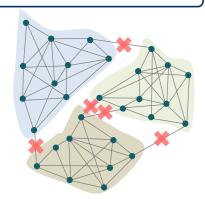


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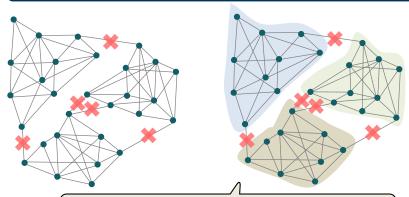


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Partition the graph into **connected components** so that any pair of vertices in the same component is connected, but vertices in different components are not.



We could obviously solve this easily using DFS/BFS, but let's see how we can tackle this using the spectrum of L!



 $\label{eq:Question: Question: What are the Eigenvectors with Eigenvalue 0 of L?}$



Question: What are the Eigenvectors with Eigenvalue 0 of L?





$\mathbf{A} =$	/0	1	1	0	0	0	0
	1	0	1	0	0	0	0
	1	1	0	0	0	0	0
A =	0	0	0	0	1	0	1
	0	0	0	1	0	1	0
	0	0	0	0	1	0	1
	0/	0	0	1	0	1	0,



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$$\mathbf{L} = \begin{pmatrix} 1 & -\frac{1}{2} & -\frac{1}{2} & 0 & 0 & 0 & 0 \\ -\frac{1}{2} & 1 & -\frac{1}{2} & 0 & 0 & 0 & 0 \\ -\frac{1}{2} & -\frac{1}{2} & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & -\frac{1}{2} & 0 & -\frac{1}{2} \\ 0 & 0 & 0 & -\frac{1}{2} & 1 & -\frac{1}{2} & 0 \\ 0 & 0 & 0 & 0 & -\frac{1}{2} & 1 & -\frac{1}{2} \\ 0 & 0 & 0 & 0 & -\frac{1}{2} & 0 & -\frac{1}{2} & 1 \end{pmatrix}$$



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Solution:

- Two smallest eigenvalues are $\lambda_1 = \lambda_2 = 0$.
- The corresponding two eigenvectors are:

$$f_1 = \begin{pmatrix} 1 \\ 1 \\ 1 \\ 0 \\ 0 \\ 0 \end{pmatrix}, \quad f_2 = \begin{pmatrix} 0 \\ 0 \\ 0 \\ 1 \\ 1 \\ 1 \end{pmatrix}$$



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Solution:

- The corresponding two eigenvectors are:

• Two smallest eigenvalues are $\lambda_1 = \lambda_2 = 0$. Thus we can easily solve the simplified clustering problem by computing the eigenvectors with eigenvalue 0

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Next Lecture: A fine-grained approach works even if the clusters are **sparsely** connected!

Let us generalise and formalise the previous example!

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Proof (multiplicity of 0 equals the no. of connected components):

1. (" \Longrightarrow " $cc(G) \le mult(0)$). We will show: G has exactly k connected comp. $C_1, \ldots, C_k \implies \lambda_1 = \cdots = \lambda_k = 0$

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- $\Rightarrow f_1, \dots, f_k$ constant on connected components
- **a** as f_1, \ldots, f_k are pairwise orthogonal, G must have k different connected components.

Let us generalise and formalise the previous example!

Proof (multiplicity of 0 equals the no. of connected components):

1. (" \Longrightarrow " $cc(G) \le mult(0)$). We will show:

G has exactly *k* connected comp.
$$C_1, \ldots, C_k \Rightarrow \lambda_1 = \cdots = \lambda_k = 0$$

- Take $\chi_{C_i} \in \{0,1\}^n$ such that $\chi_{C_i}(u) = \mathbf{1}_{u \in C_i}$ for all $u \in V$
- Clearly, the χ_{C_i} 's are orthogonal

2. (" \Leftarrow " $cc(G) \ge mult(0)$). We will show:

$$\lambda_1 = \cdots = \lambda_k = 0 \implies G$$
 has at least k connected comp. C_1, \ldots, C_k

- there exist f_1, \ldots, f_k orthonormal such that $\sum_{\{u,v\}\in E} (f_i(u) f_i(v))^2 = 0$
- $\Rightarrow f_1, \dots, f_k$ constant on connected components
- **a** as f_1, \ldots, f_k are pairwise orthogonal, G must have k different connected components.

