## SS 2017

# Efficient Algorithms and Data Structures II

Harald Räcke

Fakultät für Informatik TU München

http://www14.in.tum.de/lehre/2017SS/ea/

Summer Term 2017



# **Organizational Matters**



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### Modul: IN2004

- Name: "Efficient Algorithms and Data Structures II" "Effiziente Algorithmen und Datenstrukturen II"
- ECTS: 8 Credit points
- Lectures:
  - ► 4 SWS

Wed 12:15-13:45 (Room 00.13.009A) Fri 10:15-11:45 (MS HS3)

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## **The Lecturer**

- Harald Räcke
- Email: raecke@in.tum.de
- Room: 03.09.044
- Office hours: (per appointment)



## **Tutorials**

#### Tutor:

- Richard Stotz
- stotz@tum.de
- Room: 03.09.057
- per appointment
- Room: 03.11.018
- Time: Wed 16:00-17:30



In order to pass the module you need to pass an exam.

#### Exam:

- 2.5 hours
- Date will be announced shortly.
- There are no resources allowed, apart from a hand-written piece of paper (A4).
- Answers should be given in English, but German is also accepted.



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- An assignment sheet is usually made available on Wednesday on the module webpage.
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- You can hand in your solutions by putting them in the right folder in front of room 03.09.020.
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### **1** Contents

## Part 1: Linear Programming

# Part 2: Approximation Algorithms



1 Contents

# 2 Literatur



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Skript Optimierung, 1996

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Approximation Algorithms,

Springer 2001



David P. Williamson and David B. Shmoys: The Design of Approximation Algorithms, Cambridge University Press 2011

G. Ausiello, P. Crescenzi, G. Gambosi, V. Kann, A. Marchetti-Spaccamela, and M. Protasi: *Complexity and Approximation*, Springer, 1999



# **Linear Programming**



#### Brewery brews ale and beer.

- Production limited by supply of corn, hops and barley malt
- Recipes for ale and beer require different amounts of resources



**3 Introduction to Linear Programming** 

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	Corn (kg)	Hops (kg)	Malt (kg)	Profit (€)
ale (barrel)	5	4	35	13
beer (barrel)	15	4	20	23
supply	480	160	1190	



3 Introduction to Linear Programming

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ale (barrel)	5	4	35	13
beer (barrel)	15	4	20	23
supply	480	160	1190	

#### How can brewer maximize profits?

only brew been 32 barrels of been
 barrels ale, 255 barrels been
 barrels ale, 26 barrels been



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ale (barrel)	5	4	35	13
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supply	480	160	1190	

#### How can brewer maximize profits?

- only brew ale: 34 barrels of ale
- only brew beer: 32 barrels of beer
- 7.5 barrels ale, 29.5 barrels beer
- 12 barrels ale, 28 barrels beer





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#### How can brewer maximize profits?

- ▶ only brew ale: 34 barrels of ale  $\Rightarrow$  442 €
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 $\Rightarrow$  //b\*

 $\rightarrow$  800  $\in$ 



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- ▶ only brew ale: 34 barrels of ale  $\Rightarrow$  442 €
- only brew beer: 32 barrels of beer  $\Rightarrow$  736  $\in$
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FADS II Harald Räcke

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- ⇒ 736€
  - ⇒ 776€
    - ⇒ 800€



### Linear Program

- Introduce subdivision and dothat define how much ale and beer to produce.
- Choose the variables in such a way that the (profit) is maximized.
- Make: sure that no consistent (due to limited supply) are violated.



### Linear Program

- Introduce variables a and b that define how much ale and beer to produce.
- Choose the variables in such a way that the objective function (profit) is maximized.
- Make sure that no constraints (due to limited supply) are violated.



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max	13a	+	23 <i>b</i>
s.t.	5a	+	$15b \leq 480$
	4 <i>a</i>	+	$4b \leq 160$
	35a	+	$20b \leq 1190$
			$a,b \geq 0$



### LP in standard form:

- output: numbers >
- #decision variables, m = #constraints
- maximize linear objective function subject to linear (in)equalities





### LP in standard form:

- input: numbers  $a_{ij}$ ,  $c_j$ ,  $b_i$
- output: numbers x<sub>j</sub>
- ▶ n = #decision variables, m = #constraints
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$$\begin{array}{|c|c|c|c|c|} \max & \sum_{j=1}^{n} c_{j} x_{j} \\ \text{s.t.} & \sum_{j=1}^{n} a_{ij} x_{j} &= b_{i} \ 1 \leq i \leq m \\ & x_{j} \geq 0 \ 1 \leq j \leq n \end{array} \qquad \begin{array}{|c|c|c|c|c|c|} \max & c^{T} x \\ \text{s.t.} & Ax &= b \\ & x \geq 0 \\ & x \geq 0 \end{array}$$



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$$\max \sum_{\substack{j=1\\n}}^{n} c_j x_j$$
  
s.t. 
$$\sum_{\substack{j=1\\j=1}}^{n} a_{ij} x_j = b_i \quad 1 \le i \le m$$
$$x_j \ge 0 \quad 1 \le j \le n$$

$$\begin{array}{rcl} \max & c^T x \\ \text{s.t.} & Ax &= b \\ & x &\geq 0 \end{array}$$



### LP in standard form:

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### **Original LP**

max	13a	+	23 <i>b</i>	
s.t.	5a	+	15 <i>b</i>	≤ 480
	4 <i>a</i>	+	4b	$\leq 160$
	35a	+	20 <i>b</i>	$\leq 1190$
			a,b	≥ 0

**Standard Form** 

Add a slack variable to every constraint.





### **Original LP**

max	13a	+	23 <i>b</i>	
s.t.	5a	+	15b	≤ 480
	4 <i>a</i>	+	4b	$\leq 160$
	35a	+	20b	$\leq 1190$
			a,b	$\geq 0$

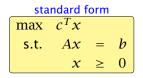
#### **Standard Form**

Add a slack variable to every constraint.

max	13a	+	23 <i>b</i>							
s.t.	5a	+	15 <i>b</i>	+	$S_C$					= 480
	4 <i>a</i>	+	4b			+	Sh			= 160
	35a	+	20 <i>b</i>					+	$S_m$	= 1190
	а	,	b	,	$S_C$	,	s <sub>h</sub>	,	Sm	≥ 0



There are different standard forms:

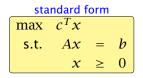








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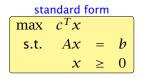


min	$c^T x$		
s.t.	Ax	=	b
	X	≥	0



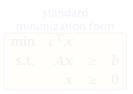


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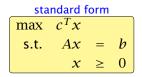
standard						
maximization form						
max	$c^T x$					
s.t.	Ax	$\leq$	b			
	x	$\geq$	0			

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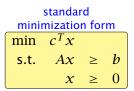


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standard					
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	X	$\geq$	0		

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s.t.	Ax	=	b
	X	$\geq$	0





It is easy to transform variants of LPs into (any) standard form:

greater or equal to equality:

min to max:



It is easy to transform variants of LPs into (any) standard form:

less or equal to equality:



min to max:

nina – 36 – 50 – – **max** – a + 36 – 50



It is easy to transform variants of LPs into (any) standard form:

less or equal to equality:

$$a - 3b + 5c \le 12 \implies a - 3b + 5c + s = 12$$
  
 $s \ge 0$ 

greater or equal to equality:

min to max:

a + a − 36 + 5c − − 26 + 36 − 5c



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greater or equal to equality:

 $a - 3b + 5c \ge 12 \implies \frac{a - 3b + 5c - s = 12}{s \ge 0}$ 

min to max:



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ai∂—db+5a — x**sm** <--- ai∂+db—5a



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min to max:

 $\min a - 3b + 5c \implies \max -a + 3b - 5c$ 



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It is easy to transform variants of LPs into (any) standard form:

equality to less or equal:

 $a - 3b + 5c = 12 \implies a - 3b + 5c \le 12$  $-a + 3b - 5c \le -12$ 

equality to greater or equal:

unrestricted to nonnegative:



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It is easy to transform variants of LPs into (any) standard form:

equality to less or equal:

$$a - 3b + 5c = 12 \implies \frac{a - 3b + 5c \le 12}{-a + 3b - 5c \le -12}$$

#### equality to greater or equal:

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$$-a + 3b - 5c \ge -12$$

unrestricted to nonnegative:



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unrestricted to nonnegative:

x unrestricted  $\Rightarrow x = x^+ - x^-, x^+ \ge 0, x^- \ge 0$ 



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#### **Observations:**

- a linear program does not contain  $x^2$ ,  $\cos(x)$ , etc.
- transformations between standard forms can be done efficiently and only change the size of the LP by a small constant factor
- for the standard minimization or maximization LPs we could include the nonnegativity constraints into the set of ordinary constraints; this is of course not possible for the standard form



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#### Definition 1 (Linear Programming Problem (LP))

Let  $A \in \mathbb{Q}^{m \times n}$ ,  $b \in \mathbb{Q}^m$ ,  $c \in \mathbb{Q}^n$ ,  $\alpha \in \mathbb{Q}$ . Does there exist  $x \in \mathbb{Q}^n$  s.t. Ax = b,  $x \ge 0$ ,  $c^T x \ge \alpha$ ?

**Questions**:

- Is LP in NP?
- Is LP in co-NP?
- Is LP in P?

Input size:



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Let  $A \in \mathbb{Q}^{m \times n}$ ,  $b \in \mathbb{Q}^m$ ,  $c \in \mathbb{Q}^n$ ,  $\alpha \in \mathbb{Q}$ . Does there exist  $x \in \mathbb{Q}^n$  s.t. Ax = b,  $x \ge 0$ ,  $c^T x \ge \alpha$ ?

#### **Questions**:

- Is LP in NP?
- Is LP in co-NP?
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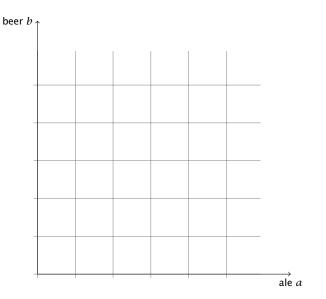
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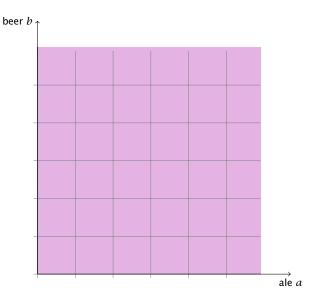
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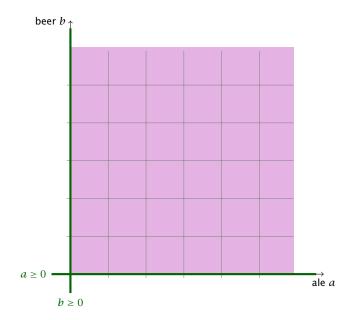
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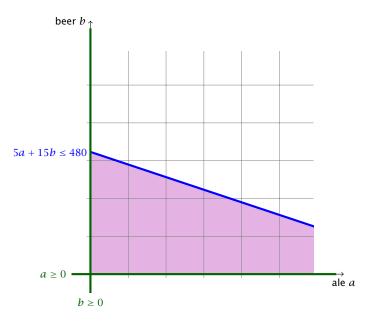
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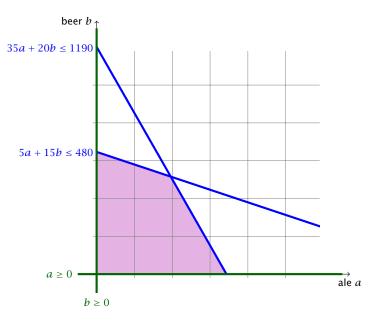


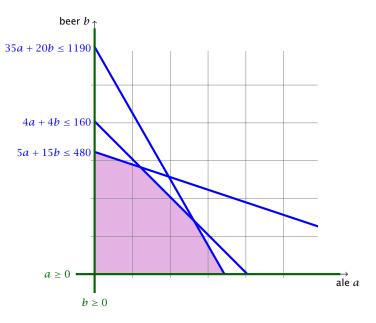


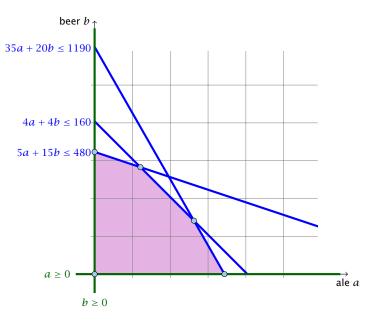


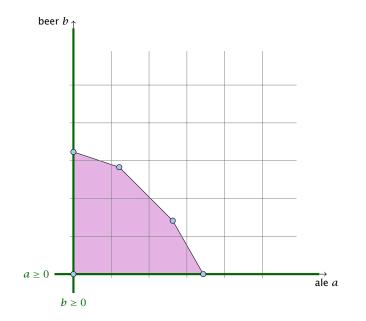


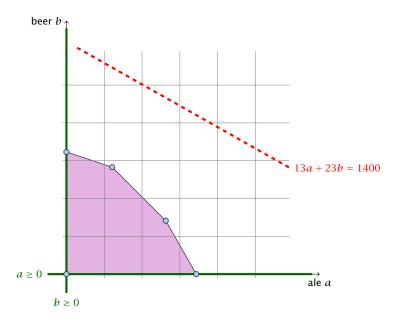


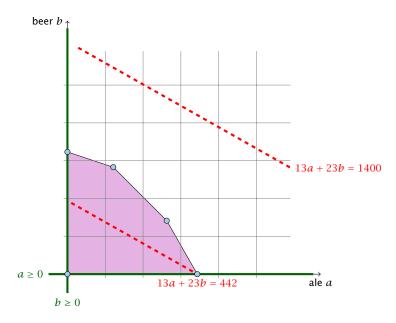


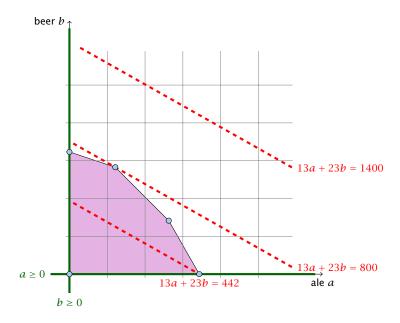


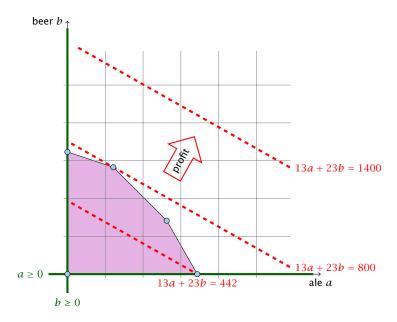


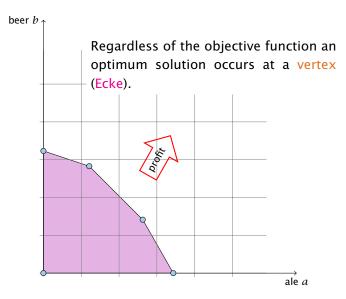












Let for a Linear Program in standard form  $P = \{x \mid Ax = b, x \ge 0\}.$ 

Is called the second constraints (Losungsraum) of the LR A point second is called a second constraint (gültige Losung). If A constraints the LP is called Second (gültige Losung).

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Given vectors/points  $x_1, \ldots, x_k \in \mathbb{R}^n$ ,  $\sum \lambda_i x_i$  is called

- linear combination if  $\lambda_i \in \mathbb{R}$ .
- affine combination if  $\lambda_i \in \mathbb{R}$  and  $\sum_i \lambda_i = 1$ .
- convex combination if  $\lambda_i \in \mathbb{R}$  and  $\sum_i \lambda_i = 1$  and  $\lambda_i \ge 0$ .
- conic combination if  $\lambda_i \in \mathbb{R}$  and  $\lambda_i \ge 0$ .

Note that a combination involves only finitely many vectors.



A set  $X \subseteq \mathbb{R}^n$  is called

- a linear subspace if it is closed under linear combinations.
- an affine subspace if it is closed under affine combinations.
- convex if it is closed under convex combinations.
- a convex cone if it is closed under conic combinations.

Note that an affine subspace is **not** a vector space



Given a set  $X \subseteq \mathbb{R}^n$ .

- span(X) is the set of all linear combinations of X (linear hull, span)
- aff(X) is the set of all affine combinations of X (affine hull)
- conv(X) is the set of all convex combinations of X (convex hull)
- cone(X) is the set of all conic combinations of X (conic hull)



## **Definition 5** A function $f : \mathbb{R}^n \to \mathbb{R}$ is convex if for $x, y \in \mathbb{R}^n$ and $\lambda \in [0, 1]$ we have

 $f(\lambda x + (1-\lambda)\gamma) \leq \lambda f(x) + (1-\lambda)f(\gamma)$ 

**Lemma 6** If  $P \subseteq \mathbb{R}^n$ , and  $f : \mathbb{R}^n \to \mathbb{R}$  convex then also

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**3 Introduction to Linear Programming** 

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#### **Dimensions**

#### **Definition 7**

The dimension dim(*A*) of an affine subspace  $A \subseteq \mathbb{R}^n$  is the dimension of the vector space  $\{x - a \mid x \in A\}$ , where  $a \in A$ .

#### **Definition 8**

The dimension  $\dim(X)$  of a convex set  $X \subseteq \mathbb{R}^n$  is the dimension of its affine hull  $\operatorname{aff}(X)$ .



#### **Definition 9** A set $H \subseteq \mathbb{R}^n$ is a hyperplane if $H = \{x \mid a^T x = b\}$ , for $a \neq 0$ .

## **Definition 10** A set $H' \subseteq \mathbb{R}^n$ is a (closed) halfspace if $H = \{x \mid a^T x \leq b\}$ , for $a \neq 0$ .



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#### **Definition 11**

A polytop is a set  $P \subseteq \mathbb{R}^n$  that is the convex hull of a finite set of points, i.e., P = conv(X) where |X| = c.



### **Definition 12**

A polyhedron is a set  $P \subseteq \mathbb{R}^n$  that can be represented as the intersection of finitely many half-spaces  $\{H(a_1, b_1), \ldots, H(a_m, b_m)\}$ , where

 $H(a_i, b_i) = \{x \in \mathbb{R}^n \mid a_i x \le b_i\} .$ 

### **Definition** 13

A polyhedron *P* is bounded if there exists *B* s.t.  $||x||_2 \le B$  for all  $x \in P$ .



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#### Theorem 14

P is a bounded polyhedron iff P is a polytop.



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# **Definition 15** Let $P \subseteq \mathbb{R}^n$ , $a \in \mathbb{R}^n$ and $b \in \mathbb{R}$ . The hyperplane

 $H(a,b) = \{x \in \mathbb{R}^n \mid a^T x = b\}$ 

## is a supporting hyperplane of *P* if $\max\{a^T x \mid x \in P\} = b$ .

### **Definition 16**

Let  $P \subseteq \mathbb{R}^n$ . F is a face of P if F = P or  $F = P \cap H$  for some supporting hyperplane H.

### **Definition 17**

Let  $P \subseteq \mathbb{R}^n$ .

- a face v is a vertex of P if  $\{v\}$  is a face of P.
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### Equivalent definition for vertex:

### **Definition 18**

Given polyhedron *P*. A point  $x \in P$  is a vertex if  $\exists c \in \mathbb{R}^n$  such that  $c^T y < c^T x$ , for all  $y \in P$ ,  $y \neq x$ .

#### **Definition 19**

Given polyhedron *P*. A point  $x \in P$  is an extreme point if  $\nexists a, b \neq x, a, b \in P$ , with  $\lambda a + (1 - \lambda)b = x$  for  $\lambda \in [0, 1]$ .

#### Lemma 20

A vertex is also an extreme point.



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#### Observation

The feasible region of an LP is a Polyhedron.



#### Theorem 21

## *If there exists an optimal solution to an LP (in standard form) then there exists an optimum solution that is an extreme point.*

- suppose of is optimal solution that is not extreme point:
- there exists direction d = 0 such that a = 0
- $\wedge dd = 0$  because  $\wedge (x + d) = b$
- Wlog. assume  $c^{(j)} (2 \leq 0$  (by taking either c or -d)
- $\sim$  Consider  $\sim + \lambda d_{\mu} \lambda > 0$



### Theorem 21

If there exists an optimal solution to an LP (in standard form) then there exists an optimum solution that is an extreme point.