Randomised Algorithms

Lecture 12: Spectral Graph Clustering

Thomas Sauerwald (tms41@cam.ac.uk)

Outline

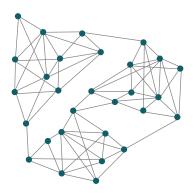
Conductance, Cheeger's Inequality and Spectral Clustering

Illustrations of Spectral Clustering and Extension to Non-Regular Graphs

Appendix: Relating Spectrum to Mixing Times (non-examinable)

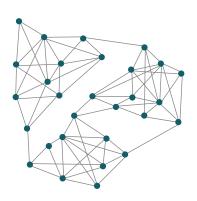
Graph Clustering

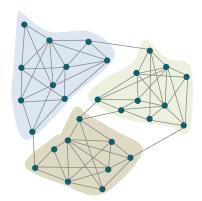
Partition the graph into **pieces (clusters)** so that vertices in the same piece have, on average, more connections among each other than with vertices in other clusters



Graph Clustering

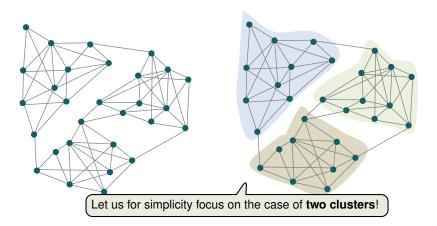
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Conductance

Let G = (V, E) be a d-regular and undirected graph and $\emptyset \neq S \subsetneq V$. The conductance (edge expansion) of S is

$$\phi(\mathcal{S}) := \frac{e(\mathcal{S}, \mathcal{S}^c)}{d \cdot |\mathcal{S}|}$$

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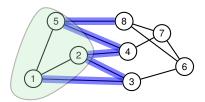
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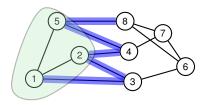


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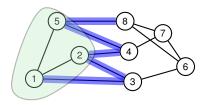
•
$$\phi(S) = ??$$

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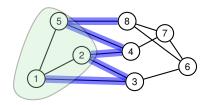
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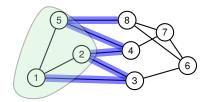
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- If G is a complete graph, then $e(S, V \setminus S) = |S| \cdot (n |S|)$ and $\phi(G) \approx 1/2$.

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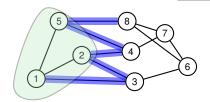
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Moreover, the conductance (edge expansion) of the graph G is

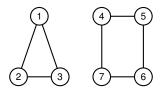
$$\phi(G) := \min_{S \subseteq V: \ 1 \le |S| \le n/2} \phi(S)$$

NP-hard to compute!

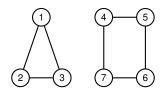


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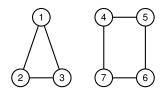
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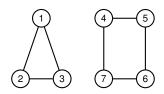
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$$\phi(G) = 0 \Leftrightarrow G \text{ is disconnected}$$



$$\phi(G) = 0 \Leftrightarrow G \text{ is disconnected } \Leftrightarrow \lambda_2(G) = 0$$



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What is the relationship between $\phi(G)$ and $\lambda_2(G)$ for **connected** graphs?

1D Grid (Path)

2D Grid

3D Grid





$$\lambda_2 \sim n^{-2}$$
 $\phi \sim n^{-1}$

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Hypercube







$$\lambda_2 \sim (\log n)^{-1}$$
$$\phi \sim (\log n)^{-1}$$



$$\lambda_2 = \Theta(1)$$
 $\phi = \Theta(1)$



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Cheeger's inequality -

Let G be a d-regular undirected graph and $\lambda_1 \leq \cdots \leq \lambda_n$ be the eigenvalues of its Laplacian matrix. Then,

$$\frac{\lambda_2}{2} \leq \phi(G) \leq \sqrt{2\lambda_2}.$$

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- no constant factor worst-case guarantee, but usually works well in practice (see examples later!)
- very fast: can be implemented in $O(|E| \log |E|)$ time

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By the Courant-Fischer Formula,

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Optimisation Problem: Embed vertices on a line such that sum of squared distances is minimised

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■ Let $S \subseteq V$ be the subset for which $\phi(G)$ is minimised. Define $y \in \mathbb{R}^n$ by:

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• Since $y \perp 1$, it follows that

$$\begin{split} \lambda_2 &\leq \frac{1}{d} \cdot \frac{\sum_{u \sim v} (y_u - y_v)^2}{\sum_u y_u^2} = \frac{1}{d} \cdot \frac{|E(S, V \setminus S)| \cdot \left(\frac{1}{|S|} + \frac{1}{|V \setminus S|}\right)^2}{\frac{1}{|S|} + \frac{1}{|V \setminus S|}} \\ &= \frac{1}{d} \cdot |E(S, V \setminus S)| \cdot \left(\frac{1}{|S|} + \frac{1}{|V \setminus S|}\right) \\ &\leq \frac{1}{d} \cdot \frac{2 \cdot |E(S, V \setminus S)|}{|S|} = 2 \cdot \phi(G). \quad \Box \end{split}$$

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$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 1 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix}$$

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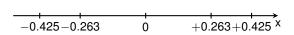
$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$

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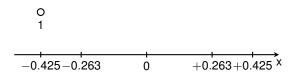
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & 1 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix} \quad \begin{pmatrix} 4 & 3 & 2 \\ 4 & 3 & 2 \\ 5 & 5 & 5 & 5 \\ 6 & 7 & 8 \end{pmatrix}$$

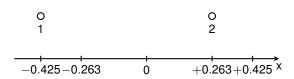
$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

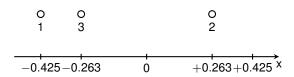
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

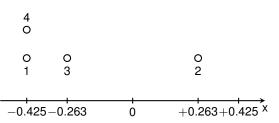
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

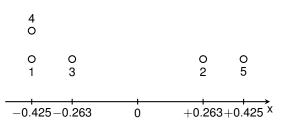
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

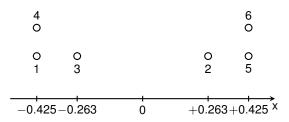
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

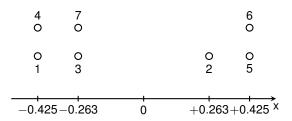
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

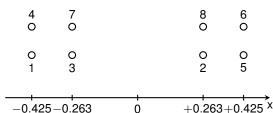
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & -\frac{1}{3} & 1 & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & 0 & 0 & 1 \\ 0 & 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

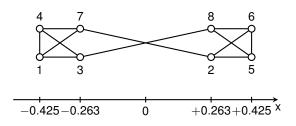
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 1 & 0 & 0 & 1 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & 1 \\ 0 & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

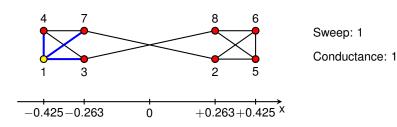
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 \end{pmatrix} \quad \mathbf{A}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

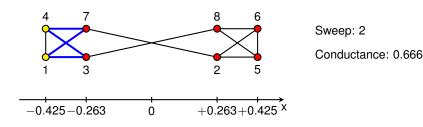
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

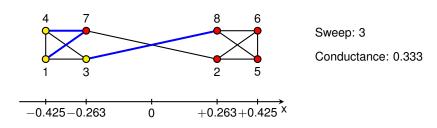
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