- suppose x is optimal solution that is not extreme point
- there exists direction  $d \neq 0$  such that  $x \pm d \in P$
- Ad = 0 because  $A(x \pm d) = b$
- Wlog. assume  $c^T d \ge 0$  (by taking either d or -d)
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**Case 1.**  $[\exists j \text{ s.t. } d_j < 0]$ 

- increase  $\lambda$  to  $\lambda$  until first component of  $\infty$  -- $\lambda/\ell$  hits 0
- $\sim \sim -\lambda/d$  is feasible. Since  $\lambda(z + \lambda/d) = b$  and  $z + \lambda/d = 0$

**Case 2.**  $[d_j \ge 0$  for all j and  $c^T d > 0$ ]

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**Case 1.**  $[\exists j \text{ s.t. } d_j < 0]$ 

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3 Introduction to Linear Programming

**Case 1.**  $[\exists j \text{ s.t. } d_j < 0]$ 

### • increase $\lambda$ to $\lambda'$ until first component of $x + \lambda d$ hits 0

- $x + \lambda' d$  is feasible. Since  $A(x + \lambda' d) = b$  and  $x + \lambda' d \ge 0$
- $x + \lambda' d$  has one more zero-component ( $d_k = 0$  for  $x_k = 0$  as  $x \pm d \in P$ )
- $c^T x' = c^T (x + \lambda' d) = c^T x + \lambda' c^T d \ge c^T x$

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- $x + \lambda d$  is feasible for all  $\lambda \ge 0$  since  $A(x + \lambda d) = b$  and  $x + \lambda d \ge x \ge 0$
- as  $\lambda \to \infty$ ,  $c^T(x + \lambda d) \to \infty$  as  $c^T d > 0$



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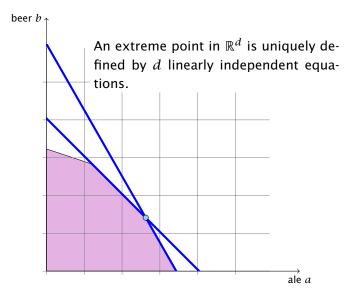
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- as  $\lambda \to \infty$ ,  $c^T(x + \lambda d) \to \infty$  as  $c^T d > 0$



# **Algebraic View**



#### Notation

Suppose  $B \subseteq \{1 \dots n\}$  is a set of column-indices. Define  $A_B$  as the subset of columns of A indexed by B.

**Theorem 22** Let  $P = \{x \mid Ax = b, x \ge 0\}$ . For  $x \in P$ , define  $B = \{j \mid x_j > 0\}$ . Then x is extreme point iff  $A_B$  has linearly independent columns.



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- has linearly dependent columns as the 0
- $d_{ij} = 0$  for all j with  $c_{ij} = 0$  as  $c = d \ge 0$
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- there exists direction d s.t.  $x \pm d \in P$
- Ad = 0 because  $A(x \pm d) = b$
- define  $B' = \{j \mid d_j \neq 0\}$
- $A_{B'}$  has linearly dependent columns as Ad = 0
- $d_j = 0$  for all j with  $x_j = 0$  as  $x \pm d \ge 0$
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- extend if to 80° by adding 0-components
- $\sim$  now,  $Ad \simeq 0$  and  $d_{1} \simeq 0$  whenever  $\alpha_{1} \simeq 0$
- for sufficiently small  $\lambda$  we have  $\infty = \lambda d = 2^{n}$
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- assume A<sub>B</sub> has linearly dependent columns
- there exists  $d \neq 0$  such that  $A_B d = 0$
- extend d to  $\mathbb{R}^n$  by adding 0-components
- now, Ad = 0 and  $d_j = 0$  whenever  $x_j = 0$
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Let  $P = \{x \mid Ax = b, x \ge 0\}$ . For  $x \in P$ , define  $B = \{j \mid x_j > 0\}$ . If  $A_B$  has linearly independent columns then x is a vertex of P.

• define 
$$c_j = \begin{cases} 0 & j \in B \\ -1 & j \notin B \end{cases}$$

• then  $c^T x = 0$  and  $c^T y \le 0$  for  $y \in P$ 

- assume  $c^T y = 0$ ; then  $y_j = 0$  for all  $j \notin B$
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- assume that mode(A) = m
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- if now  $b_1 = \sum_{i=1}^{n} b_i > b_i$  then
- then the LP is infeasible, since for all is the LP is infeasible, since for all is that fulfill constraints device we have

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# From now on we will always assume that the constraint matrix of a standard form LP has full row rank.



**3 Introduction to Linear Programming** 

Given  $P = \{x \mid Ax = b, x \ge 0\}$ . x is extreme point iff there exists  $B \subseteq \{1, ..., n\}$  with |B| = m and

- A<sub>B</sub> is non-singular
- $\bullet \ x_B = A_B^{-1}b \ge 0$
- $x_N = 0$

where  $N = \{1, \ldots, n\} \setminus B$ .

**Proof** Take  $B = \{j \mid x_j > 0\}$  and augment with linearly independent columns until |B| = m; always possible since rank(A) = m.



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#### Proof

Take  $B = \{j \mid x_j > 0\}$  and augment with linearly independent columns until |B| = m; always possible since rank(A) = m.



 $x \in \mathbb{R}^n$  is called basic solution (Basislösung) if Ax = b and  $\operatorname{rank}(A_J) = |J|$  where  $J = \{j \mid x_j \neq 0\}$ ;

x is a basic **feasible** solution (gültige Basislösung) if in addition  $x \ge 0$ .

A basis (Basis) is an index set  $B \subseteq \{1, ..., n\}$  with  $rank(A_B) = m$ and |B| = m.



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A BFS fulfills the m equality constraints.

In addition, at least n - m of the  $x_i$ 's are zero. The corresponding non-negativity constraint is fulfilled with equality.

#### Fact:

In a BFS at least n constraints are fulfilled with equality.

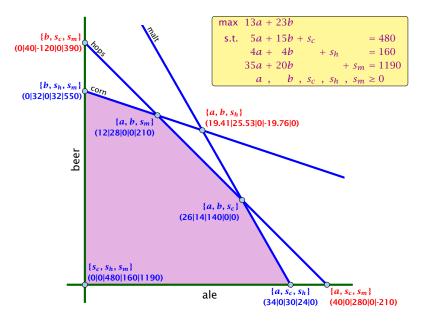


#### **Definition 25**

For a general LP (max{ $c^T x | Ax \le b$ }) with n variables a point x is a basic feasible solution if x is feasible and there exist n (linearly independent) constraints that are tight.



# **Algebraic View**



#### **Fundamental Questions**

#### Linear Programming Problem (LP)

Let  $A \in \mathbb{Q}^{m \times n}$ ,  $b \in \mathbb{Q}^m$ ,  $c \in \mathbb{Q}^n$ ,  $\alpha \in \mathbb{Q}$ . Does there exist  $x \in \mathbb{Q}^n$  s.t. Ax = b,  $x \ge 0$ ,  $c^T x \ge \alpha$ ?

**Questions**:

- Is LP in NP? yes!
- ▶ Is LP in co-NP?
- Is LP in P?

**Proof**:

Given a basis B we can compute the associated basis solution by calculating A<sub>B</sub><sup>-1</sup>b in polynomial time; then we can also compute the profit.



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#### Observation

We can compute an optimal solution to a linear program in time  $\mathcal{O}\left(\binom{n}{m} \cdot \operatorname{poly}(n,m)\right)$ .

- there are only  $\binom{n}{m}$  different bases.
- compute the profit of each of them and take the maximum

What happens if LP is unbounded?



# Enumerating all basic feasible solutions (BFS), in order to find the optimum is slow.

**Simplex Algorithm** [George Dantzig 1947] Move from BFS to adjacent BFS, without decreasing objective function.

Two BFSs are called adjacent if the bases just differ in one variable.



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 $\begin{array}{ll} \max \ 13a + 23b \\ \text{s.t.} \ 5a + 15b + s_c &= 480 \\ 4a + 4b &+ s_h &= 160 \\ 35a + 20b &+ s_m = 1190 \\ a , b , s_c , s_h , s_m \ge 0 \end{array}$ 





#### **4 Simplex Algorithm**

 $\begin{array}{ll} \max \ 13a + 23b \\ \text{s.t.} \ 5a + 15b + s_c &= 480 \\ 4a + 4b &+ s_h &= 160 \\ 35a + 20b &+ s_m = 1190 \\ a , b , s_c , s_h , s_m \ge 0 \end{array}$ 

max Z	ba
$13a + 23b \qquad -Z = 0$	A
$5a + 15b + s_c = 480$	
$4a + 4b + s_h = 160$	S <sub>C</sub>
$35a + 20b + s_m = 1190$	) $\begin{vmatrix} S_h \\ S_m \end{vmatrix}$
$a$ , $b$ , $s_c$ , $s_h$ , $s_m \ge 0$	

$$basis = \{s_c, s_h, s_m\} A = B = 0 Z = 0 s_c = 480 s_h = 160 s_m = 1190$$



**4 Simplex Algorithm** 

max Z	
13a + 23b -	Z = 0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
a, b, s <sub>c</sub> , s <sub>h</sub> , s <sub>m</sub>	≥ 0

basis = 
$$\{s_c, s_h, s_m\}$$
  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

- choose variable to bring into the basis
- chosen variable should have positive coefficient in objective function
- apply devices test to find out by how much the variable can be increased
- pivot on row found by min-ratio test
- the existing basis variable in this row leaves the basis

max Z	
13a + 23b	-Z=0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
a, b, s <sub>c</sub> , s <sub>h</sub> , s <sub>m</sub>	≥ 0

basis = 
$$\{s_c, s_h, s_m\}$$
  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

#### choose variable to bring into the basis

- chosen variable should have positive coefficient in objective function
- apply min-ratio test to find out by how much the variable can be increased
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- the existing basis variable in this row leaves the basis

max Z	
13a + 23b	-Z=0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
a, b, s <sub>c</sub> , s <sub>h</sub> , s <sub>m</sub>	≥ 0

basis = 
$$\{s_c, s_h, s_m\}$$
  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

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13a + 23b	-Z=0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
a, b,s <sub>c</sub> ,s <sub>h</sub> ,s <sub>m</sub>	≥ 0

**basis** = {
$$s_c$$
,  $s_h$ ,  $s_m$ }  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

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$35a + 20b + s_m$	= 1190
a, b, s <sub>c</sub> , s <sub>h</sub> , s <sub>m</sub>	≥ 0

**basis** = {
$$s_c$$
,  $s_h$ ,  $s_m$ }  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

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$5a + 15b + s_c$	= 480
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a, b, s <sub>c</sub> , s <sub>h</sub> , s <sub>m</sub>	≥ 0

basis = 
$$\{s_c, s_h, s_m\}$$
  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

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- chosen variable should have positive coefficient in objective function
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max Z	
13a + 23b	-Z=0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
$a, b, s_c, s_h, s_m$	≥ 0

$basis = \{s_c, s_h, s_m\}$
a = b = 0
Z = 0
$s_c = 480$
$s_h = 160$
$s_m = 1190$

max Z	<b>basis</b> = { $s_c$ , $s_h$ , $s_m$ }
13a + 23b $- Z = 0$	a = b = 0
$5a + 15b + s_c = 480$	Z = 0
$4a + 4b + s_h = 160$	$s_c = 480$ $s_h = 160$
$35a + 20b + s_m = 1190$	$s_h = 100$ $s_m = 1190$
$a$ , $b$ , $s_c$ , $s_h$ , $s_m \ge 0$	

• Choose variable with coefficient > 0 as entering variable.

max Z		<b>basis</b> = { $s_c, s_h, s_n$
13a + 23 <b>b</b>	-Z = 0	a = b = 0
$5a + 15b + s_c$	= 480	Z = 0
$4a + 4b + s_h$	= 160	$s_c = 480$
35 <i>a</i> + 20 <b>b</b> + s	$S_m = 1190$	$s_h = 160$ $s_m = 1190$
$a, b, s_c, s_h, s_c$	$S_m \geq 0$	

n

- Choose variable with coefficient > 0 as entering variable.
- If we keep a = 0 and increase b from 0 to θ > 0 s.t. all constraints (Ax = b, x ≥ 0) are still fulfilled the objective value Z will strictly increase.

max Z		<b>basis</b> = { $s_c$ , $s_h$ , $s_h$
13a + 23 <b>b</b>	-Z = 0	a = b = 0
$5a + 15b + s_c$	= 480	Z = 0
$4a + 4b + s_h$	= 160	$s_c = 480$
$35a + 20b + s_{r}$	n = 1190	$s_h = 160$ $s_m = 1190$
$a, b, s_c, s_h, s_r$	$n \ge 0$	

m

- Choose variable with coefficient > 0 as entering variable.
- If we keep a = 0 and increase b from 0 to θ > 0 s.t. all constraints (Ax = b, x ≥ 0) are still fulfilled the objective value Z will strictly increase.
- For maintaining Ax = b we need e.g. to set  $s_c = 480 15\theta$ .

max Z		<b>basis</b> = { $s_c$ , $s_h$ ,
13a + 23b	-Z = 0	a = b = 0
$5a + 15b + s_c$	= 480	Z = 0
$4a + 4b + s_h$	= 160	$s_c = 480$
35a + 20b + 35a	$s_m = 1190$	$s_h = 160$ $s_m = 1190$
$a, b, s_c, s_h, s_c$	$s_m \geq 0$	

Sm

- Choose variable with coefficient > 0 as entering variable.
- If we keep a = 0 and increase b from 0 to θ > 0 s.t. all constraints (Ax = b, x ≥ 0) are still fulfilled the objective value Z will strictly increase.
- For maintaining Ax = b we need e.g. to set  $s_c = 480 15\theta$ .
- Choosing \(\theta\) = min{480/15, 160/4, 1190/20}\) ensures that in the new solution one current basic variable becomes 0, and no variable goes negative.

max Z		<b>basis</b> = { $s_c$ , $s_h$ , $s_h$
13a + 23b	-Z = 0	a = b = 0
5a + 15 <b>b</b> + sc	= 480	Z = 0
$4a + 4b + s_h$	= 160	$s_c = 480$
$35a + 20b + s_m$	n = 1190	$s_h = 160$ $s_m = 1190$
$a, b, s_c, s_h, s_n$	$\iota \geq 0$	

- Choose variable with coefficient > 0 as entering variable.
- If we keep a = 0 and increase b from 0 to θ > 0 s.t. all constraints (Ax = b, x ≥ 0) are still fulfilled the objective value Z will strictly increase.
- For maintaining Ax = b we need e.g. to set  $s_c = 480 15\theta$ .
- Choosing \(\theta\) = min{480/15, 160/4, 1190/20}\) ensures that in the new solution one current basic variable becomes 0, and no variable goes negative.
- The basic variable in the row that gives min{480/15,160/4,1190/20} becomes the leaving variable.

max Z	
13a + 23b	-Z = 0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
35a + 20b -	$+ s_m = 1190$
a, b, s <sub>c</sub> , s <sub>h</sub>	, $s_m \geq 0$

$$basis = \{s_c, s_h, s_m\} a = b = 0 Z = 0 s_c = 480 s_h = 160 s_m = 1190$$

max Z	
13 <i>a</i> + 23 <b>b</b>	-Z = 0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
$a, b, s_c, s_h, s_m$	≥ 0

$$basis = \{s_c, s_h, s_m\} a = b = 0 Z = 0 s_c = 480 s_h = 160 s_m = 1190$$

Substitute  $b = \frac{1}{15}(480 - 5a - s_c)$ .

max Z	
13 <i>a</i> + 23 <b>b</b>	-Z=0
$5a + 15b + s_c$	= 480
$4a + 4b + s_h$	= 160
$35a + 20b + s_m$	= 1190
$a, b, s_c, s_h, s_m$	≥ 0

basis = 
$$\{s_c, s_h, s_m\}$$
  
 $a = b = 0$   
 $Z = 0$   
 $s_c = 480$   
 $s_h = 160$   
 $s_m = 1190$ 

Substitute  $b = \frac{1}{15}(480 - 5a - s_c)$ .

basis = 
$$\{b, s_h, s_m\}$$
  
 $a = s_c = 0$   
 $Z = 736$   
 $b = 32$   
 $s_h = 32$   
 $s_m = 550$ 

max Z			
$\frac{16}{3}a - \frac{16}{3}a$	$\frac{23}{15}s_c$ -	-Z = -736	
$\frac{1}{3}a + b + b$	$\frac{1}{15}S_c$	= 32	
$\frac{8}{3}a - \frac{1}{3}a$	$\frac{4}{15}s_c + s_h$	= 32	
$\frac{85}{3}a$ –	$\frac{4}{3}s_c + s_m$	= 550	l
a , b ,	$S_c$ , $S_h$ , $S_m$	≥ 0	

<b>basis</b> = $\{b, s_h, s_m\}$
$a = s_c = 0$
Z = 736
<i>b</i> = 32
$s_h = 32$
$s_m = 550$

ſ	max Z	
		<b>basis</b> = { $b, s_h, s_m$ }
	$\frac{16}{3}a - \frac{23}{15}s_c - Z = -736$	$a = s_c = 0$
	$\frac{1}{3}a + b + \frac{1}{15}s_c = 32$	Z = 736
	$\frac{8}{3}a - \frac{4}{15}s_c + s_h = 32$	b = 32
	$\frac{85}{3}a - \frac{4}{3}s_c + s_m = 550$	$s_h = 32$ $s_m = 550$
	$a$ , $b$ , $s_c$ , $s_h$ , $s_m \ge 0$	

max Z			
$\frac{16}{3}a$	$-\frac{23}{15}s_{c}$	-Z = -736	$basis = \{b, s_h, s_m\}$
5	15		$a = s_c = 0$
$\frac{1}{3}a +$	$b + \frac{1}{15}s_c$	= 32	Z = 736
8	$-\frac{4}{15}s_{c}+s_{h}$	= 32	b = 32
0	10	52	$s_h = 32$
$\frac{85}{3}a$	$-\frac{4}{3}s_c$ $+s_r$	n = 550	$s_m = 550$
a		m > 0	
<b>"</b> ,	$b$ , $s_c$ , $s_h$ , $s_r$	$n \ge 0$	

Computing  $min{3 \cdot 32, 3 \cdot 32/8, 3 \cdot 550/85}$  means pivot on line 2.

max Z	
$\frac{16}{3}a - \frac{23}{15}s_c - Z = -736$	$basis = \{b, s_h, s_m\}$
$\frac{3}{3}u - \frac{15}{15}s_c - 2 - 750$	$a = s_c = 0$
$\frac{1}{3}a + b + \frac{1}{15}s_c = 32$	Z = 736
$\frac{8}{3}a - \frac{4}{15}s_c + s_h = 32$	b = 32
5 15	$s_h = 32$
$\frac{85}{3}a - \frac{4}{3}s_c + s_m = 550$	$s_m = 550$
$a, b, s_c, s_h, s_m > 0$	
$a, b, s_c, s_h, s_m \geq 0$	J

Computing min{3 · 32, 3·32/8, 3·550/85} means pivot on line 2. Substitute  $a = \frac{3}{8}(32 + \frac{4}{15}s_c - s_h)$ .

max Z	
$\frac{16}{3}a - \frac{23}{15}s_c - Z = -736$	basis = $\{b, s_h, s_n\}$
5 15	$a = s_c = 0$
$\frac{1}{3}a + b + \frac{1}{15}s_c = 32$	Z = 736
5 15	b = 32
$\frac{8}{3}a - \frac{4}{15}s_c + s_h = 32$	$s_h = 32$
$\frac{85}{3}a - \frac{4}{3}s_c + s_m = 550$	
$_{3}u_{35c} + s_{m} = 550$	$s_m = 550$
$a$ , $b$ , $s_c$ , $s_h$ , $s_m \ge 0$	

Computing  $\min\{3 \cdot 32, \frac{3 \cdot 32}{8}, \frac{3 \cdot 550}{85}\}$  means pivot on line 2. Substitute  $a = \frac{3}{8}(32 + \frac{4}{15}s_c - s_h)$ .

max Z  $- s_c - 2s_h - Z = -800$  $b + \frac{1}{10}s_c - \frac{1}{8}s_h = 28$  Z = 800= 12 b = 28 $a - \frac{1}{10}s_c + \frac{3}{8}s_h$  $\frac{3}{2}s_c - \frac{85}{8}s_h + s_m = 210$  $a, b, s_c, s_h, s_m$  $\geq 0$ 

basis =  $\{a, b, s_m\}$  $s_{c} = s_{h} = 0$ a = 12 $s_m = 210$ 

# Pivoting stops when all coefficients in the objective function are non-positive.

- any feasible solution satisfies all equations in the tableaux
  - in particular:  $Z = 300 s_1 2s_0$ ,  $s \ge 0$ ,  $s_1 \ge 0$
  - hence optimum solution value is at most 800
  - the current solution has value 8000



Pivoting stops when all coefficients in the objective function are non-positive.