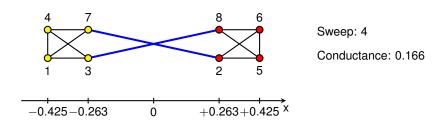
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 1 & 0 & 0 & 1 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & 1 \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & 1 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix} \quad \mathbf{A}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

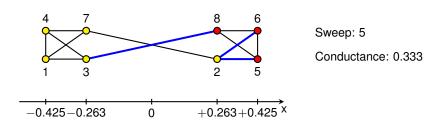
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 1 & 0 & 0 & 1 \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 & 0 \\ 0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 1 \end{pmatrix}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

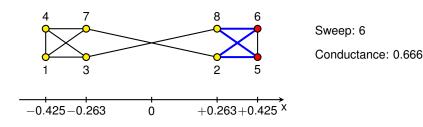
$$v = (-0.425, +0.263, -0.263, -0.425, +0.425, +0.425, -0.263, +0.263)^T$$



$$\mathbf{A} = \begin{pmatrix} 0 & 0 & 1 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 & 1 & 0 & 0 \end{pmatrix} \quad \mathbf{L} = \begin{pmatrix} 1 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & 1 & 0 & 0 & -\frac{1}{3} & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} \\ -\frac{1}{3} & 0 & 1 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & 1 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} \\ 0 & -\frac{1}{3} & 0 & 0 & 0 & -\frac{1}{3} & 1 & 0 & -\frac{1}{3} \\ 0 & 0 & -\frac{1}{3} & 0 & 0 & -\frac{1}{3} & 0 & 0 & 1 \end{pmatrix} \quad \mathbf{E}$$

$$\lambda_2 = 1 - \sqrt{5}/3 \approx 0.25$$

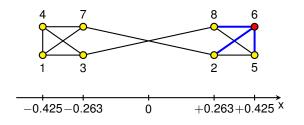
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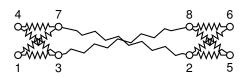


Sweep: 7

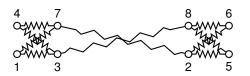
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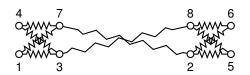


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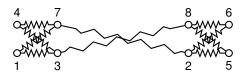
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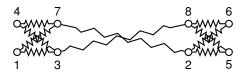
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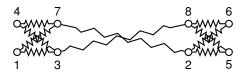
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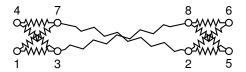


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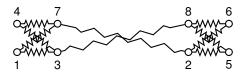


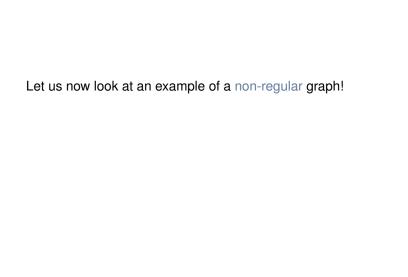
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The Laplacian Matrix (General Version)

The (normalised) Laplacian matrix of G = (V, E, w) is the n by n matrix

$$L = I - D^{-1/2}AD^{-1/2}$$

where **D** is a diagonal $n \times n$ matrix such that $\mathbf{D}_{uu} = deg(u) = \sum_{v : \{u,v\} \in \mathbf{E}} w(u,v)$, and **A** is the weighted adjacency matrix of G.

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- $\mathbf{L}_{uv} = -\frac{w(u,v)}{\sqrt{d_u d_v}}$ for $u \neq v$
- L is symmetric
- If G is d-regular, $\mathbf{L} = \mathbf{I} \frac{1}{d} \cdot \mathbf{A}$.