Solution is optimal:

any feasible solution satisfies all equations in the tableaux in particular: <a href="https://www.solution.com">https://www.solution.com</a> hence optimum solution value is at most <a href="https://www.solution.com">https://www.solution.com</a> the current solution has value <a href="https://www.solution.com">https://www.solution.com</a>



Pivoting stops when all coefficients in the objective function are non-positive.

- any feasible solution satisfies all equations in the tableaux
- in particular:  $Z = 800 s_c 2s_h, s_c \ge 0, s_h \ge 0$
- hence optimum solution value is at most 800
- the current solution has value 800



Pivoting stops when all coefficients in the objective function are non-positive.

- any feasible solution satisfies all equations in the tableaux
- in particular:  $Z = 800 s_c 2s_h$ ,  $s_c \ge 0$ ,  $s_h \ge 0$
- hence optimum solution value is at most 800
- the current solution has value 800



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- any feasible solution satisfies all equations in the tableaux
- in particular:  $Z = 800 s_c 2s_h$ ,  $s_c \ge 0$ ,  $s_h \ge 0$
- hence optimum solution value is at most 800
- the current solution has value 800



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- any feasible solution satisfies all equations in the tableaux
- in particular:  $Z = 800 s_c 2s_h$ ,  $s_c \ge 0$ ,  $s_h \ge 0$
- hence optimum solution value is at most 800
- the current solution has value 800



#### **Matrix View**

Let our linear program be

$$\begin{array}{rclcrcrc} c_B^T x_B &+& c_N^T x_N &=& Z\\ A_B x_B &+& A_N x_N &=& b\\ x_B &, & x_N &\geq & 0 \end{array}$$

The simplex tableaux for basis *B* is

$$\begin{array}{rcl} (c_{N}^{T}-c_{B}^{T}A_{B}^{-1}A_{N})x_{N} &=& Z-c_{B}^{T}A_{B}^{-1}b\\ Ix_{B} &+& A_{B}^{-1}A_{N}x_{N} &=& A_{B}^{-1}b\\ x_{B} &,& x_{N} &\geq& 0 \end{array}$$

The BFS is given by  $x_N = 0$ ,  $x_B = A_B^{-1}b$ .

If  $(c_N^T - c_B^T A_B^{-1} A_N) \le 0$  we know that we have an optimum solution.



#### 4 Simplex Algorithm

#### **Matrix View**

Let our linear program be

$$c_B^T x_B + c_N^T x_N = Z$$
  

$$A_B x_B + A_N x_N = b$$
  

$$x_B , x_N \ge 0$$

The simplex tableaux for basis *B* is

$$(c_N^T - c_B^T A_B^{-1} A_N) x_N = Z - c_B^T A_B^{-1} b$$
  

$$Ix_B + A_B^{-1} A_N x_N = A_B^{-1} b$$
  

$$x_B , x_N \ge 0$$

The BFS is given by  $x_N = 0, x_B = A_B^{-1}b$ .

If  $(c_N^T - c_B^T A_B^{-1} A_N) \le 0$  we know that we have an optimum solution.



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$$c_B^T x_B + c_N^T x_N = Z$$
  

$$A_B x_B + A_N x_N = b$$
  

$$x_B , x_N \ge 0$$

The simplex tableaux for basis *B* is

$$(c_{N}^{T} - c_{B}^{T}A_{B}^{-1}A_{N})x_{N} = Z - c_{B}^{T}A_{B}^{-1}b$$
  

$$Ix_{B} + A_{B}^{-1}A_{N}x_{N} = A_{B}^{-1}b$$
  

$$x_{B} , \qquad x_{N} \ge 0$$

The BFS is given by  $x_N = 0, x_B = A_B^{-1}b$ .

If  $(c_N^T - c_B^T A_B^{-1} A_N) \le 0$  we know that we have an optimum solution.



#### **Matrix View**

Let our linear program be

$$c_B^T x_B + c_N^T x_N = Z$$
  

$$A_B x_B + A_N x_N = b$$
  

$$x_B , x_N \ge 0$$

The simplex tableaux for basis *B* is

$$(c_{N}^{T} - c_{B}^{T}A_{B}^{-1}A_{N})x_{N} = Z - c_{B}^{T}A_{B}^{-1}b$$
  

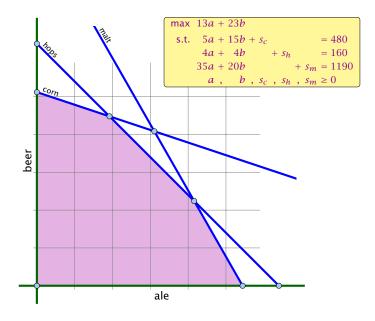
$$Ix_{B} + A_{B}^{-1}A_{N}x_{N} = A_{B}^{-1}b$$
  

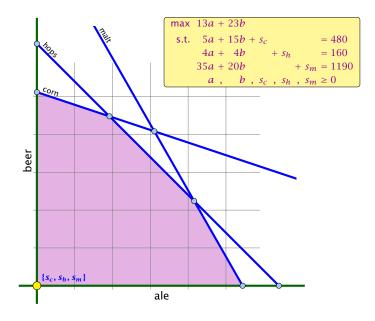
$$x_{B} , \qquad x_{N} \ge 0$$

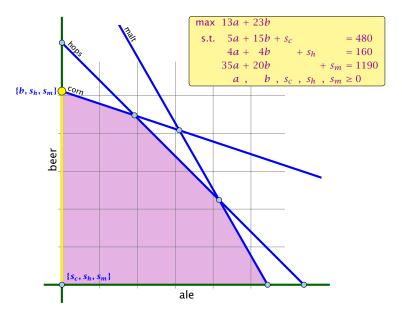
The BFS is given by  $x_N = 0, x_B = A_B^{-1}b$ .

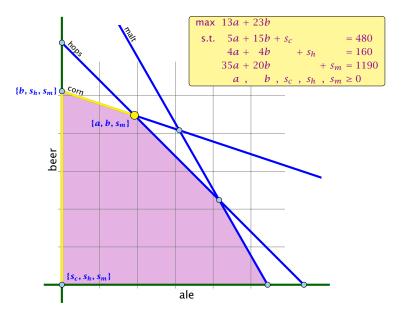
If  $(c_N^T - c_B^T A_B^{-1} A_N) \le 0$  we know that we have an optimum solution.

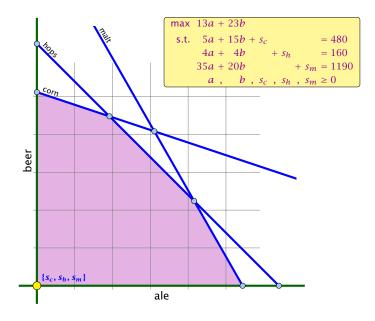
EADS II Harald Räcke **4 Simplex Algorithm** 

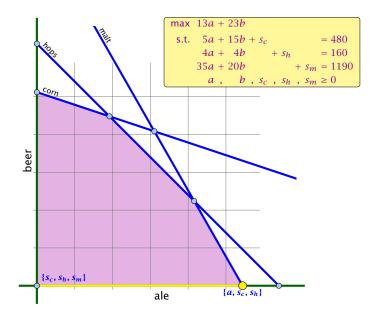


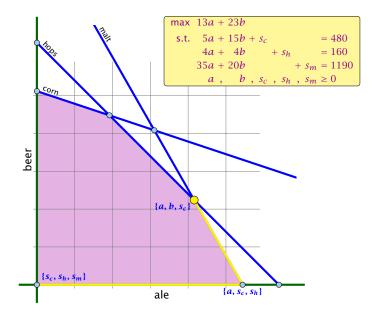


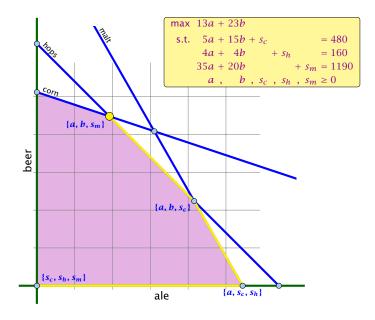












• Given basis *B* with BFS  $x^*$ .

- Choose index  $j \notin B$  in order to increase  $x_j^*$  from 0 to  $\theta > 0$ . Other numbers to variables should store at the statement of the store of the
- Go from  $x^*$  to  $x^* + \theta \cdot d$ .

- $d_{1} = 0$  (normalization)
- $A(x^* \mapsto \partial d) = b$  must hold. Hence Ad = 0.
- Altogether: And preday and a graduate dy which gives



- Given basis *B* with BFS  $x^*$ .
- Choose index  $j \notin B$  in order to increase  $x_i^*$  from 0 to  $\theta > 0$ .
  - Other non-basis variables should stay at 0.
  - Basis variables change to maintain feasibility.
- Go from  $x^*$  to  $x^* + \theta \cdot d$ .

- $d_{1} = 0$  (normalization)
- $dy=0,\, d\in B,\, d=j$
- $A(x^* \rightarrow \partial u) = b$  must hold. Hence Au = 0.
- Altogether: Altogether



- Given basis *B* with BFS  $x^*$ .
- Choose index  $j \notin B$  in order to increase  $x_i^*$  from 0 to  $\theta > 0$ .
  - Other non-basis variables should stay at 0.
  - Basis variables change to maintain feasibility.
- Go from  $x^*$  to  $x^* + \theta \cdot d$ .

Requirements for *d*:

 $d_{1} = 0$  (normalization)

- $(z, b, q) = 0, \beta \in \mathbb{R}, \beta = q h = q$
- $A(x^* \rightarrow \partial u) = b$  must hold. Hence Au = 0.
- Altogether: Altogether



- Given basis *B* with BFS  $x^*$ .
- Choose index  $j \notin B$  in order to increase  $x_i^*$  from 0 to  $\theta > 0$ .
  - Other non-basis variables should stay at 0.
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#### Definition 26 (*j*-th basis direction)

Let *B* be a basis, and let  $j \notin B$ . The vector *d* with  $d_j = 1$  and  $d_{\ell} = 0, \ell \notin B, \ell \neq j$  and  $d_B = -A_B^{-1}A_{*j}$  is called the *j*-th basis direction for *B*.

Going from  $x^*$  to  $x^* + \theta \cdot d$  the objective function changes by

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#### **Definition 27 (Reduced Cost)**

For a basis B the value

$$\tilde{c}_j = c_j - c_B^T A_B^{-1} A_{*j}$$

is called the reduced cost for variable  $x_j$ .

Note that this is defined for every j. If  $j \in B$  then the above term is 0.



Let our linear program be

$$c_B^T x_B + c_N^T x_N = Z$$
  

$$A_B x_B + A_N x_N = b$$
  

$$x_B , x_N \ge 0$$

The simplex tableaux for basis *B* is

$$\begin{array}{rcl} (c_{N}^{T}-c_{B}^{T}A_{B}^{-1}A_{N})x_{N} &=& Z-c_{B}^{T}A_{B}^{-1}b\\ Ix_{B} &+& A_{B}^{-1}A_{N}x_{N} &=& A_{B}^{-1}b\\ x_{B} &,& x_{N} &\geq& 0 \end{array}$$

The BFS is given by  $x_N = 0$ ,  $x_B = A_B^{-1}b$ .

If  $(c_N^T - c_B^T A_B^{-1} A_N) \le 0$  we know that we have an optimum solution.



#### 4 Simplex Algorithm

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The simplex tableaux for basis B is

$$(c_N^T - c_B^T A_B^{-1} A_N) x_N = Z - c_B^T A_B^{-1} b$$
  

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$$x_B , \qquad x_N \ge 0$$

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EADS II Harald Räcke **4 Simplex Algorithm** 

#### **Questions:**

- What happens if the minimatio test fails to give us a value of by which we can safely increase the entering variable? How do we find the initial basic feasible solution?
- Is there always a basis 8 such that

- Then we can terminate because we know that the solution is optimal.
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For this, one computes  $b_i/A_{ie}$  for all constraints i and calculates the minimum positive value.

What does it mean that the ratio  $b_i/A_{ie}$  (and hence  $A_{ie}$ ) is negative for a constraint?

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The set of inequalities is degenerate (also the basis is degenerate).

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A BFS  $x^*$  is called degenerate if the set  $J = \{j \mid x_j^* > 0\}$  fulfills |J| < m.



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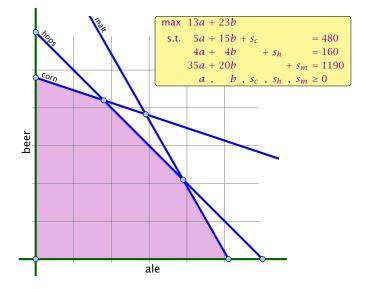
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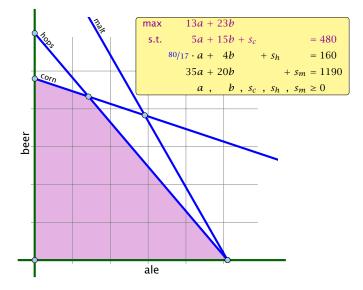
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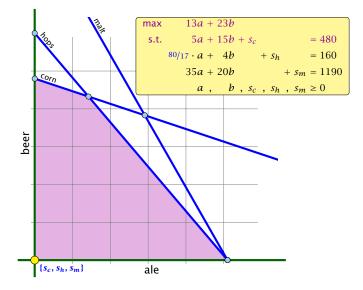
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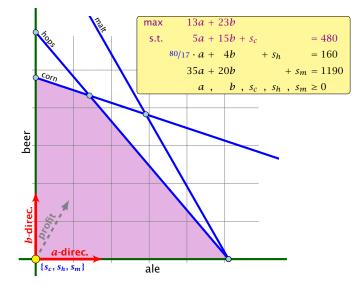


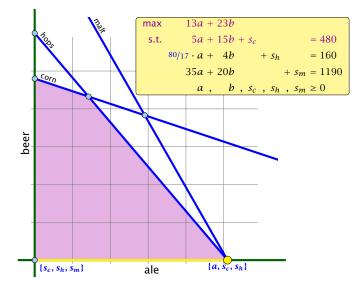
## Non Degenerate Example

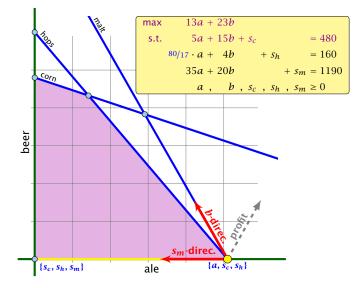


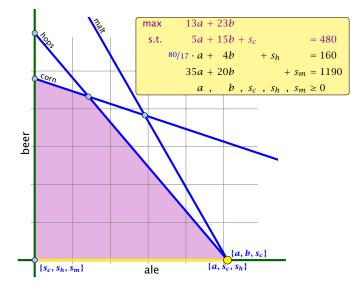


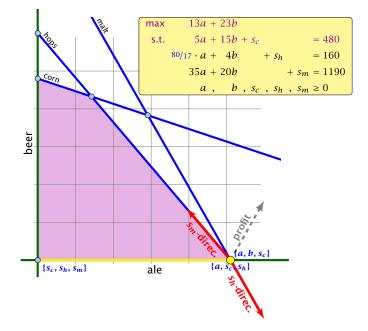


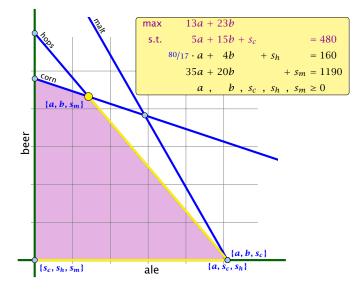


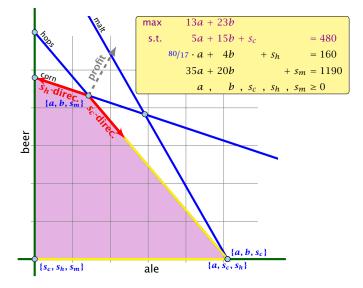












- ► We can choose a column *e* as an entering variable if *c̃<sub>e</sub>* > 0 (*c̃<sub>e</sub>* is reduced cost for *x<sub>e</sub>*).
- The standard choice is the column that maximizes  $\tilde{c}_e$ .
- If  $A_{ie} \leq 0$  for all  $i \in \{1, ..., m\}$  then the maximum is not bounded.
- Otw. choose a leaving variable  $\ell$  such that  $b_{\ell}/A_{\ell e}$  is minimal among all variables *i* with  $A_{ie} > 0$ .
- ► If several variables have minimum b<sub>ℓ</sub>/A<sub>ℓe</sub> you reach a degenerate basis.
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#### What do we have so far?

Suppose we are given an initial feasible solution to an LP. If the LP is non-degenerate then Simplex will terminate.

Note that we either terminate because the min-ratio test fails and we can conclude that the LP is <u>unbounded</u>, or we terminate because the vector of reduced cost is non-positive. In the latter case we have an <u>optimum solution</u>.



•  $Ax \leq b, x \geq 0$ , and  $b \geq 0$ .

- The standard slack from for this problem is  $Ax + Is = b, x \ge 0, s \ge 0$ , where *s* denotes the vector of slack variables.
- Then s = b, x = 0 is a basic feasible solution (how?).
- We directly can start the simplex algorithm.



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# **Optimality**

#### Lemma 29

Let *B* be a basis and  $x^*$  a BFS corresponding to basis *B*.  $\tilde{c} \le 0$  implies that  $x^*$  is an optimum solution to the LP.



#### How do we get an upper bound to a maximization LP?

 $\max 13a + 23b$ s.t.  $5a + 15b \le 480$  $4a + 4b \le 160$  $35a + 20b \le 1190$  $a, b \ge 0$ 

Note that a lower bound is easy to derive. Every choice of  $a, b \ge 0$  gives us a lower bound (e.g. a = 12, b = 28 gives us a lower bound of 800).

If you take a conic combination of the rows (multiply the *i*-th row with  $y_i \ge 0$ ) such that  $\sum_i y_i a_{ij} \ge c_j$  then  $\sum_i y_i b_i$  will be an upper bound.



#### 5.1 Weak Duality

#### How do we get an upper bound to a maximization LP?

Note that a lower bound is easy to derive. Every choice of  $a, b \ge 0$  gives us a lower bound (e.g. a = 12, b = 28 gives us a lower bound of 800).

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#### **Definition 30**

Let  $z = \max\{c^T x \mid Ax \le b, x \ge 0\}$  be a linear program P (called the primal linear program).

The linear program D defined by

$$w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$$

is called the dual problem.



#### **Lemma 31** The dual of the dual problem is the primal problem.

#### Proof:

#### The dual problem is

 $0 = 2 - max[c^{1}x^{1}x^{1}x^{2}x + b]$ 



#### Lemma 31

#### The dual of the dual problem is the primal problem.

#### Proof:

- $w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$
- $w = -\max\{-b^T y \mid -A^T y \leq -c, y \geq 0\}$

#### The dual problem is

- $0 = \alpha (k \alpha ) (k \beta ) (k -$
- 0 = 2 0 + 2 = 0 + 2 + 2 + 2 = 0



#### Lemma 31

The dual of the dual problem is the primal problem.

Proof:

• 
$$w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$$

• 
$$w = -\max\{-b^T \gamma \mid -A^T \gamma \leq -c, \gamma \geq 0\}$$

The dual problem is

 $0 \leq x_1 d + z \leq y_1 d + z \leq y_2 d + z \leq 0$ 

 $0 < \infty, 0 > \infty, 0 > \infty, 0 > \infty$ 



#### Lemma 31

The dual of the dual problem is the primal problem.

**Proof:** 

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$$w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$$

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The dual problem is

- $z = -\min\{-c^T x \mid -Ax \ge -b, x \ge 0\}$
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Let  $z = \max\{c^T x \mid Ax \le b, x \ge 0\}$  and  $w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$  be a primal dual pair.

x is primal feasible iff  $x \in \{x \mid Ax \le b, x \ge 0\}$ 

y is dual feasible, iff  $y \in \{y \mid A^T y \ge c, y \ge 0\}$ .

Theorem 32 (Weak Duality)

Let  $\hat{x}$  be primal feasible and let  $\hat{y}$  be dual feasible. Then

 $c^T \hat{x} \leq z \leq w \leq b^T \hat{y} \; .$ 



Let  $z = \max\{c^T x \mid Ax \le b, x \ge 0\}$  and  $w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$  be a primal dual pair.

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Let  $\hat{x}$  be primal feasible and let  $\hat{y}$  be dual feasible. Then

 $c^T \hat{x} \leq z \leq w \leq b^T \hat{y} \ .$ 



 $A^{T}\hat{y} \ge c \Rightarrow \hat{x}^{T}A^{T}\hat{y} \ge \hat{x}^{T}c \ (\hat{x} \ge 0)$  $A\hat{x} \le b \Rightarrow y^{T}A\hat{x} \le \hat{y}^{T}b \ (\hat{y} \ge 0)$ This gives

Since, there exists primal feasible  $\hat{x}$  with  $c^T \hat{x} = z$ , and dual feasible  $\hat{y}$  with  $b^T y = w$  we get  $z \le w$ .



 $A^T \hat{\gamma} \ge c \Rightarrow \hat{\chi}^T A^T \hat{\gamma} \ge \hat{\chi}^T c \ (\hat{\chi} \ge 0)$ 

This gives

Since, there exists primal feasible  $\hat{x}$  with  $c^T \hat{x} = z$ , and dual feasible  $\hat{y}$  with  $b^T y = w$  we get  $z \le w$ .



 $A^T \hat{y} \ge c \Rightarrow \hat{x}^T A^T \hat{y} \ge \hat{x}^T c \ (\hat{x} \ge 0)$ 

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$$A^{T}\hat{y} \ge c \Rightarrow \hat{x}^{T}A^{T}\hat{y} \ge \hat{x}^{T}c \ (\hat{x} \ge 0)$$
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This gives

#### $c^T \hat{x} \leq \hat{y}^T A \hat{x} \leq b^T \hat{y} \ .$

Since, there exists primal feasible  $\hat{x}$  with  $c^T \hat{x} = z$ , and dual feasible  $\hat{y}$  with  $b^T y = w$  we get  $z \le w$ .



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## 5.2 Simplex and Duality

The following linear programs form a primal dual pair:

$$z = \max\{c^T x \mid Ax = b, x \ge 0\}$$
$$w = \min\{b^T y \mid A^T y \ge c\}$$

This means for computing the dual of a standard form LP, we do not have non-negativity constraints for the dual variables.



#### Primal:

 $\max\{c^T x \mid Ax = b, x \ge 0\}$ 



#### Primal:

$$\max\{c^T x \mid Ax = b, x \ge 0\}$$
$$= \max\{c^T x \mid Ax \le b, -Ax \le -b, x \ge 0\}$$



#### Primal:

$$\max\{c^{T}x \mid Ax = b, x \ge 0\}$$
  
=  $\max\{c^{T}x \mid Ax \le b, -Ax \le -b, x \ge 0\}$   
=  $\max\{c^{T}x \mid \begin{bmatrix} A \\ -A \end{bmatrix} x \le \begin{bmatrix} b \\ -b \end{bmatrix}, x \ge 0\}$ 



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=  $\max\{c^{T}x \mid \begin{bmatrix} A \\ -A \end{bmatrix} x \le \begin{bmatrix} b \\ -b \end{bmatrix}, x \ge 0\}$ 

#### Dual:

$$\min\{[b^T - b^T]y \mid [A^T - A^T]y \ge c, y \ge 0\}$$



#### Primal:

$$\max\{c^{T}x \mid Ax = b, x \ge 0\}$$
  
=  $\max\{c^{T}x \mid Ax \le b, -Ax \le -b, x \ge 0\}$   
=  $\max\{c^{T}x \mid \begin{bmatrix} A \\ -A \end{bmatrix} x \le \begin{bmatrix} b \\ -b \end{bmatrix}, x \ge 0\}$ 

#### Dual:

$$\min\{\begin{bmatrix} b^T & -b^T \end{bmatrix} y \mid \begin{bmatrix} A^T & -A^T \end{bmatrix} y \ge c, y \ge 0\}$$
$$= \min\left\{\begin{bmatrix} b^T & -b^T \end{bmatrix} \cdot \begin{bmatrix} y^+ \\ y^- \end{bmatrix} \mid \begin{bmatrix} A^T & -A^T \end{bmatrix} \cdot \begin{bmatrix} y^+ \\ y^- \end{bmatrix} \ge c, y^- \ge 0, y^+ \ge 0\right\}$$



5.2 Simplex and Duality

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$$\max\{c^{T}x \mid Ax = b, x \ge 0\}$$
  
=  $\max\{c^{T}x \mid Ax \le b, -Ax \le -b, x \ge 0\}$   
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#### Dual:

$$\min\{\begin{bmatrix} b^T & -b^T \end{bmatrix} y \mid \begin{bmatrix} A^T & -A^T \end{bmatrix} y \ge c, y \ge 0\}$$
  
= 
$$\min\left\{\begin{bmatrix} b^T & -b^T \end{bmatrix} \cdot \begin{bmatrix} y^+ \\ y^- \end{bmatrix} \mid \begin{bmatrix} A^T & -A^T \end{bmatrix} \cdot \begin{bmatrix} y^+ \\ y^- \end{bmatrix} \ge c, y^- \ge 0, y^+ \ge 0\right\}$$
  
= 
$$\min\left\{b^T \cdot (y^+ - y^-) \mid A^T \cdot (y^+ - y^-) \ge c, y^- \ge 0, y^+ \ge 0\right\}$$



#### Primal:

$$\max\{c^{T}x \mid Ax = b, x \ge 0\}$$
  
=  $\max\{c^{T}x \mid Ax \le b, -Ax \le -b, x \ge 0\}$   
=  $\max\{c^{T}x \mid \begin{bmatrix} A \\ -A \end{bmatrix} x \le \begin{bmatrix} b \\ -b \end{bmatrix}, x \ge 0\}$ 

#### Dual:

$$\min\{\begin{bmatrix} b^T & -b^T\end{bmatrix} y \mid \begin{bmatrix} A^T & -A^T\end{bmatrix} y \ge c, y \ge 0\}$$
  
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= 
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#### Suppose that we have a basic feasible solution with reduced cost

 $\tilde{c} = c^T - c_B^T A_B^{-1} A \le 0$ 

This is equivalent to  $A^T (A_B^{-1})^T c_B \ge c$ 

 $y^* = (A_B^{-1})^T c_B$  is solution to the dual  $\min\{b^T y | A^T y \ge c\}$ .



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## **Proof of Optimality Criterion for Simplex**

Suppose that we have a basic feasible solution with reduced cost

 $\tilde{c} = c^T - c_B^T A_B^{-1} A \le 0$ 

This is equivalent to  $A^T (A_B^{-1})^T c_B \ge c$ 

 $y^{*} = (A_{B}^{-1})^{T} c_{B} \text{ is solution to the dual } \min\{b^{T} y | A^{T} y \ge c\}.$  $b^{T} y^{*} = (A x^{*})^{T} y^{*} = (A_{B} x^{*}_{B})^{T} y^{*}$  $= (A_{B} x^{*}_{B})^{T} (A^{-1}_{B})^{T} c_{B} = (x^{*}_{B})^{T} A^{T}_{B} (A^{-1}_{B})^{T} c_{B}$  $= c^{T} x^{*}$ 

Hence, the solution is optimal.



## **Proof of Optimality Criterion for Simplex**

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Hence, the solution is optimal.



## 5.3 Strong Duality

 $P = \max\{c^T x \mid Ax \le b, x \ge 0\}$ 

 $n_A$ : number of variables,  $m_A$ : number of constraints

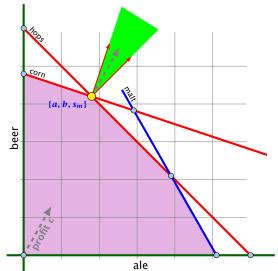
We can put the non-negativity constraints into A (which gives us unrestricted variables):  $\bar{P} = \max\{c^T x \mid \bar{A}x \leq \bar{b}\}$ 

 $n_{ar{A}}=n_A$ ,  $m_{ar{A}}=m_A+n_A$ 

Dual 
$$D = \min\{\bar{b}^T \gamma \mid \bar{A}^T \gamma = c, \gamma \ge 0\}.$$



#### 5.3 Strong Duality



The profit vector c lies in the cone generated by the normals for the hops and the corn constraint (the tight constraints).

### **Strong Duality**

#### Theorem 33 (Strong Duality)

Let P and D be a primal dual pair of linear programs, and let  $z^*$  and  $w^*$  denote the optimal solution to P and D, respectively. Then

 $z^* = w^*$ 



#### Lemma 34 (Weierstrass)

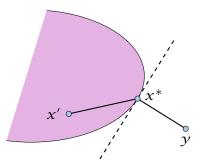
Let X be a compact set and let f(x) be a continuous function on X. Then  $\min\{f(x) : x \in X\}$  exists.

#### (without proof)



#### Lemma 35 (Projection Lemma)

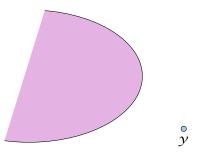
Let  $X \subseteq \mathbb{R}^m$  be a non-empty convex set, and let  $y \notin X$ . Then there exist  $x^* \in X$  with minimum distance from y. Moreover for all  $x \in X$  we have  $(y - x^*)^T (x - x^*) \le 0$ .





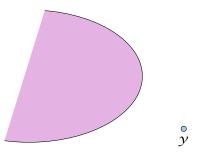
• Define f(x) = ||y - x||.

- We want to apply Weierstrass but *X* may not be bounded.
- $X \neq \emptyset$ . Hence, there exists  $x' \in X$ .
- Define  $X' = \{x \in X \mid ||y x|| \le ||y x'||\}$ . This set is closed and bounded.
- Applying Weierstrass gives the existence.



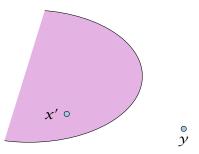


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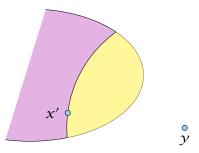


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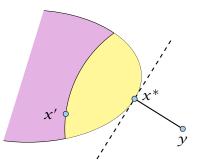


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5.3 Strong Duality



 $x^*$  is minimum. Hence  $\|y - x^*\|^2 \le \|y - x\|^2$  for all  $x \in X$ .



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By convexity:  $x \in X$  then  $x^* + \epsilon(x - x^*) \in X$  for all  $0 \le \epsilon \le 1$ .



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 $\|y-x^*\|^2$ 



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By convexity:  $x \in X$  then  $x^* + \epsilon(x - x^*) \in X$  for all  $0 \le \epsilon \le 1$ .

$$\|y - x^*\|^2 \le \|y - x^* - \epsilon(x - x^*)\|^2$$



 $x^*$  is minimum. Hence  $\|y - x^*\|^2 \le \|y - x\|^2$  for all  $x \in X$ .

By convexity:  $x \in X$  then  $x^* + \epsilon(x - x^*) \in X$  for all  $0 \le \epsilon \le 1$ .

$$\begin{aligned} \|y - x^*\|^2 &\leq \|y - x^* - \epsilon(x - x^*)\|^2 \\ &= \|y - x^*\|^2 + \epsilon^2 \|x - x^*\|^2 - 2\epsilon(y - x^*)^T (x - x^*) \end{aligned}$$



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Hence,  $(y - x^*)^T (x - x^*) \le \frac{1}{2} \epsilon ||x - x^*||^2$ .



5.3 Strong Duality

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Hence,  $(y - x^*)^T (x - x^*) \le \frac{1}{2} \epsilon ||x - x^*||^2$ .

Letting  $\epsilon \to 0$  gives the result.



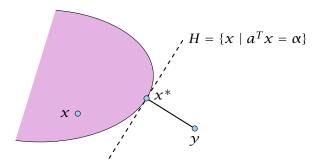
#### **Theorem 36 (Separating Hyperplane)**

Let  $X \subseteq \mathbb{R}^m$  be a non-empty closed convex set, and let  $y \notin X$ . Then there exists a separating hyperplane  $\{x \in \mathbb{R} : a^T x = \alpha\}$ where  $a \in \mathbb{R}^m$ ,  $\alpha \in \mathbb{R}$  that separates y from X.  $(a^T y < \alpha; a^T x \ge \alpha$  for all  $x \in X$ )



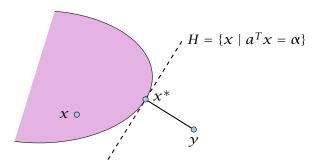
• Let  $x^* \in X$  be closest point to y in X.

- ▶ By previous lemma  $(y x^*)^T (x x^*) \le 0$  for all  $x \in X$ .
- Choose  $a = (x^* y)$  and  $\alpha = a^T x^*$ .
- For  $x \in X$ :  $a^T(x x^*) \ge 0$ , and, hence,  $a^T x \ge \alpha$ .
- Also,  $a^T y = a^T (x^* a) = \alpha ||a||^2 < \alpha$



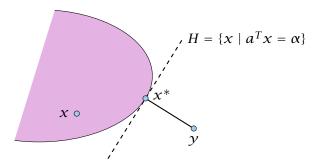


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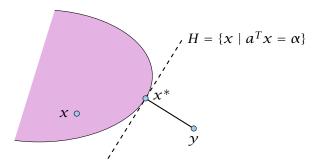
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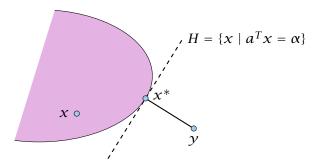
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- Choose  $a = (x^* y)$  and  $\alpha = a^T x^*$ .
- For  $x \in X$ :  $a^T(x x^*) \ge 0$ , and, hence,  $a^T x \ge \alpha$ .

• Also,  $a^T y = a^T (x^* - a) = \alpha - ||a||^2 < \alpha$ 





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- ▶ By previous lemma  $(y x^*)^T (x x^*) \le 0$  for all  $x \in X$ .
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#### Lemma 37 (Farkas Lemma)

Let A be an  $m \times n$  matrix,  $b \in \mathbb{R}^m$ . Then exactly one of the following statements holds.

- **1.**  $\exists x \in \mathbb{R}^n$  with Ax = b,  $x \ge 0$
- **2.**  $\exists y \in \mathbb{R}^m$  with  $A^T y \ge 0$ ,  $b^T y < 0$

Assume  $\hat{x}$  satisfies 1. and  $\hat{y}$  satisfies 2. Then

 $0 > y^T b = y^T A x \ge 0$ 

Hence, at most one of the statements can hold.



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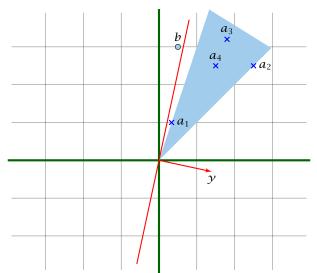
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Hence, at most one of the statements can hold.



#### Farkas Lemma



If b is not in the cone generated by the columns of A, there exists a hyperplane y that separates b from the cone.

Now, assume that 1. does not hold.

Consider  $S = \{Ax : x \ge 0\}$  so that S closed, convex,  $b \notin S$ .

We want to show that there is y with  $A^T y \ge 0$ ,  $b^T y < 0$ .

Let  $\gamma$  be a hyperplane that separates b from S. Hence,  $\gamma^T b < \alpha$ and  $\gamma^T s \ge \alpha$  for all  $s \in S$ .

 $0 \in S \Rightarrow \alpha \le 0 \Rightarrow y^T b < 0$ 

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#### Lemma 38 (Farkas Lemma; different version)

Let A be an  $m \times n$  matrix,  $b \in \mathbb{R}^m$ . Then exactly one of the following statements holds.