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Let G = (V, E, w) and $\emptyset \subsetneq S \subsetneq V$. The conductance (edge expansion) of S is

$$\phi(\mathcal{S}) := rac{\textit{w}(\mathcal{S}, \mathcal{S}^c)}{\min\{\mathsf{vol}(\mathcal{S}), \mathsf{vol}(\mathcal{S}^c)\}},$$

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Stochastic Block Model and 1D-Embedding

Stochastic Block Model

$$\textit{G} = (\textit{V}, \textit{E})$$
 with clusters $\textit{S}_1, \textit{S}_2 \subseteq \textit{V}, 0 \leq \textit{q} < \textit{p} \leq 1$

$$\mathbf{P}[\{u,v\} \in E] = \begin{cases} p & \text{if } u,v \in S_i, \\ q & \text{if } u \in S_i,v \in S_j, i \neq j. \end{cases}$$

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Number of Edges: 919
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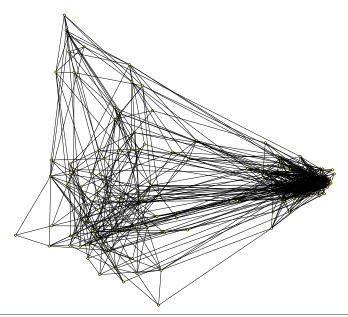
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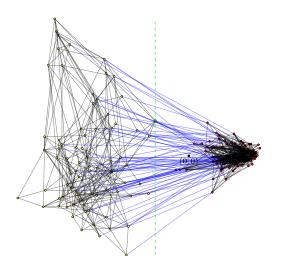


Drawing the 2D-Embedding



Spectral Clustering

Best Solution found by Spectral Clustering

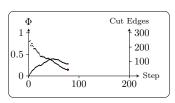


• Step: 78

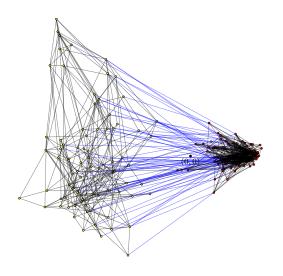
Threshold: -0.0336
Partition Sizes: 78/122

• Cut Edges: 84

• Conductance: 0.1448



Clustering induced by Blocks



• Step: -

• Threshold: -

• Partition Sizes: 80/120

• Cut Edges: 88

• Conductance: 0.1486

Graph
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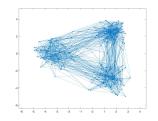
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$$|V| = 300, |S_i| = 100$$

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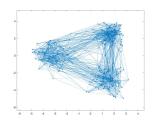


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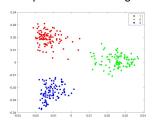
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Spectral embedding

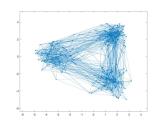


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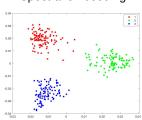
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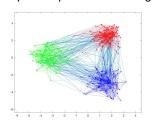
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Spectral embedding



Output of Spectral Clustering



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 - small λ_k means there exist k sparsely connected subsets in the graph (recall: $\lambda_1 = \ldots = \lambda_k = 0$ means there are k connected components)

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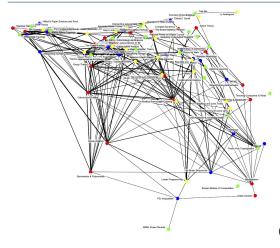
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- In the former example $\lambda = \{0, 0.15, 0.37, 0.40, 0.43, \dots\} \implies k = 2$.
- For k = 2 use sweep-cut extract clusters. For k ≥ 3 use embedding in k-dimensional space and apply k-means (geometric clustering)

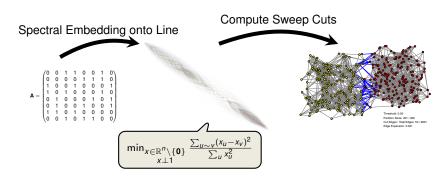
Another Example



(many thanks to Kalina Jasinska)

- nodes represent math topics taught within 4 weeks of a Mathcamp
- node colours represent to the week in which they thought
- teachers were asked to assign weights in 0 10 indicating how closely related two classes are

Summary: Spectral Clustering



- Given any graph (adjacency matrix)
- Graph Spectrum (computable in poly-time)
 - λ_2 (relates to connectivity)
 - λ_n (relates to bipartiteness)