- **1.**  $\exists x \in \mathbb{R}^n$  with  $Ax \leq b$ ,  $x \geq 0$
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```
Rewrite the conditions:
```

**1.** 
$$\exists x \in \mathbb{R}^n$$
 with  $\begin{bmatrix} A \ I \end{bmatrix} \cdot \begin{bmatrix} x \\ s \end{bmatrix} = b, x \ge 0, s \ge 0$   
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$$P: z = \max\{c^T x \mid Ax \le b, x \ge 0\}$$

$$D: w = \min\{b^T y \mid A^T y \ge c, y \ge 0\}$$

#### **Theorem 39 (Strong Duality)**

Let P and D be a primal dual pair of linear programs, and let z and w denote the optimal solution to P and D, respectively (i.e., P and D are non-empty). Then

z = w .





 $z \leq w$ : follows from weak duality



- $z \leq w$ : follows from weak duality
- $z \ge w$ :



- $z \leq w$ : follows from weak duality
- $z \ge w$ :
- We show  $z < \alpha$  implies  $w < \alpha$ .



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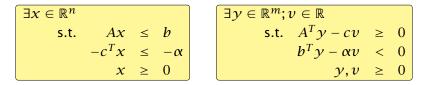
$\exists x \in \mathbb{R}^n$			
s.t.	Ax	$\leq$	b
	$-c^T x$	$\leq$	$-\alpha$
	x	$\geq$	0



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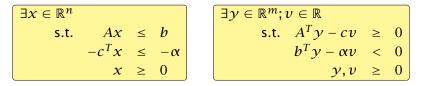




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From the definition of  $\alpha$  we know that the first system is infeasible; hence the second must be feasible.



$$\exists y \in \mathbb{R}^{m}; v \in \mathbb{R}$$
  
s.t.  $A^{T}y - cv \geq 0$   
 $b^{T}y - \alpha v < 0$   
 $y, v \geq 0$ 



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If the solution y, v has v = 0 we have that

$$\exists y \in \mathbb{R}^m \\ \text{s.t.} \quad A^T y \ge 0 \\ b^T y < 0 \\ y \ge 0$$

is feasible.



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is feasible. By Farkas lemma this gives that LP P is infeasible. Contradiction to the assumption of the lemma.



Hence, there exists a solution y, v with v > 0.

We can rescale this solution (scaling both y and v) s.t. v = 1.

Then  $\gamma$  is feasible for the dual but  $b^T \gamma < \alpha$ . This means that  $w < \alpha$ .



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#### Definition 40 (Linear Programming Problem (LP))

Let  $A \in \mathbb{Q}^{m \times n}$ ,  $b \in \mathbb{Q}^m$ ,  $c \in \mathbb{Q}^n$ ,  $\alpha \in \mathbb{Q}$ . Does there exist  $x \in \mathbb{Q}^n$  s.t. Ax = b,  $x \ge 0$ ,  $c^T x \ge \alpha$ ?

#### Questions:

- Is LP in NP?
- Is LP in co-NP? yes!
- Is LP in P?

#### **Proof**:

- Given a primal maximization problem Pland a parameter of Suppose that of some open Plan
- We can prove this by providing an optimal basis for the dual.
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EADS II Harald Räcke

# **Complementary Slackness**

#### Lemma 41

Assume a linear program  $P = \max\{c^T x \mid Ax \le b; x \ge 0\}$  has solution  $x^*$  and its dual  $D = \min\{b^T y \mid A^T y \ge c; y \ge 0\}$  has solution  $y^*$ .

- **1.** If  $x_i^* > 0$  then the *j*-th constraint in *D* is tight.
- **2.** If the *j*-th constraint in *D* is not tight than  $x_i^* = 0$ .
- **3.** If  $y_i^* > 0$  then the *i*-th constraint in *P* is tight.
- **4.** If the *i*-th constraint in *P* is not tight than  $y_i^* = 0$ .



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- **3.** If  $y_i^* > 0$  then the *i*-th constraint in *P* is tight.
- **4.** If the *i*-th constraint in *P* is not tight than  $y_i^* = 0$ .

If we say that a variable  $x_j^*$  ( $y_i^*$ ) has slack if  $x_j^* > 0$  ( $y_i^* > 0$ ), (i.e., the corresponding variable restriction is not tight) and a contraint has slack if it is not tight, then the above says that for a primal-dual solution pair it is not possible that a constraint **and** its corresponding (dual) variable has slack.



## **Proof: Complementary Slackness**

Analogous to the proof of weak duality we obtain

 $c^T x^* \leq y^{*T} A x^* \leq b^T y^*$ 



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From the constraint of the dual it follows that  $y^T A \ge c^T$ . Hence the left hand side is a sum over the product of non-negative numbers. Hence, if e.g.  $(y^T A - c^T)_j > 0$  (the *j*-th constraint in the dual is not tight) then  $x_j = 0$  (2.). The result for (1./3./4.) follows similarly.



Brewer: find mix of ale and beer that maximizes profits

Entrepeneur: buy resources from brewer at minimum cost C, H, M: unit price for corn, hops and malt.

Note that brewer won't sell (at least not all) if e.g. 5C + 4H + 35M < 13 as then brewing ale would be advantageous.

Brewer: find mix of ale and beer that maximizes profits

 $\max 13a + 23b$ s.t.  $5a + 15b \le 480$  $4a + 4b \le 160$  $35a + 20b \le 1190$  $a, b \ge 0$ 

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s.t.	5 <i>C</i>	+	4H	+	35 <i>M</i>	$\geq 13$
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Note that brewer won't sell (at least not all) if e.g. 5C + 4H + 35M < 13 as then brewing ale would be advantageous.

### **Marginal Price:**

- How much money is the brewer willing to pay for additional amount of Corn, Hops, or Malt?
- We are interested in the marginal price, i.e., what happens if we increase the amount of Corn, Hops, and Malt by ε<sub>C</sub>, ε<sub>H</sub>, and ε<sub>M</sub>, respectively.

The profit increases to  $\max\{c^T x \mid Ax \le b + \varepsilon; x \ge 0\}$ . Because of strong duality this is equal to

$$\begin{array}{rcl} \min & (b^T + \epsilon^T) y \\ \text{s.t.} & A^T y &\geq c \\ & y &\geq 0 \end{array}$$



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Therefore we can interpret the dual variables as marginal prices.

Note that with this interpretation, complementary slackness becomes obvious.

- If the brewer has slack of some resource (e.g. corn) then he is not willing to pay anything for it (corresponding dual variable is zero).
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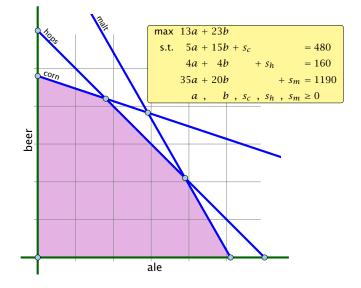
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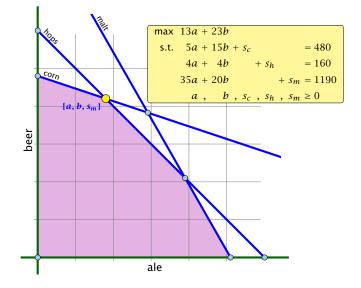
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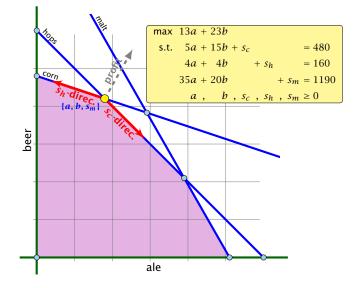
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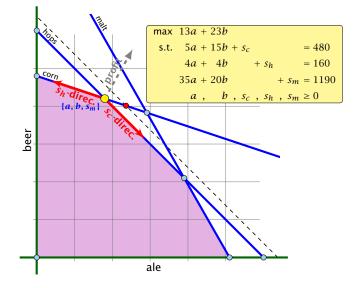
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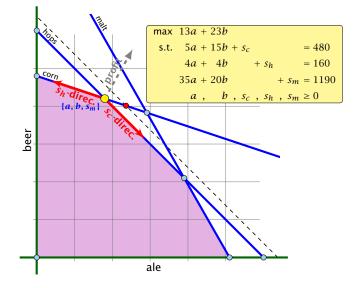


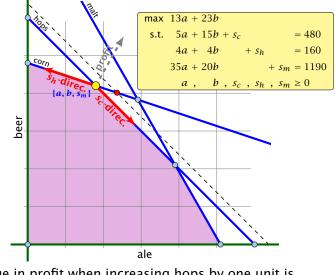




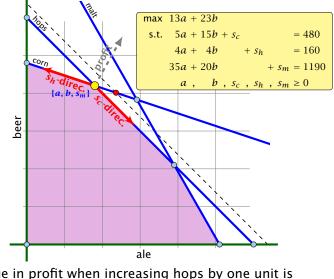








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The change in profit when increasing hops by one unit is  $=\underbrace{c_B^T A_B^{-1}}_{\gamma^*}e_h.$  Of course, the previous argument about the increase in the primal objective only holds for the non-degenerate case.

If the optimum basis is degenerate then increasing the supply of one resource may not allow the objective value to increase.



#### **Definition 42**

An (s, t)-flow in a (complete) directed graph  $G = (V, V \times V, c)$  is a function  $f : V \times V \mapsto \mathbb{R}_0^+$  that satisfies

**1.** For each edge (x, y)

 $0 \leq f_{xy} \leq c_{xy}$  .

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**2.** For each  $v \in V \setminus \{s, t\}$ 

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max		$\sum_{z} f_{sz} - \sum_{z} f_{zs}$			
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s.t.	$f_{xy}(x, y \neq s, t)$ :	$1\ell_{xy}-1p_x+1p_y$	$\geq$	0
	$f_{sy}(y \neq s,t)$ :	$1\ell_{sy}$ $+1p_y$	$\geq$	1
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5.5 Computing Duals



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We can interpret the  $\ell_{xy}$  value as assigning a length to every edge.

The value  $p_x$  for a variable, then can be seen as the distance of x to t (where the distance from s to t is required to be 1 since  $p_s = 1$ ).



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# One can show that there is an optimum LP-solution for the dual problem that gives an integral assignment of variables.

This means  $p_x = 1$  or  $p_x = 0$  for our case. This gives rise to a cut in the graph with vertices having value 1 on one side and the other vertices on the other side. The objective function then evaluates the capacity of this cut.

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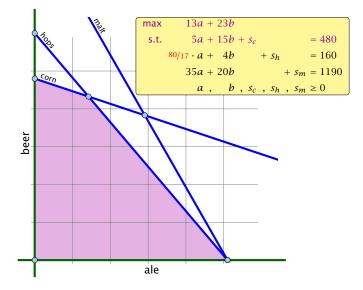
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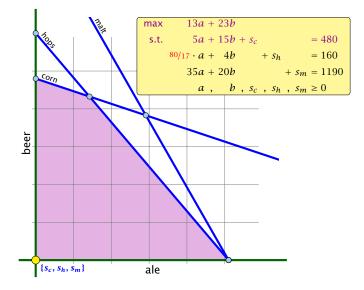


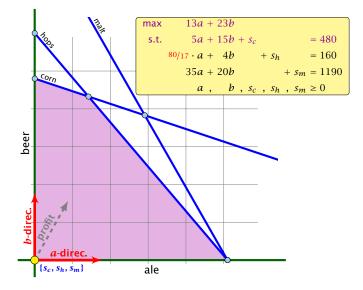
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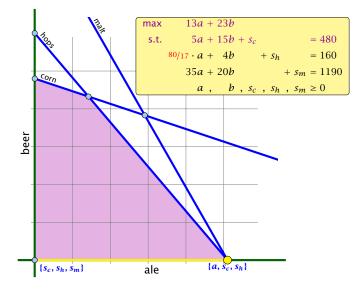
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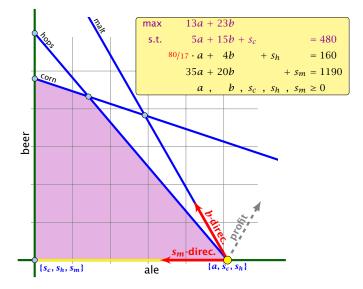


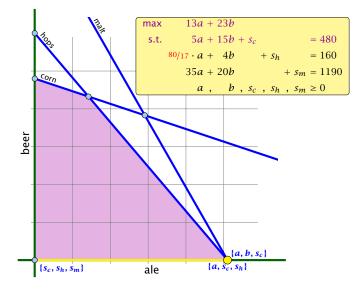


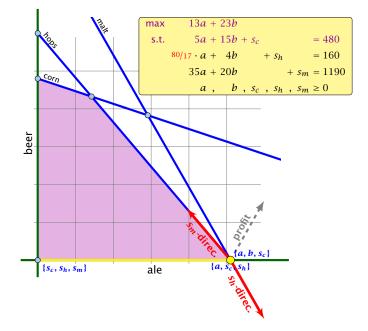




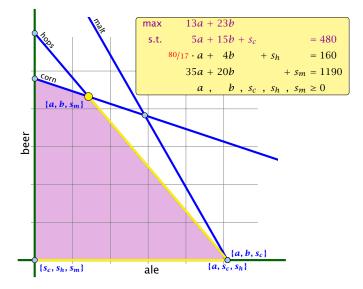




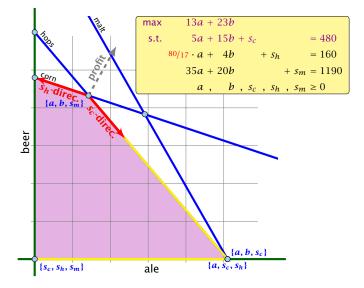




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Idea:

Given feasible LP :=  $\max\{c^T x, Ax = b; x \ge 0\}$ . Change it into LP' :=  $\max\{c^T x, Ax = b', x \ge 0\}$  such that

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If a set 0 of basis variables corresponds to an object of basis variables corresponds to an infeasible basis (i.e. 45,25,20), then 0 corresponds to an infeasible basis in 0.21 (note that columns in 0), are linearly independent).

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- II. If a set *B* of basis variables corresponds to an infeasible basis (i.e.  $A_B^{-1}b \neq 0$ ) then *B* corresponds to an infeasible basis in LP' (note that columns in  $A_B$  are linearly independent).
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#### Perturbation

Let *B* be index set of some basis with basic solution

 $x_{B}^{*} = A_{B}^{-1}b \ge 0, x_{N}^{*} = 0$  (i.e. *B* is feasible)

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 for  $\varepsilon > 0$ .

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The new LP is feasible because the set B of basis variables provides a feasible basis:

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Hence,  $\tilde{B}$  is not feasible.



Let  $\tilde{B}$  be a basis. It has an associated solution

$$x_{\tilde{B}}^* = A_{\tilde{B}}^{-1}b + A_{\tilde{B}}^{-1}A_B\begin{pmatrix}\varepsilon\\\vdots\\\varepsilon^m\end{pmatrix}$$

#### in the perturbed instance.

We can view each component of the vector as a polynom with variable  $\varepsilon$  of degree at most m.

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▶ If it terminates because it finds a variable  $x_j$  with  $\tilde{c}_j > 0$  for which the *j*-th basis direction *d*, fulfills  $d \ge 0$  we know that LP' is unbounded. The basis direction does not depend on *b*. Hence, we also know that LP is unbounded.



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We choose the entering variable arbitrarily as before ( $\tilde{c}_e > 0$ , of course).

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In the following we assume that  $b \ge 0$ . This can be obtained by replacing the initial system  $(A \mid b)$  by  $(A_B^{-1}A \mid A_B^{-1}b)$  where *B* is the index set of a feasible basis (found e.g. by the first phase of the Two-phase algorithm).

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6 Degeneracy Revisited

#### **Matrix View**

Let our linear program be

$$c_B^T x_B + c_N^T x_N = Z$$
  

$$A_B x_B + A_N x_N = b$$
  

$$x_B , x_N \ge 0$$

The simplex tableaux for basis *B* is

$$(c_{N}^{T} - c_{B}^{T}A_{B}^{-1}A_{N})x_{N} = Z - c_{B}^{T}A_{B}^{-1}b$$
  

$$Ix_{B} + A_{B}^{-1}A_{N}x_{N} = A_{B}^{-1}b$$
  

$$x_{B} , \qquad x_{N} \ge 0$$

The BFS is given by  $x_N = 0, x_B = A_B^{-1}b$ .

If  $(c_N^T - c_B^T A_B^{-1} A_N) \le 0$  we know that we have an optimum solution.

EADS II Harald Räcke 6 Degeneracy Revisited

# **Lexicographic Pivoting**

# LP chooses an arbitrary leaving variable that has $\hat{A}_{\ell e} > 0$ and minimizes $\theta_{\ell} = \frac{\hat{h}_{\ell}}{\hat{A}_{\ell e}} = \frac{(A_{\ell}^{-1}b)_{\ell}}{(A_{\ell}^{-1}A_{\ell}c)_{\ell}}$ .

 $\ell$  is the index of a leaving variable within *B*. This means if e.g. *B* = {1,3,7,14} and leaving variable is 3 then  $\ell$  = 2.



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#### **Definition 44**

 $u \leq_{\mathsf{lex}} v$  if and only if the first component in which u and v differ fulfills  $u_i \leq v_i$ .



 $LP^\prime$  chooses an index that minimizes

 $\theta_\ell$ 



LP' chooses an index that minimizes

$$\theta_{\ell} = \frac{\left(A_B^{-1}\left(b + \begin{pmatrix} \varepsilon \\ \vdots \\ \varepsilon^m \end{pmatrix}\right)\right)_{\ell}}{(A_B^{-1}A_{*\ell})_{\ell}}$$



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$$= \frac{\ell \text{-th row of } A_B^{-1}(b \mid I)}{(A_B^{-1}A_{*e})_{\ell}} \begin{pmatrix} 1 \\ \varepsilon \\ \vdots \\ \varepsilon^m \end{pmatrix}$$



This means you can choose the variable/row  $\ell$  for which the vector

 $\frac{\ell\text{-th row of }A_B^{-1}(b \mid I)}{(A_B^{-1}A_{*e})_\ell}$ 

is lexicographically minimal.

Of course only including rows with  $(A_B^{-1}A_{*e})_{\ell} > 0$ .

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7 Klee Minty Cube

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#### Can we obtain a better analysis?



#### Observation

Simplex visits every feasible basis at most once.



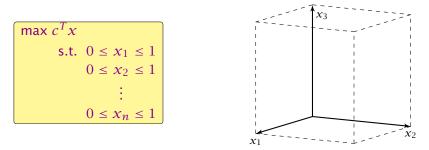
#### Observation

Simplex visits every feasible basis at most once.

However, also the number of feasible bases can be very large.



## Example

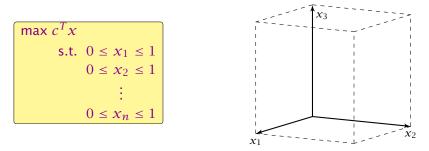


2n constraint on n variables define an n-dimensional hypercube as feasible region.

The feasible region has  $2^n$  vertices.



## Example



However, Simplex may still run quickly as it usually does not visit all feasible bases.

In the following we give an example of a feasible region for which there is a bad Pivoting Rule.

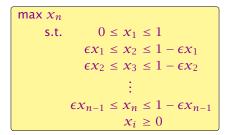


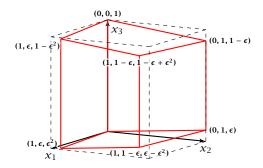
A Pivoting Rule defines how to choose the entering and leaving variable for an iteration of Simplex.

In the non-degenerate case after choosing the entering variable the leaving variable is unique.



#### **Klee Minty Cube**





- ▶ We have 2*n* constraints, and 3*n* variables (after adding slack variables to every constraint).
- Every basis is defined by 2n variables, and n non-basic variables.
- There exist degenerate vertices.
- The degeneracies come from the non-negativity constraints, which are superfluous.
- ▶ In the following all variables *x*<sub>i</sub> stay in the basis at all times.
- Then, we can uniquely specify a basis by choosing for each variable whether it should be equal to its lower bound, or equal to its upper bound (the slack variable corresponding to the non-tight constraint is part of the basis).
- We can also simply identify each basis/vertex with the corresponding hypercube vertex obtained by letting ε → 0.

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- In the following we specify a sequence of bases (identified by the corresponding hypercube node) along which the objective function strictly increases.
- The basis  $(0, \ldots, 0, 1)$  is the unique optimal basis.
- ► Our sequence S<sub>n</sub> starts at (0,...,0) ends with (0,...,0,1) and visits every node of the hypercube.
- An unfortunate Pivoting Rule may choose this sequence, and, hence, require an exponential number of iterations.



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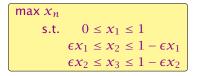
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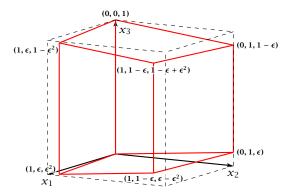


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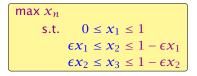


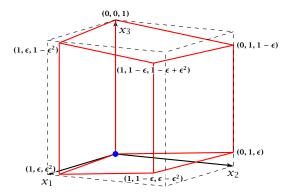
#### **Klee Minty Cube**

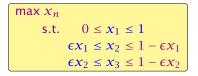


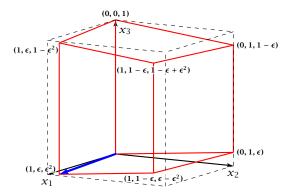


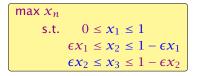
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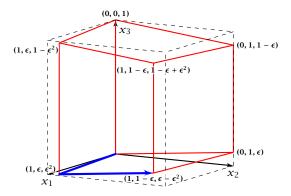


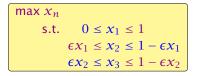


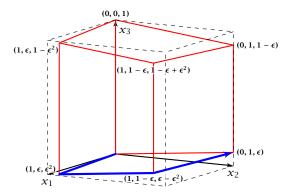


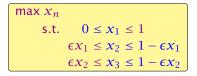


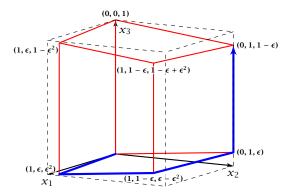


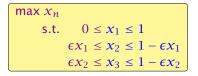


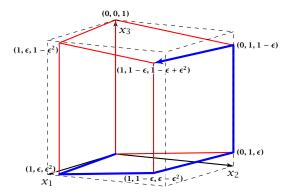


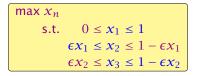


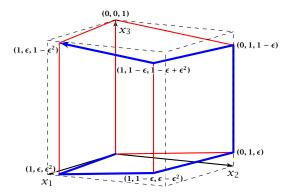


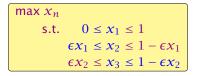


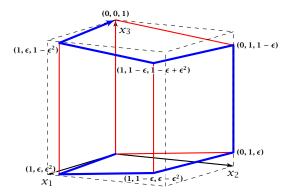












The sequence  $S_n$  that visits every node of the hypercube is defined recursively

$$(0, ..., 0, 0, 0)$$

$$\begin{cases} S_{n-1} \\ (0, ..., 0, 1, 0) \\ 0, ..., 0, 1, 1) \\ \\ \\ S_{n-1}^{\mathsf{rev}} \\ (0, ..., 0, 0, 1) \end{cases}$$

The non-recursive case is  $S_1 = 0 \rightarrow 1$ 



#### Lemma 45

The objective value  $x_n$  is increasing along path  $S_n$ .

#### **Proof by induction:**

n = 1: obvious, since  $S_1 = 0 \rightarrow 1$ , and 1 > 0.

- For the first part the value of  $S_{2} = S_{2} = 0$
- By induction hypothesis ....... is increasing along ....... hence, also .....
- Going from (0), and (100 to (0), and 0.5 (1) increases on for small enough a
- For the remaining path  $\beta_{1}^{\rm res}$  we have  $\beta_{2}$   $\beta$   $\beta$   $\beta$  .
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#### Observation

The simplex algorithm takes at most  $\binom{n}{m}$  iterations. Each iteration can be implemented in time  $\mathcal{O}(mn)$ .

In practise it usually takes a linear number of iterations.



#### Theorem

For almost all known deterministic pivoting rules (rules for choosing entering and leaving variables) there exist lower bounds that require the algorithm to have exponential running time ( $\Omega(2^{\Omega(n)})$ ) (e.g. Klee Minty 1972).



#### Theorem

For some standard randomized pivoting rules there exist subexponential lower bounds ( $\Omega(2^{\Omega(n^{\alpha})})$  for  $\alpha > 0$ ) (Friedmann, Hansen, Zwick 2011).



**Conjecture** (Hirsch 1957)

The edge-vertex graph of an m-facet polytope in d-dimensional Euclidean space has diameter no more than m - d.

The conjecture has been proven wrong in 2010.

But the question whether the diameter is perhaps of the form O(poly(m, d)) is open.



- Suppose we want to solve  $\min\{c^T x \mid Ax \ge b; x \ge 0\}$ , where  $x \in \mathbb{R}^d$  and we have *m* constraints.
- ▶ In the worst-case Simplex runs in time roughly  $O(m(m+d)\binom{m+d}{m}) \approx (m+d)^m$ . (slightly better bounds on the running time exist, but will not be discussed here).
- ▶ If *d* is much smaller than *m* one can do a lot better.
- ► In the following we develop an algorithm with running time O(d! · m), i.e., linear in m.



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- ▶ If *d* is much smaller than *m* one can do a lot better.
- ► In the following we develop an algorithm with running time O(d! · m), i.e., linear in m.



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#### Setting:

We assume an LP of the form

$$\begin{array}{cccc} \min & c^T x \\ \text{s.t.} & Ax &\geq b \\ & x &\geq 0 \end{array}$$

• We assume that the LP is **bounded**.



# **Ensuring Conditions**

Given a standard minimization LP

$$\begin{array}{rcl} \min & c^T x \\ \text{s.t.} & Ax &\geq b \\ & x &\geq 0 \end{array}$$

how can we obtain an LP of the required form?

Compute a lower bound on c<sup>T</sup>x for any basic feasible solution.



# Let s denote the smallest common multiple of all denominators of entries in A, b.

Multiply entries in A, b by s to obtain integral entries. This does not change the feasible region.

Add slack variables to A; denote the resulting matrix with  $ar{A}.$ 



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#### Theorem 46 (Cramers Rule)

Let M be a matrix with  $det(M) \neq 0$ . Then the solution to the system Mx = b is given by

 $x_i = rac{\det(M_j)}{\det(M)}$  ,

where  $M_i$  is the matrix obtained from M by replacing the *i*-th column by the vector b.



Define

Further, we have



Define

$$X_{i} = \begin{pmatrix} | & | & | & | & | \\ e_{1} \cdots e_{i-1} \mathbf{x} e_{i+1} \cdots e_{n} \\ | & | & | & | \end{pmatrix}$$

Note that expanding along the *i*-th column gives that  $det(X_i) = x_i$ .

$$MX_{j} = \begin{pmatrix} | & | & | & | \\ Me_{1} \cdots Me_{i-1} Mx Me_{i+1} \cdots Me_{n} \\ | & | & | \end{pmatrix} = M_{i}$$
  
Hence,  
$$x_{i} = \det(X_{i}) = \frac{\det(M_{i})}{\det(M)}$$



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$$\blacktriangleright \text{ Hence,}$$

$$\det(M_{i})$$

$$x_i = \det(X_i) = \frac{\det(M_i)}{\det(M)}$$



Let Z be the maximum absolute entry occuring in  $\bar{A}$ ,  $\bar{b}$  or c. Let C denote the matrix obtained from  $\bar{A}_B$  by replacing the *j*-th column with vector  $\bar{b}$  (for some *j*).

Observe that

 $|\det(C)|$ 



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$$|\det(C)| = \left| \sum_{\pi \in S_m} \operatorname{sgn}(\pi) \prod_{1 \le i \le m} C_{i\pi(i)} \right|$$



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#### Alternatively, Hadamards inequality gives

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 $|\det(C)| \leq \prod_{i=1}^m \|C_{*i}\|$ 



#### Alternatively, Hadamards inequality gives

$$|\det(C)| \le \prod_{i=1}^{m} ||C_{*i}|| \le \prod_{i=1}^{m} (\sqrt{m}Z)$$

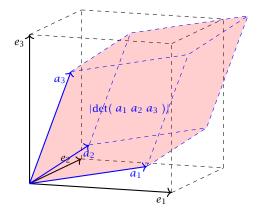


#### Alternatively, Hadamards inequality gives

$$|\det(C)| \le \prod_{i=1}^{m} ||C_{*i}|| \le \prod_{i=1}^{m} (\sqrt{m}Z)$$
$$\le m^{m/2} Z^m .$$



## **Hadamards Inequality**



Hadamards inequality says that the volume of the red parallelepiped (Spat) is smaller than the volume in the black cube (if  $||e_1|| = ||a_1||$ ,  $||e_2|| = ||a_2||$ ,  $||e_3|| = ||a_3||$ ).



## **Ensuring Conditions**

#### Given a standard minimization LP

$$\begin{array}{ccc} \min & c^T x \\ \text{s.t.} & Ax \geq b \\ & x \geq 0 \end{array}$$

how can we obtain an LP of the required form?

Compute a lower bound on c<sup>T</sup>x for any basic feasible solution. Add the constraint c<sup>T</sup>x ≥ -mZ(m! · Z<sup>m</sup>) - 1. Note that this constraint is superfluous unless the LP is unbounded.

## **Ensuring Conditions**

Compute an optimum basis for the new LP.

- ► If the cost is  $c^T x = -(mZ)(m! \cdot Z^m) 1$  we know that the original LP is unbounded.
- Otw. we have an optimum basis.



We give a routine SeidelLP( $\mathcal{H}, d$ ) that is given a set  $\mathcal{H}$  of explicit, non-degenerate constraints over d variables, and minimizes  $c^T x$  over all feasible points.



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- 12: **if**  $\hat{x}^*$  = infeasible **then**
- 13: return infeasible

14: **else** 

15: add the value of  $x_\ell$  to  $\hat{x}^*$  and return the solution

- If d = 1 we can solve the 1-dimensional problem in time  $O(\max\{m, 1\})$ .
- If d > 1 and m = 0 we take time O(d) to return d-dimensional vector x.
- The first recursive call takes time T(m-1, d) for the call plus O(d) for checking whether the solution fulfills h.
- ▶ If we are unlucky and  $\hat{x}^*$  does not fulfill h we need time  $\mathcal{O}(d(m+1)) = \mathcal{O}(dm)$  to eliminate  $x_\ell$ . Then we make a recursive call that takes time T(m-1, d-1).
- The probability of being unlucky is at most d/m as there are at most d constraints whose removal will decrease the objective function



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This gives the recurrence

$$T(m,d) = \begin{cases} \mathcal{O}(\max\{1,m\}) & \text{if } d = 1\\ \mathcal{O}(d) & \text{if } d > 1 \text{ and } m = 0\\ \mathcal{O}(d) + T(m-1,d) + \\ \frac{d}{m}(\mathcal{O}(dm) + T(m-1,d-1)) & \text{otw.} \end{cases}$$

Note that T(m, d) denotes the expected running time.



Let *C* be the largest constant in the  $\mathcal{O}$ -notations.

$$T(m,d) = \begin{cases} C \max\{1,m\} & \text{if } d = 1\\ Cd & \text{if } d > 1 \text{ and } m = 0\\ Cd + T(m-1,d) + & \\ \frac{d}{m}(Cdm + T(m-1,d-1)) & \text{otw.} \end{cases}$$

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d > 1; m = 0:

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$$\leq Cd + Cf(d)(m-1) + Cd^2 + \frac{d}{m}Cf(d-1)(m-1)$$
  

$$\leq 2Cd^2 + Cf(d)(m-1) + dCf(d-1)$$
  

$$\leq Cf(d)m$$



d > 1; m > 1: (by induction hypothesis statm. true for d' < d,  $m' \ge 0$ ; and for d' = d, m' < m)

$$T(m,d) = \mathcal{O}(d) + T(m-1,d) + \frac{d}{m} \Big( \mathcal{O}(dm) + T(m-1,d-1) \Big)$$
  

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if  $f(d) \ge df(d-1) + 2d^2$ .



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since  $\sum_{i\geq 1} \frac{i^2}{i!}$  is a constant.



# Complexity

#### LP Feasibility Problem (LP feasibility)

Given  $A \in \mathbb{Z}^{m \times n}$ ,  $b \in \mathbb{Z}^m$ . Does there exist  $x \in \mathbb{R}$  with Ax = b,  $x \ge 0$ ?



#### Input size

▶ The number of bits to represent a number  $a \in \mathbb{Z}$  is

## $\lceil \log_2(|a|) \rceil + 1$

$$\langle M \rangle := \sum_{i,j} \lceil \log_2(|m_{ij}|) + 1 \rceil$$

- In the following we assume that input matrices are encoded in a standard way, where each number is encoded in binary and then suitable separators are added in order to separate distinct number from each other.
- Then the input length is  $L = \Theta(\langle A \rangle + \langle b \rangle)$ .

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- In the following we sometimes refer to L := ⟨A⟩ + ⟨b⟩ as the input size (even though the real input size is something in Θ(⟨A⟩ + ⟨b⟩)).
- In order to show that LP-decision is in NP we show that if there is a solution x then there exists a small solution for which feasibility can be verified in polynomial time (polynomial in L).



#### Suppose that Ax = b; $x \ge 0$ is feasible.

Then there exists a basic feasible solution. This means a set B of basic variables such that

 $x_B = A_B^{-1}b$ 

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Lemma 47

Let  $M \in \mathbb{Z}^{m \times m}$  be an invertible matrix and let  $b \in \mathbb{Z}^m$ . Further define  $L = \langle M \rangle + \langle b \rangle + n \log_2 n$ . Then a solution to Mx = b has rational components  $x_j$  of the form  $\frac{D_j}{D}$ , where  $|D_j| \le 2^L$  and  $|D| \le 2^L$ .

**Proof:** Cramers rules says that we can compute  $x_j$  as

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Analogously for  $det(M_j)$ .



Given an LP max{ $c^T x | Ax = b; x \ge 0$ } do a binary search for the optimum solution

(Add constraint  $c^T x - \delta = M$ ;  $\delta \ge 0$  or  $(c^T x \ge M)$ . Then checking for feasibility shows whether optimum solution is larger or smaller than M).

If the LP is feasible then the binary search finishes in at most

$$\log_2\left(rac{2n2^{2L'}}{1/2^{L'}}
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Let  $M_{\text{max}} = n2^{2L'}$  be an upper bound on the objective value of a basic feasible solution.

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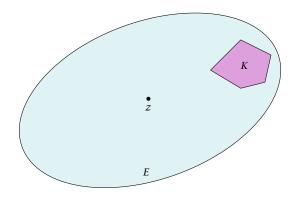


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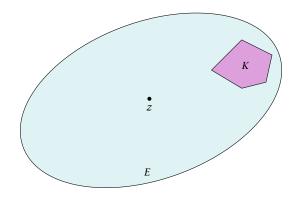
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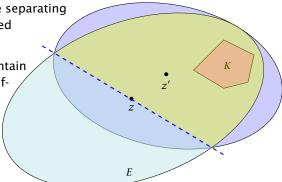
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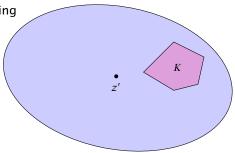
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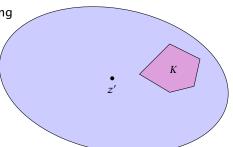


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- REPEAT





#### Issues/Questions:

- How do you choose the first Ellipsoid? What is its volume?
- How do you measure progress? By how much does the volume decrease in each iteration?
- When can you stop? What is the minimum volume of a non-empty polytop?



A mapping  $f : \mathbb{R}^n \to \mathbb{R}^n$  with f(x) = Lx + t, where *L* is an invertible matrix is called an affine transformation.



A ball in  $\mathbb{R}^n$  with center *c* and radius *r* is given by

$$B(c,r) = \{x \mid (x-c)^T (x-c) \le r^2\} \\ = \{x \mid \sum_i (x-c)_i^2 / r^2 \le 1\}$$

B(0,1) is called the unit ball.



An affine transformation of the unit ball is called an ellipsoid.



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From f(x) = Lx + t follows  $x = L^{-1}(f(x) - t)$ .

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#### **Definition 50**

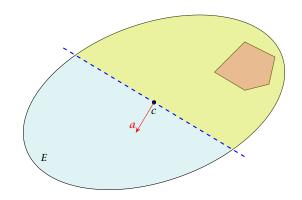
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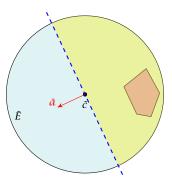
where  $Q = LL^T$  is an invertible matrix.





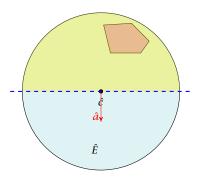


• Use  $f^{-1}$  (recall that f = Lx + t is the affine transformation of the unit ball) to rotate/distort the ellipsoid (back) into the unit ball.



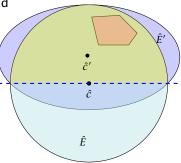


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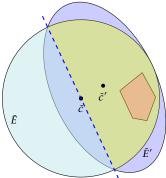


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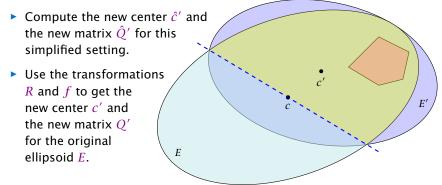


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- Use the transformations *R* and *f* to get the new center *c'* and the new matrix *Q'* for the original ellipsoid *E*.

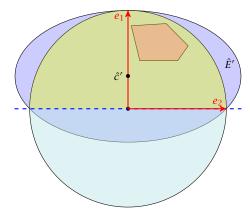




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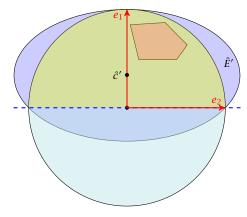




• The new center lies on axis  $x_1$ . Hence,  $\hat{c}' = te_1$  for t > 0.

▶ The vectors  $e_1, e_2, ...$  have to fulfill the ellipsoid constraint with equality. Hence  $(e_i - \hat{c}')^T \hat{Q}'^{-1} (e_i - \hat{c}') = 1$ .





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EADS II Harald Räcke

- To obtain the matrix  $\hat{Q}'^{-1}$  for our ellipsoid  $\hat{E}'$  note that  $\hat{E}'$  is axis-parallel.
- Let *a* denote the radius along the  $x_1$ -axis and let *b* denote the (common) radius for the other axes.
- The matrix

$$\hat{L}' = \begin{pmatrix} a & 0 & \dots & 0 \\ 0 & b & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ 0 & \dots & 0 & b \end{pmatrix}$$

maps the unit ball (via function  $\hat{f}'(x) = \hat{L}'x$ ) to an axis-parallel ellipsoid with radius a in direction  $x_1$  and b in all other directions.



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As  $\hat{Q}' = \hat{L}' \hat{L}'^{t}$  the matrix  $\hat{Q}'^{-1}$  is of the form  $\hat{Q}'^{-1} = \begin{pmatrix} \frac{1}{a^{2}} & 0 & \dots & 0\\ 0 & \frac{1}{b^{2}} & \ddots & \vdots\\ \vdots & \ddots & \ddots & 0\\ 0 & \dots & 0 & \frac{1}{b^{2}} \end{pmatrix}$ 



$$\begin{pmatrix} e_1 - \hat{c}' \end{pmatrix}^T \hat{Q}'^{-1} (e_1 - \hat{c}') = 1 \text{ gives} \\ \begin{pmatrix} 1 - t \\ 0 \\ \vdots \\ 0 \end{pmatrix}^T \cdot \begin{pmatrix} \frac{1}{a^2} & 0 & \dots & 0 \\ 0 & \frac{1}{b^2} & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ 0 & \dots & 0 & \frac{1}{b^2} \end{pmatrix} \cdot \begin{pmatrix} 1 - t \\ 0 \\ \vdots \\ 0 \end{pmatrix} = 1$$

• This gives  $(1 - t)^2 = a^2$ .



For  $i \neq 1$  the equation  $(e_i - \hat{c}')^T \hat{Q}'^{-1} (e_i - \hat{c}') = 1$  looks like (here i = 2)

$$\begin{pmatrix} -t \\ 1 \\ 0 \\ \vdots \\ 0 \end{pmatrix}^{T} \cdot \begin{pmatrix} \frac{1}{a^{2}} & 0 & \dots & 0 \\ 0 & \frac{1}{b^{2}} & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ 0 & \dots & 0 & \frac{1}{b^{2}} \end{pmatrix} \cdot \begin{pmatrix} -t \\ 1 \\ 0 \\ \vdots \\ 0 \end{pmatrix} = 1$$

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$$\frac{t^2}{a^2} + \frac{1}{b^2} = 1$$
, and hence

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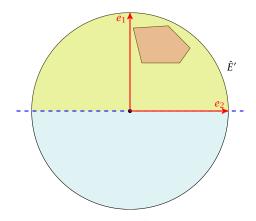
#### **Summary**

So far we have

$$a = 1 - t$$
 and  $b = \frac{1 - t}{\sqrt{1 - 2t}}$ 

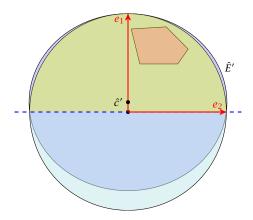


We still have many choices for *t*:



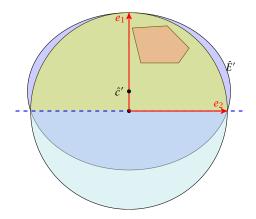


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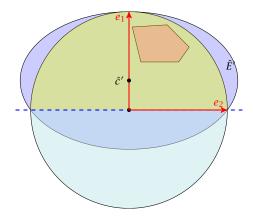


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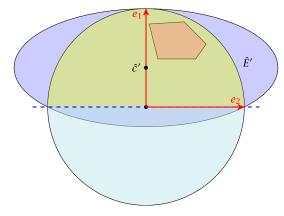


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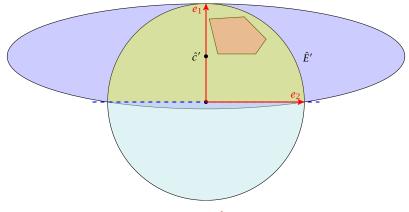


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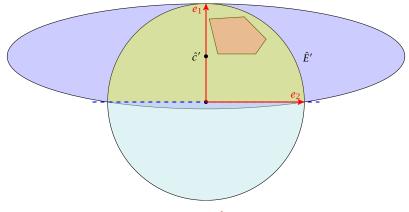


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#### We want to choose t such that the volume of $\hat{E}'$ is minimal.

Lemma 51 Let L be an affine transformation and  $K \subseteq \mathbb{R}^n$ . Then

 $\operatorname{vol}(L(K)) = |\det(L)| \cdot \operatorname{vol}(K)$ .



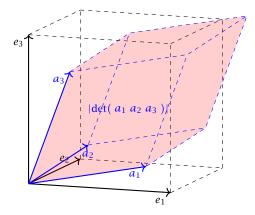
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#### **Lemma 51** Let *L* be an affine transformation and $K \subseteq \mathbb{R}^n$ . Then

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# n-dimensional volume





#### • We want to choose t such that the volume of $\hat{E}'$ is minimal.

 $\operatorname{vol}(\hat{E}') = \operatorname{vol}(B(0,1)) \cdot |\operatorname{det}(\hat{L}')|$  ,



Note that a and b in the above equations depend on t, by the previous equations.



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## $\mathrm{vol}(\hat{E}')$



 $\operatorname{vol}(\hat{E}') = \operatorname{vol}(B(0,1)) \cdot |\operatorname{det}(\hat{L}')|$ 



 $vol(\hat{E}') = vol(B(0,1)) \cdot |det(\hat{L}')|$  $= vol(B(0,1)) \cdot ab^{n-1}$ 



$$vol(\hat{E}') = vol(B(0,1)) \cdot |det(\hat{L}')|$$
  
= vol(B(0,1)) \cdot ab^{n-1}  
= vol(B(0,1)) \cdot (1-t) \cdot \left( \frac{1-t}{\sqrt{1-2t}} \right)^{n-1}



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=  $vol(B(0,1)) \cdot \frac{(1-t)^n}{(\sqrt{1-2t})^{n-1}}$ 

We use the shortcut  $\Phi := \operatorname{vol}(B(0, 1))$ .









$$\frac{\mathrm{d}\operatorname{vol}(\hat{E}')}{\mathrm{d}\,t} = \frac{\mathrm{d}}{\mathrm{d}\,t} \left( \Phi \frac{(1-t)^n}{(\sqrt{1-2t})^{n-1}} \right)$$



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$$= \frac{\Phi}{N^2}$$
$$\boxed{N = \text{denominator}}$$



$$\frac{\mathrm{d}\operatorname{vol}(\hat{E}')}{\mathrm{d}t} = \frac{\mathrm{d}}{\mathrm{d}t} \left( \Phi \frac{(1-t)^n}{(\sqrt{1-2t})^{n-1}} \right)$$
$$= \frac{\Phi}{N^2} \cdot \left( \frac{(-1) \cdot n(1-t)^{n-1}}{(\operatorname{derivative of numerator})^{n-1}} \right)$$



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denominator



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outer derivative



$$\begin{aligned} \frac{\mathrm{d}\operatorname{vol}(\hat{E}')}{\mathrm{d}\,t} &= \frac{\mathrm{d}}{\mathrm{d}\,t} \left( \Phi \frac{(1-t)^n}{(\sqrt{1-2t})^{n-1}} \right) \\ &= \frac{\Phi}{N^2} \cdot \left( (-1) \cdot n(1-t)^{n-1} \cdot (\sqrt{1-2t})^{n-1} - (n-1)(\sqrt{1-2t})^{n-2} \cdot \frac{1}{2\sqrt{1-2t}} \cdot (-2) \right) \\ &\quad \left( \operatorname{inner derivative} \right) \end{aligned}$$



$$\frac{\mathrm{d}\operatorname{vol}(\hat{E}')}{\mathrm{d}t} = \frac{\mathrm{d}}{\mathrm{d}t} \left( \Phi \frac{(1-t)^n}{(\sqrt{1-2t})^{n-1}} \right)$$
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$$\underbrace{\operatorname{numerator}}_{\text{numerator}}$$



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- We obtain the minimum for  $t = \frac{1}{n+1}$ .
- For this value we obtain





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 and  $b =$ 



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#### To see the equation for b, observe that

 $b^2$ 



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$$b^{2} = \frac{(1-t)^{2}}{1-2t} = \frac{(1-\frac{1}{n+1})^{2}}{1-\frac{2}{n+1}}$$



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Let  $\gamma_n = \frac{\operatorname{vol}(\hat{E}')}{\operatorname{vol}(B(0,1))} = ab^{n-1}$  be the ratio by which the volume changes:

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$$\begin{split} \gamma_n^2 &= \Big(\frac{n}{n+1}\Big)^2 \Big(\frac{n^2}{n^2-1}\Big)^{n-1} \\ &= \Big(1 - \frac{1}{n+1}\Big)^2 \Big(1 + \frac{1}{(n-1)(n+1)}\Big)^{n-1} \\ &\le e^{-2\frac{1}{n+1}} \cdot e^{\frac{1}{n+1}} \\ &= e^{-\frac{1}{n+1}} \end{split}$$



Let  $\gamma_n = \frac{\operatorname{vol}(\hat{E}')}{\operatorname{vol}(B(0,1))} = ab^{n-1}$  be the ratio by which the volume changes:

$$y_n^2 = \left(\frac{n}{n+1}\right)^2 \left(\frac{n^2}{n^2 - 1}\right)^{n-1}$$
  
=  $\left(1 - \frac{1}{n+1}\right)^2 \left(1 + \frac{1}{(n-1)(n+1)}\right)^{n-1}$   