- Cheeger's Inequality

 - relates λ_2 to conductance unbounded approximation ratio

 - effective in practice

Outline

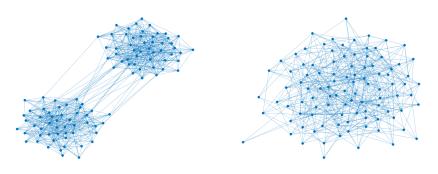
Conductance, Cheeger's Inequality and Spectral Clustering

Illustrations of Spectral Clustering and Extension to Non-Regular Graphs

Appendix: Relating Spectrum to Mixing Times (non-examinable)

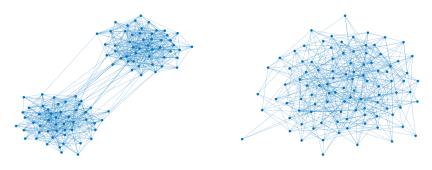
Relation between Clustering and Mixing (non-examinable)

Which graph has a "cluster-structure"?



Relation between Clustering and Mixing (non-examinable)

- Which graph has a "cluster-structure"?
- Which graph mixes faster?



Recall: If the underlying graph G is connected, undirected and d-regular, then the random walk converges towards the stationary distribution $\pi = (1/n, \dots, 1/n)$, which satisfies $\pi P = \pi$.

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Lemma

Consider a lazy random walk on a connected, undirected and d-regular graph. Then for any initial distribution x,

$$\left\| \mathbf{x}\mathbf{P}^{t} - \pi \right\|_{2} \leq \lambda^{t},$$

with $1 = \lambda_1 > \lambda_2 \ge \cdots \ge \lambda_n$ as eigenvalues and $\lambda := \max\{|\lambda_2|, |\lambda_n|\}$.

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Lemma

Consider a lazy random walk on a connected, undirected and d-regular graph. Then for any initial distribution x,

$$\left\| x \mathbf{P}^t - \pi \right\|_2 \le \lambda^t,$$

with $1 = \lambda_1 > \lambda_2 \ge \cdots \ge \lambda_n$ as eigenvalues and $\lambda := \max\{|\lambda_2|, |\lambda_n|\}$. \Rightarrow This implies for $t = \mathcal{O}(\frac{\log n}{\log(1/\lambda)}) = \mathcal{O}(\frac{\log n}{1-\lambda})$,

 $\left\|x\mathbf{P}^t - \pi\right\|_{t_t} \leq \frac{1}{4}.$

due to laziness, $\lambda_n \geq 0$

Express x in terms of the orthonormal basis of **P**, $v_1 = \pi, v_2, \dots, v_n$:

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$$||x\mathbf{P} - \pi||_2^2$$

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$$\|x\mathbf{P} - \pi\|_{2}^{2} = \left\|\left(\sum_{i=1}^{n} \alpha_{i} v_{i}\right) \mathbf{P} - \pi\right\|_{2}^{2}$$

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$$x = \sum_{i=1}^{n} \alpha_i v_i.$$

$$\Rightarrow \|x\mathbf{P} - \pi\|_{2}^{2} = \left\| \left(\sum_{i=1}^{n} \alpha_{i} \mathbf{v}_{i} \right) \mathbf{P} - \pi \right\|_{2}^{2}$$
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• Since x is a probability vector and all $v_i \ge 2$ are orthogonal to π , $\alpha_1 = 1$.