 $\leq e^{-2\frac{1}{n+1}} \cdot e^{\frac{1}{n+1}}$   
=  $e^{-\frac{1}{n+1}}$ 

where we used  $(1 + x)^a \le e^{ax}$  for  $x \in \mathbb{R}$  and a > 0.



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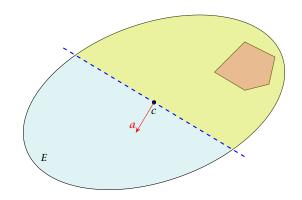
$$y_n^2 = \left(\frac{n}{n+1}\right)^2 \left(\frac{n^2}{n^2 - 1}\right)^{n-1}$$
  
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where we used  $(1 + x)^a \le e^{ax}$  for  $x \in \mathbb{R}$  and a > 0.

This gives  $\gamma_n \leq e^{-\frac{1}{2(n+1)}}$ .

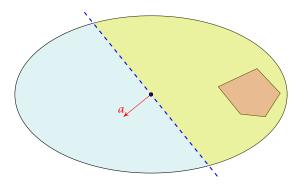


## How to Compute the New Ellipsoid



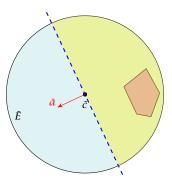


• Use  $f^{-1}$  (recall that f = Lx + t is the affine transformation of the unit ball) to rotate/distort the ellipsoid (back) into the unit ball.



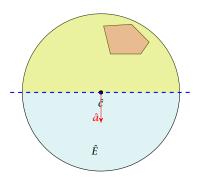


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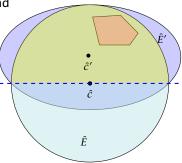


- Use  $f^{-1}$  (recall that f = Lx + t is the affine transformation of the unit ball) to rotate/distort the ellipsoid (back) into the unit ball.
- ► Use a rotation R<sup>-1</sup> to rotate the unit ball such that the normal vector of the halfspace is parallel to e<sub>1</sub>.



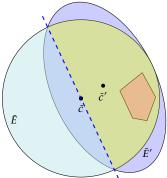


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- Compute the new center ĉ' and the new matrix Q̂' for this simplified setting.



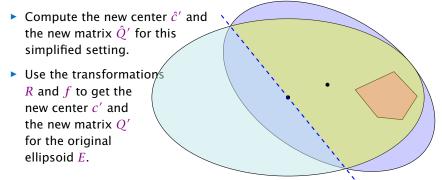


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- Use the transformations *R* and *f* to get the new center *c'* and the new matrix *Q'* for the original ellipsoid *E*.



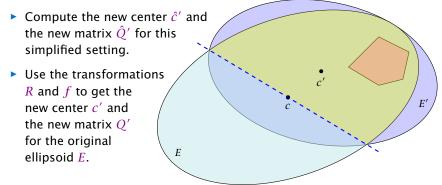


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$$e^{-\frac{1}{2(n+1)}}$$



$$e^{-\frac{1}{2(n+1)}} \ge \frac{\operatorname{vol}(\hat{E}')}{\operatorname{vol}(B(0,1))}$$



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Here it is important that mapping a set with affine function f(x) = Lx + t changes the volume by factor det(*L*).



How to Compute The New Parameters?



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This means  $\bar{a} = L^T a$ .



After rotating back (applying  $R^{-1}$ ) the normal vector of the halfspace points in negative  $x_1$ -direction. Hence,

$$R^{-1}\left(\frac{L^{T}a}{\|L^{T}a\|}\right) = -e_{1} \quad \Rightarrow \quad -\frac{L^{T}a}{\|L^{T}a\|} = R \cdot e_{1}$$

 $\bar{c}'$ 

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$$= -\frac{1}{n+1}L\frac{L^{T}a}{\|L^{T}a\|} + c$$
$$= c - \frac{1}{n+1}\frac{Qa}{\sqrt{a^{T}Qa}}$$

For computing the matrix Q' of the new ellipsoid we assume in the following that  $\hat{E}', \bar{E}'$  and E' refer to the ellipsoids centered in the origin.



### Recall that

$$\hat{Q}' = \begin{pmatrix} a^2 & 0 & \dots & 0 \\ 0 & b^2 & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ 0 & \dots & 0 & b^2 \end{pmatrix}$$

This gives

$$\hat{Q}' = \frac{n^2}{n^2 - 1} \left( I - \frac{2}{n+1} e_1 e_1^T \right)$$

because for  $a^2 = n^2/(n+1)^2$  and  $b^2 = n^2/n^2-1$ 

$$\hat{Q}' = \begin{pmatrix} a^2 & 0 & \dots & 0 \\ 0 & b^2 & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ 0 & \dots & 0 & b^2 \end{pmatrix}$$

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$$\hat{Q}' = \frac{n^2}{n^2 - 1} \left( I - \frac{2}{n+1} e_1 e_1^T \right)$$

$$b^{2} - b^{2} \frac{2}{n+1} = \frac{n^{2}}{n^{2}-1} - \frac{2n^{2}}{(n-1)(n+1)^{2}}$$
$$= \frac{n^{2}(n+1) - 2n^{2}}{(n-1)(n+1)^{2}} - \frac{n^{2}(n-1)}{(n-1)(n+1)^{2}} = a^{2}$$

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 $\bar{E}'$ 



 $\bar{E}' = R(\hat{E}')$ 



$$\bar{E}' = R(\hat{E}')$$

$$= \{R(x) \mid x^T \hat{Q'}^{-1} x \le 1\}$$



$$\bar{E}' = R(\hat{E}')$$
  
= {R(x) |  $x^T \hat{Q}'^{-1} x \le 1$ }  
= { $y | (R^{-1}y)^T \hat{Q}'^{-1} R^{-1} y \le 1$ }



$$\begin{split} \bar{E}' &= R(\hat{E}') \\ &= \{ R(x) \mid x^T \hat{Q'}^{-1} x \le 1 \} \\ &= \{ \gamma \mid (R^{-1} \gamma)^T \hat{Q'}^{-1} R^{-1} \gamma \le 1 \} \\ &= \{ \gamma \mid \gamma^T (R^T)^{-1} \hat{Q'}^{-1} R^{-1} \gamma \le 1 \} \end{split}$$



$$\begin{split} \bar{E}' &= R(\hat{E}') \\ &= \{ R(x) \mid x^T \hat{Q}'^{-1} x \le 1 \} \\ &= \{ y \mid (R^{-1} y)^T \hat{Q}'^{-1} R^{-1} y \le 1 \} \\ &= \{ y \mid y^T (R^T)^{-1} \hat{Q}'^{-1} R^{-1} y \le 1 \} \\ &= \{ y \mid y^T (\underline{R} \hat{Q}' R^T)^{-1} y \le 1 \} \\ &= \{ y \mid y^T (\underline{R} \hat{Q}' R^T)^{-1} y \le 1 \} \end{split}$$



Hence,

 $\bar{Q}'$ 



Hence,

 $\bar{Q}' = R\hat{Q}'R^T$ 



$$\bar{Q}' = R\hat{Q}'R^T$$
$$= R \cdot \frac{n^2}{n^2 - 1} \left(I - \frac{2}{n+1}e_1e_1^T\right) \cdot R^T$$



1

$$\begin{split} \bar{Q}' &= R\hat{Q}'R^T \\ &= R \cdot \frac{n^2}{n^2 - 1} \Big( I - \frac{2}{n+1} e_1 e_1^T \Big) \cdot R^T \\ &= \frac{n^2}{n^2 - 1} \Big( R \cdot R^T - \frac{2}{n+1} (Re_1) (Re_1)^T \Big) \end{split}$$



$$\begin{split} \bar{Q}' &= R\hat{Q}'R^T \\ &= R \cdot \frac{n^2}{n^2 - 1} \Big( I - \frac{2}{n+1} e_1 e_1^T \Big) \cdot R^T \\ &= \frac{n^2}{n^2 - 1} \Big( R \cdot R^T - \frac{2}{n+1} (Re_1) (Re_1)^T \Big) \\ &= \frac{n^2}{n^2 - 1} \Big( I - \frac{2}{n+1} \frac{L^T a a^T L}{\|L^T a\|^2} \Big) \end{split}$$



E'



 $E' = L(\bar{E}')$ 



$$E' = L(\bar{E}') = \{L(x) \mid x^T \bar{Q}'^{-1} x \le 1\}$$



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= { $y$  |  $(L^{-1}y)^T \bar{Q}'^{-1} L^{-1} y \le 1$ }  
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= { $y$  |  $y^T (\underline{L} \bar{Q}' L^T)^{-1} y \le 1$ }



Hence,

Q'



Hence,

 $Q' = L\bar{Q}'L^T$ 



$$Q' = L\bar{Q}'L^T$$
$$= L \cdot \frac{n^2}{n^2 - 1} \left(I - \frac{2}{n+1} \frac{L^T a a^T L}{a^T Q a}\right) \cdot L^T$$



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$$= \frac{n^{2}}{n^{2}-1} \left(Q - \frac{2}{n+1} \frac{Qaa^{T}Q}{a^{T}Qa}\right)$$



#### **Incomplete Algorithm**

#### Algorithm 1 ellipsoid-algorithm

- 1: **input:** point  $c \in \mathbb{R}^n$ , convex set  $K \subseteq \mathbb{R}^n$
- 2: **output:** point  $x \in K$  or "K is empty"
- 3: *Q* ← ???

4: repeat

5: **if** 
$$c \in K$$
 **then return**  $c$ 

6: else

7: choose a violated hyperplane *a* 

8: 
$$c \leftarrow c - \frac{1}{n+1} \frac{Qa}{\sqrt{a^T Qa}}$$

9: 
$$Q \leftarrow \frac{n^2}{n^2 - 1} \Big( Q - \frac{2}{n+1} \frac{Qaa^T Q}{a^T Qaa} \Big)$$

10: **endif** 

11: until ???

12: return "K is empty"

#### **Repeat: Size of basic solutions**

#### Lemma 52

Let  $P = \{x \in \mathbb{R}^n \mid Ax \le b\}$  be a bounded polyhedron. Let  $\langle a_{\max} \rangle$  be the maximum encoding length of an entry in A, b. Then every entry  $x_j$  in a basic solution fulfills  $|x_j| = \frac{D_j}{D}$  with  $D_j, D \le 2^{2n\langle a_{\max} \rangle + 2n\log_2 n}$ .

In the following we use  $\delta := 2^{2n(a_{\max})+2n\log_2 n}$ .

Note that here we have  $P = \{x \mid Ax \le b\}$ . The previous lemmas we had about the size of feasible solutions were slightly different as they were for different polytopes.



#### **Repeat: Size of basic solutions**

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#### **Repeat: Size of basic solutions**

**Proof:** Let  $\overline{A} = \begin{bmatrix} A & -A & I_m \end{bmatrix}$ , *b*, be the matrix and right-hand vector after transforming the system to standard form.

The determinant of the matrices  $\bar{A}_B$  and  $\bar{M}_j$  (matrix obt. when replacing the *j*-th column of  $\bar{A}_B$  by *b*) can become at most

 $\begin{aligned} \det(\bar{A}_B), \det(\bar{M}_j) &\leq \|\vec{\ell}_{\max}\|^{2n} \\ &\leq (\sqrt{2n} \cdot 2^{\langle a_{\max} \rangle})^{2n} \leq 2^{2n \langle a_{\max} \rangle + 2n \log_2 n} \end{aligned}$ 

where  $\vec{\ell}_{max}$  is the longest column-vector that can be obtained after deleting all but 2n rows and columns from  $\bar{A}$ .

This holds because columns from  $I_m$  selected when going from  $\overline{A}$  to  $\overline{A}_B$  do not increase the determinant. Only the at most 2n columns from matrices A and -A that  $\overline{A}$  consists of contribute.

For feasibility checking we can assume that the polytop P is bounded; it is sufficient to consider basic solutions.

Every entry  $x_i$  in a basic solution fulfills  $|x_i| \le \delta$ .

Hence, *P* is contained in the cube  $-\delta \le x_i \le \delta$ .

A vector in this cube has at most distance  $R:=\sqrt{n}\delta$  from the origin.



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# How do we find the first ellipsoid?

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Starting with the ball  $E_0 := B(0, R)$  ensures that P is completely contained in the initial ellipsoid. This ellipsoid has volume at most  $R^n \operatorname{vol}(B(0, 1)) \le (n\delta)^n \operatorname{vol}(B(0, 1))$ .



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# When can we terminate?

Let  $P := \{x \mid Ax \le b\}$  with  $A \in \mathbb{Z}$  and  $b \in \mathbb{Z}$  be a bounded polytop. Let  $\langle a_{\max} \rangle$  be the encoding length of the largest entry in A or b.

Consider the following polyhedron

$$P_{\lambda} := \left\{ x \mid Ax \le b + \frac{1}{\lambda} \begin{pmatrix} 1 \\ \vdots \\ 1 \end{pmatrix} \right\} ,$$

where  $\lambda = \delta^2 + 1$ .



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# **Lemma 53** $P_{\lambda}$ is feasible if and only if P is feasible.

 $\Leftarrow$ : obvious!



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←: obvious!



 $\Rightarrow$ :

Consider the polyhedrons

$$\bar{P} = \left\{ x \mid \left[ A - A I_m \right] x = b; x \ge 0 \right\}$$

and

$$ar{P}_{\lambda} = \left\{ x \mid \left[ A - A I_m \right] x = b + rac{1}{\lambda} \begin{pmatrix} 1 \\ \vdots \\ 1 \end{pmatrix}; x \ge 0 
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P is feasible if and only if  $\bar{P}$  is feasible, and  $P_{\lambda}$  feasible if and only if  $\bar{P}_{\lambda}$  feasible.

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Let 
$$\overline{A} = \begin{bmatrix} A & -A & I_m \end{bmatrix}$$
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 $\bar{P}_{\lambda}$  feasible implies that there is a basic feasible solution represented by

$$\chi_B = \bar{A}_B^{-1}b + \frac{1}{\lambda}\bar{A}_B^{-1}\begin{pmatrix}1\\\vdots\\1\end{pmatrix}$$

(The other *x*-values are zero)

The only reason that this basic feasible solution is not feasible for  $\overline{P}$  is that one of the basic variables becomes negative.

Hence, there exists i with

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and

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where  $\bar{M}_j$  is obtained by replacing the *j*-th column of  $\bar{A}_B$  by  $\vec{1}$ .

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If  $P_{\lambda}$  is feasible then it contains a ball of radius  $r := 1/\delta^3$ . This has a volume of at least  $r^n \operatorname{vol}(B(0,1)) = \frac{1}{\delta^{3n}} \operatorname{vol}(B(0,1))$ .



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## Proof:

If  $P_{\lambda}$  feasible then also *P*. Let *x* be feasible for *P*.



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 $\le b_i + \frac{\sqrt{n} \cdot 2^{\langle a_{\max} \rangle}}{\delta^3} \le b_i + \frac{1}{\delta^2 + 1} \le b_i + \frac{1}{\lambda}$ 



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 $\le b_i + \frac{\sqrt{n} \cdot 2^{\langle a_{\max} \rangle}}{\delta^3} \le b_i + \frac{1}{\delta^2 + 1} \le b_i + \frac{1}{\lambda}$ 

Hence,  $x + \vec{\ell}$  is feasible for  $P_{\lambda}$  which proves the lemma.







$$e^{-\frac{i}{2(n+1)}} \cdot \operatorname{vol}(B(0,R)) < \operatorname{vol}(B(0,r))$$



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Hence,

i



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Hence,

$$i > 2(n+1)\ln\left(\frac{\operatorname{vol}(B(0,R))}{\operatorname{vol}(B(0,r))}\right)$$



$$e^{-\frac{i}{2(n+1)}} \cdot \operatorname{vol}(B(0,R)) < \operatorname{vol}(B(0,r))$$

Hence,

$$\begin{split} i &> 2(n+1) \ln \left( \frac{\operatorname{vol}(B(0,R))}{\operatorname{vol}(B(0,r))} \right) \\ &= 2(n+1) \ln \left( n^n \delta^n \cdot \delta^{3n} \right) \end{split}$$



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How many iterations do we need until the volume becomes too small?

$$e^{-\frac{i}{2(n+1)}} \cdot \operatorname{vol}(B(0,R)) < \operatorname{vol}(B(0,r))$$

Hence,

$$\begin{split} i &> 2(n+1) \ln \left( \frac{\operatorname{vol}(B(0,R))}{\operatorname{vol}(B(0,r))} \right) \\ &= 2(n+1) \ln \left( n^n \delta^n \cdot \delta^{3n} \right) \\ &= 8n(n+1) \ln(\delta) + 2(n+1)n \ln(n) \\ &= \mathcal{O}(\operatorname{poly}(n, \langle a_{\max} \rangle)) \end{split}$$



### Algorithm 1 ellipsoid-algorithm

1: **input:** point  $c \in \mathbb{R}^n$ , convex set  $K \subseteq \mathbb{R}^n$ , radii *R* and *r* 

- 2: with  $K \subseteq B(c, R)$ , and  $B(x, r) \subseteq K$  for some x
- 3: **output:** point  $x \in K$  or "K is empty"

4: 
$$Q \leftarrow \operatorname{diag}(R^2, \dots, R^2) // \text{ i.e., } L = \operatorname{diag}(R, \dots, R)$$

5: repeat

6: **if** 
$$c \in K$$
 **then return**  $c$ 

С

7: else

- 8: choose a violated hyperplane *a*
- 9:

$$\leftarrow c - \frac{1}{n+1} \frac{Qa}{\sqrt{a^T Qa}}$$

10: 
$$Q \leftarrow \frac{n^2}{n^2 - 1} \left( Q - \frac{2}{n+1} \frac{Qaa^T Q}{a^T Qaa} \right)$$

### 11: endif

12: **until** 
$$det(Q) \le r^{2n} // i.e., det(L) \le r^n$$

13: return "K is empty"

Let  $K \subseteq \mathbb{R}^n$  be a convex set. A separation oracle for K is an algorithm A that gets as input a point  $x \in \mathbb{R}^n$  and either

- certifies that  $x \in K$ ,
- or finds a hyperplane separating x from K.

We will usually assume that A is a polynomial-time algorithm.

In order to find a point in K we need

- a guarantee that a ball of radius  $\sim$  is contained in % ,
- $\sim$  an initial ball  $\beta(c, d)$  with radius  $\beta$  that contains  $\beta_{c}$
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- a guarantee that a ball of radius  $\sigma$  is contained in  $\mathbb{K}_{r}$
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We will usually assume that A is a polynomial-time algorithm.

### In order to find a point in K we need

- a guarantee that a ball of radius  $\sim$  is contained in  $\%_{0}$
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- » a separation oracle for %...



Let  $K \subseteq \mathbb{R}^n$  be a convex set. A separation oracle for K is an algorithm A that gets as input a point  $x \in \mathbb{R}^n$  and either

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- inequalities  $Ax \leq b$ ;  $m \times n$  matrix A with rows  $a_i^T$
- $P = \{x \mid Ax \le b\}; P^{\circ} := \{x \mid Ax < b\}$
- interior point algorithm:  $x \in P^\circ$  throughout the algorithm
- for  $x \in P^\circ$  define

$$s_i(x) := b_i - a_i^T x$$

as the slack of the *i*-th constraint

logarithmic barrier function:

$$\phi(x) = -\sum_{i=1}^m \log(s_i(x))$$

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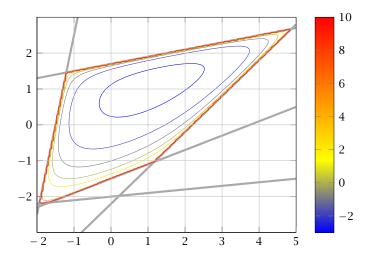
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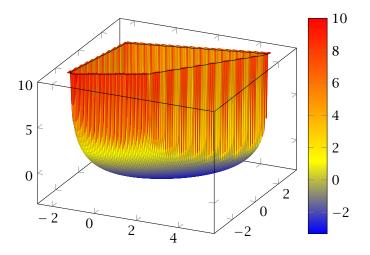
## **Penalty Function**





10 Karmarkars Algorithm

## **Penalty Function**





## **Gradient and Hessian**

### **Taylor approximation:**

$$\phi(x+\epsilon) \approx \phi(x) + \nabla \phi(x)^T \epsilon + \frac{1}{2} \epsilon^T \nabla^2 \phi(x) \epsilon$$

Gradient:

$$\nabla \phi(x) = \sum_{i=1}^{m} \frac{1}{s_i(x)} \cdot a_i = A^T d_x$$

where  $d_x^T = (1/s_1(x), \dots, 1/s_m(x))$ . ( $d_x$  vector of inverse slacks)

Hessian:

$$H_{\mathbf{x}} := \nabla^2 \phi(\mathbf{x}) = \sum_{i=1}^m \frac{1}{s_i(\mathbf{x})^2} a_i a_i^T = A^T D_{\mathbf{x}}^2 A_i$$

with  $D_x = \operatorname{diag}(d_x)$ .

## **Gradient and Hessian**

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# **Proof for Gradient**

$$\begin{split} \frac{\partial \phi(x)}{\partial x_i} &= \frac{\partial}{\partial x_i} \left( -\sum_r \ln(s_r(x)) \right) \\ &= -\sum_r \frac{\partial}{\partial x_i} \left( \ln(s_r(x)) \right) = -\sum_r \frac{1}{s_r(x)} \frac{\partial}{\partial x_i} \left( s_r(x) \right) \\ &= -\sum_r \frac{1}{s_r(x)} \frac{\partial}{\partial x_i} \left( b_r - a_r^T x \right) = \sum_r \frac{1}{s_r(x)} \frac{\partial}{\partial x_i} \left( a_r^T x \right) \\ &= \sum_r \frac{1}{s_r(x)} A_{ri} \end{split}$$

The *i*-th entry of the gradient vector is  $\sum_{r} 1/s_r(x) \cdot A_{ri}$ . This gives that the gradient is

$$\nabla \phi(x) = \sum_{r} 1/s_{r}(x)a_{r} = A^{T}d_{x}$$

## **Proof for Hessian**

$$\frac{\partial}{\partial x_j} \left( \sum_r \frac{1}{s_r(x)} A_{ri} \right) = \sum_r A_{ri} \left( -\frac{1}{s_r(x)^2} \right) \cdot \frac{\partial}{\partial x_j} \left( s_r(x) \right)$$
$$= \sum_r A_{ri} \frac{1}{s_r(x)^2} A_{rj}$$

Note that  $\sum_{r} A_{ri}A_{rj} = (A^{T}A)_{ij}$ . Adding the additional factors  $1/s_{r}(x)^{2}$  can be done with a diagonal matrix.

Hence the Hessian is

$$H_X = A^T D^2 A$$

 $H_X$  is positive semi-definite for  $x \in P^\circ$ 

 $u^{T}H_{x}u = u^{T}A^{T}D_{x}^{2}Au = ||D_{x}Au||_{2}^{2} \ge 0$ 

This gives that  $\phi(x)$  is convex.

If rank(A) = n,  $H_x$  is positive definite for  $x \in P^\circ$  $u^T H_x u = \|D_x A u\|_2^2 > 0$  for  $u \neq 0$ 

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 $E_{x} = \{ y \mid (y - x)^{T} H_{x} (y - x) \leq 1 \} = \{ y \mid ||y - x||_{H_{x}} \leq 1 \}$ 

Points in Ex are feasible!!!

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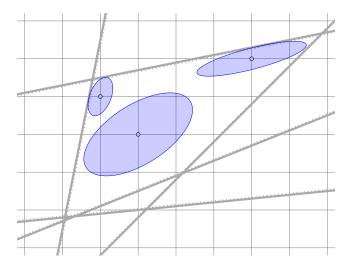
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10 Karmarkars Algorithm

# **Analytic Center**

 $x_{\mathrm{ac}} := \operatorname{arg\,min}_{x \in P^\circ} \phi(x)$ 

•  $x_{ac}$  is solution to

$$\nabla \phi(x) = \sum_{i=1}^m \frac{1}{s_i(x)} a_i = 0$$

- depends on the description of the polytope
- $x_{ac}$  exists and is unique iff  $P^{\circ}$  is nonempty and bounded



# **Central Path**

In the following we assume that the LP and its dual are strictly feasible and that rank(A) = n.