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$$x = \sum_{i=1}^{n} \alpha_i v_i.$$

• Since x is a probability vector and all $v_i \ge 2$ are orthogonal to π , $\alpha_1 = 1$.

$$\|\mathbf{x}\mathbf{P} - \pi\|_{2}^{2} = \left\| \left(\sum_{i=1}^{n} \alpha_{i} v_{i} \right) \mathbf{P} - \pi \right\|_{2}^{2}$$

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 \Rightarrow

• Express x in terms of the orthonormal basis of P, $v_1 = \pi, v_2, \dots, v_n$:

$$x=\sum_{i=1}^n\alpha_iv_i.$$

Since
$$x$$
 is a probability vector and all $v_i \ge 2$ are orthogonal to π , $\alpha_1 = 1$.

$$\|x\mathbf{P} - \pi\|_2^2 = \left\|\left(\sum_{i=1}^n \alpha_i v_i\right)\mathbf{P} - \pi\right\|_2^2$$

$$= \left\|\pi + \sum_{i=2}^n \alpha_i \lambda_i v_i - \pi\right\|_2^2$$
Since the v_i 's are orthogonal
$$= \sum_{i=2}^n \|\alpha_i \lambda_i v_i\|_2^2$$

$$\leq \lambda^2 \sum_{i=2}^n \|\alpha_i v_i\|_2^2 = \lambda^2 \left\|\sum_{i=2}^n \alpha_i v_i\right\|_2^2 = \lambda^2 \|x - \pi\|_2^2$$

Express x in terms of the orthonormal basis of **P**, $v_1 = \pi, v_2, \dots, v_n$:

$$x = \sum_{i=1}^{n} \alpha_i v_i.$$

• Since x is a probability vector and all $v_i \ge 2$ are orthogonal to π , $\alpha_1 = 1$.

$$\Rightarrow$$

$$\|x\mathbf{P} - \pi\|_{2}^{2} = \left\|\left(\sum_{i=1}^{n} \alpha_{i} v_{i}\right) \mathbf{P} - \pi\right\|_{2}^{2}$$

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■ Hence $||x\mathbf{P}^t - \pi||_2^2$

Express *x* in terms of the orthonormal basis of **P**, $v_1 = \pi, v_2, \dots, v_n$:

$$x = \sum_{i=1}^{n} \alpha_i v_i.$$

• Since x is a probability vector and all $v_i \ge 2$ are orthogonal to π , $\alpha_1 = 1$.

$$\Rightarrow$$

$$\|\mathbf{x}\mathbf{P} - \pi\|_{2}^{2} = \left\| \left(\sum_{i=1}^{n} \alpha_{i} \mathbf{v}_{i} \right) \mathbf{P} - \pi \right\|_{2}^{2}$$
$$= \left\| \pi + \sum_{i=2}^{n} \alpha_{i} \lambda_{i} \mathbf{v}_{i} - \pi \right\|_{2}^{2}$$

since the
$$v_i$$
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$$= \left\| \sum_{i=2}^n \alpha_i \lambda_i v_i \right\|_2^2$$

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$$\leq \lambda^2 \sum_{i=2}^n \|\alpha_i v_i\|_2^2 = \lambda^2 \left\| \sum_{i=2}^n \alpha_i v_i \right\|_2^2 = \lambda^2 \|x - \pi\|_2^2$$

• Hence $||x\mathbf{P}^t - \pi||_2^2 \le \lambda^{2t} \cdot ||x - \pi||_2^2$

Express x in terms of the orthonormal basis of **P**, $v_1 = \pi, v_2, \dots, v_n$:

$$x = \sum_{i=1}^{n} \alpha_i v_i.$$

$$\Rightarrow$$

$$\|x\mathbf{P} - \pi\|_{2}^{2} = \left\| \left(\sum_{i=1}^{n} \alpha_{i} v_{i} \right) \mathbf{P} - \pi \right\|_{2}^{2}$$
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$$\| \text{Hence } \|x\mathbf{P}^t - \pi\|_2^2 \leq \lambda^{2t} \cdot \|x - \pi\|_2^2 \leq \lambda^{2t} \cdot 1.$$

$$\|x - \pi\|_2^2 + \|\pi\|_2^2 = \|x\|_2^2 \leq 1$$

Some References on Spectral Graph Theory and Clustering



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Thank you and Best Wishes for the Exam!

I'm very interested to hear your feedback about the slides and the course more generally. You can use the student feedback form or send me an email during or after the course (tms41@cam.ac.uk).