```
Central Path:
Set of points \{x^*(t) \mid t > 0\} with
```

 $x^*(t) = \operatorname{argmin}_x \{ t c^T x + \phi(x) \}$ 

- t = 0: analytic center
- $t = \infty$ : optimum solution

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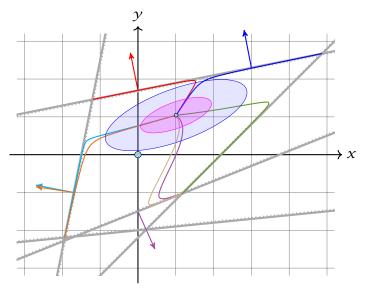
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# **Different Central Paths**





10 Karmarkars Algorithm

# **Central Path**

#### Intuitive Idea:

Find point on central path for large value of t. Should be close to optimum solution.

#### Questions:

- Is this really true? How large a t do we need?
- How do we find corresponding point  $x^*(t)$  on central path?



# The Dual

primal-dual pair:

#### Assumptions

- primal and dual problems are strictly feasible;
- $\operatorname{rank}(A) = n$ .

# **Force Field Interpretation**

Point  $x^*(t)$  on central path is solution to  $tc + \nabla \phi(x) = 0$ 

- We can view each constraint as generating a repelling force. The combination of these forces is represented by ∇φ(x).
- In addition there is a force tc pulling us towards the optimum solution.



Point  $x^*(t)$  on central path is solution to  $tc + \nabla \phi(x) = 0$ .

This means

$$tc + \sum_{i=1}^{m} \frac{1}{s_i(x^*(t))} a_i = 0$$

or

$$c + \sum_{i=1}^{m} z_i^*(t) a_i = 0$$
 with  $z_i^*(t) = \frac{1}{t s_i(x^*(t))}$ 

2000 is strictly dual feasible: (2022-5-0-0) 225-0)

if gap is less than  $1/2^{1010}$  we can snap to optimum point:

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*z*\*(*t*) is strictly dual feasible: (*A<sup>T</sup>z*\* + *c* = 0; *z*\* > 0)
 duality gap between *x* := *x*\*(*t*) and *z* := *z*\*(*t*) is

$$c^T x + b^T z = (b - Ax)^T z = \frac{m}{t}$$

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# How to find $x^*(t)$

#### First idea:

- start somewhere in the polytope
- use iterative method (Newtons method) to minimize  $f_t(x) := tc^T x + \phi(x)$



Quadratic approximation of  $f_t$ 

$$f_t(x + \epsilon) \approx f_t(x) + \nabla f_t(x)^T \epsilon + \frac{1}{2} \epsilon^T H_{f_t}(x) \epsilon$$

Suppose this were exact:

$$f_t(x + \epsilon) = f_t(x) + \nabla f_t(x)^T \epsilon + \frac{1}{2} \epsilon^T H_{f_t}(x) \epsilon$$

Then gradient is given by:

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10 Karmarkars Algorithm

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We want to move to a point where this gradient is 0:

**Newton Step** at  $x \in P^{\circ}$ 

$$\Delta x_{\mathsf{nt}} = -H_{f_t}^{-1}(x)\nabla f_t(x)$$
  
=  $-H_{f_t}^{-1}(x)(tc + \nabla \phi(x))$   
=  $-(A^T D_x^2 A)^{-1}(tc + A^T d_x)$ 

**Newton Iteration:** 

$$x := x + \Delta x_{nt}$$

## **Measuring Progress of Newton Step**

Newton decrement:

 $\lambda_t(x) = \|D_x A \Delta x_{\mathsf{nt}}\| \\ = \|\Delta x_{\mathsf{nt}}\|_{H_x}$ 

Square of Newton decrement is linear estimate of reduction if we do a Newton step:

 $-\lambda_t(x)^2 = \nabla f_t(x)^T \Delta x_{\mathsf{nt}}$ 

•  $\lambda_t(x) = 0$  iff  $x = x^*(t)$ 

•  $\lambda_t(x)$  is measure of proximity of x to  $x^*(t)$ 

## **Measuring Progress of Newton Step**

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•  $\lambda_t(x)$  is measure of proximity of x to  $x^*(t)$ 

## **Measuring Progress of Newton Step**

Newton decrement:

 $\lambda_t(x) = \|D_x A \Delta x_{\mathsf{nt}}\| \\ = \|\Delta x_{\mathsf{nt}}\|_{H_x}$ 

Square of Newton decrement is linear estimate of reduction if we do a Newton step:

 $-\lambda_t(x)^2 = \nabla f_t(x)^T \Delta x_{\mathsf{nt}}$ 

•  $\lambda_t(x) = 0$  iff  $x = x^*(t)$ 

•  $\lambda_t(x)$  is measure of proximity of x to  $x^*(t)$ 

## Theorem 55

If  $\lambda_t(x) < 1$  then

- $x_+ := x + \Delta x_{nt} \in P^\circ$  (new point feasible)
- $\lambda_t(x_+) \leq \lambda_t(x)^2$

This means we have quadratic convergence. Very fast.

#### feasibility:

►  $\lambda_t(x) = \|\Delta x_{nt}\|_{H_x} < 1$ ; hence  $x_+$  lies in the Dikin ellipsoid around x.

bound on  $\lambda_t(x^+)$ : we use  $D := D_x = \text{diag}(d_x)$  and  $D_+ := D_{x^+} = \text{diag}(d_{x^+})$ 

To see the last equality we use Pythagoras

 $||a||^2 + ||a + b||^2 = ||b||^2$ 

bound on  $\lambda_t(x^+)$ : we use  $D := D_x = \operatorname{diag}(d_x)$  and  $D_+ := D_{x^+} = \operatorname{diag}(d_{x^+})$ 

 $\lambda_{t}(\boldsymbol{x}^{+})^{2} = \|D_{+}A \Delta x_{nt}^{+}\|^{2}$  $\leq \|D_{+}A \Delta x_{nt}^{+}\|^{2} + \|D_{+}A \Delta x_{nt}^{+} + (I - D_{+}^{-1}D)DA \Delta x_{nt}\|^{2}$  $= \|(I - D_{+}^{-1}D)DA \Delta x_{nt}\|^{2}$ 

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$$= \|(I - D_+^{-1}D) D A \Delta x_{\mathsf{nt}}\|^2$$

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 $||a||^2 + ||a + b||^2 = ||b||^2$ 

 $DA\Delta x_{\mathsf{nt}} = DA(x^+ - x)$  $= D(b - Ax - (b - Ax^+))$  $= D(D^{-1}\vec{1} - D^{-1}_{+}\vec{1})$  $= (I - D^{-1}_{+}D)\vec{1}$ 

$$a^{T}(a+b)$$

$$= \Delta x_{\mathsf{nt}}^{+T} A^{T} D_{+} \left( D_{+} A \Delta x_{\mathsf{nt}}^{+} + (I - D_{+}^{-1} D) D A \Delta x_{\mathsf{nt}} \right)$$

$$= \Delta x_{\mathsf{nt}}^{+T} \left( A^{T} D_{+}^{2} A \Delta x_{\mathsf{nt}}^{+} - A^{T} D^{2} A \Delta x_{\mathsf{nt}} + A^{T} D_{+} D A \Delta x_{\mathsf{nt}} \right)$$

$$= \Delta x_{\mathsf{nt}}^{+T} \left( H_{+} \Delta x_{\mathsf{nt}}^{+} - H \Delta x_{\mathsf{nt}} + A^{T} D_{+} \vec{1} - A^{T} D \vec{1} \right)$$

$$= \Delta x_{\mathsf{nt}}^{+T} \left( -\nabla f_{t}(x^{+}) + \nabla f_{t}(x) + \nabla \phi(x^{+}) - \nabla \phi(x) \right)$$

$$= 0$$

 $DA\Delta x_{nt} = DA(x^{+} - x)$ =  $D(b - Ax - (b - Ax^{+}))$ =  $D(D^{-1}\vec{1} - D^{-1}_{+}\vec{1})$ =  $(I - D^{-1}_{+}D)\vec{1}$ 

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bound on  $\lambda_t(x^+)$ : we use  $D := D_x = \operatorname{diag}(d_x)$  and  $D_+ := D_{x^+} = \operatorname{diag}(d_{x^+})$ 

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The second inequality follows from  $\sum_i y_i^4 \leq \left(\sum_i y_i^2
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The second inequality follows from  $\sum_i y_i^4 \le (\sum_i y_i^2)^2$ 

If  $\lambda_t(x)$  is large we do not have a guarantee.

Try to avoid this case!!!



## **Path-following Methods**

Try to slowly travel along the central path.

Algorithm 1 PathFollowing

- 1: start at analytic center
- 2: while solution not good enough do
- 3: make step to improve objective function
- 4: recenter to return to central path

#### simplifying assumptions:

- a first central point  $x^*(t_0)$  is given
- $x^*(t)$  is computed exactly in each iteration

#### $\epsilon$ is approximation we are aiming for

start at  $t = t_0$ , repeat until  $m/t \le \epsilon$ 

- compute  $x^*(\mu t)$  using Newton starting from  $x^*(t)$
- ► *t* := *µt*

where  $\mu = 1 + 1/(2\sqrt{m})$ 

gradient of  $f_{t^+}$  at ( $x = x^*(t)$ )

$$\nabla f_{t^+}(x) = \nabla f_t(x) + (\mu - 1)tc$$
$$= -(\mu - 1)A^T D_X \vec{1}$$

This holds because  $0 = \nabla f_t(x) = tc + A^T D_x \vec{1}$ .

The Newton decrement is

$$\begin{split} \lambda_{t^{+}}(x)^{2} &= \nabla f_{t^{+}}(x)^{T} H^{-1} \nabla f_{t^{+}}(x) \\ &= (\mu - 1)^{2} \vec{1}^{T} B (B^{T} B)^{-1} B^{T} \vec{1} \qquad B = D_{x}^{T} A \\ &\leq (\mu - 1)^{2} m \\ &= 1/4 \end{split}$$

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## **Number of Iterations**

the number of Newton iterations per outer iteration is very small; in practise only 1 or 2

#### Number of outer iterations:

We need  $t_k = \mu^k t_0 \ge m/\epsilon$ . This holds when

 $k \geq \frac{\log(m/(\epsilon t_0))}{\log(\mu)}$ 

We get a bound of

$$\mathcal{O}\left(\sqrt{m}\log\frac{m}{\epsilon t_0}\right)$$

We show how to get a starting point with  $t_0 = 1/2^L$ . Together with  $\epsilon \approx 2^{-L}$  we get  $\mathcal{O}(L\sqrt{m})$  iterations.

EADS II Harald Räcke

For  $x \in P^{\circ}$  and direction  $v \neq 0$  define

$$\sigma_{x}(v) := \max_{i} \frac{a_{i}^{T} v}{s_{i}(x)}$$

**Observation:** 

$$x + \alpha v \in P$$
 for  $\alpha \in \{0, 1/\sigma_x(v)\}$ 



Suppose that we move from x to  $x + \alpha v$ . The linear estimate says that  $f_t(x)$  should change by  $\nabla f_t(x)^T \alpha v$ .

The following argument shows that  $f_t$  is well behaved. For small  $\alpha$  the reduction of  $f_t(x)$  is close to linear estimate.

 $f_t(x+\alpha v)-f_t(x)=tc^T\alpha v+\phi(x+\alpha v)-\phi(x)$ 

 $\phi(x + \alpha v) - \phi(x)$ 

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 $f_t(x + \alpha v) - f_t(x) = tc^T \alpha v + \phi(x + \alpha v) - \phi(x)$ 

 $\phi(x + \alpha v) - \phi(x)$ 



Suppose that we move from x to  $x + \alpha v$ . The linear estimate says that  $f_t(x)$  should change by  $\nabla f_t(x)^T \alpha v$ .

The following argument shows that  $f_t$  is well behaved. For small  $\alpha$  the reduction of  $f_t(x)$  is close to linear estimate.

 $f_t(x + \alpha v) - f_t(x) = tc^T \alpha v + \phi(x + \alpha v) - \phi(x)$ 

 $\phi(x + \alpha v) - \phi(x) = -\sum_{i} \log(s_i(x + \alpha v)) + \sum_{i} \log(s_i(x))$  $= -\sum_{i} \log(s_i(x + \alpha v)/s_i(x))$  $= -\sum_{i} \log(1 - a_i^T \alpha v/s_i(x))$ 



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Define  $w_i = a_i^T v / s_i(x)$  and  $\sigma = \max_i w_i$ . Then

 $f_t(x + \alpha v) - f_t(x) - \nabla f_t(x)^T \alpha v$ 

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$$\leq -\sum_{w_i > 0} (\alpha w_i + \log(1 - \alpha w_i)) + \sum_{w_i < 0} \frac{\alpha^2 w_i^2}{2}$$

$$\leq -\sum_{w_i > 0} \frac{w_i^2}{\sigma^2} (\alpha \sigma + \log(1 - \alpha \sigma)) + \frac{(\alpha \sigma)^2}{2} \sum_{w_i < 0} \frac{w_i^2}{\sigma^2}$$

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$$\leq -\sum_{i} \frac{w_{i}^{2}}{\sigma^{2}} \left( \alpha \sigma + \log(1 - \alpha \sigma) \right)$$
$$= -\frac{1}{\sigma^{2}} \|v\|_{H_{x}}^{2} \left( \alpha \sigma + \log(1 - \alpha \sigma) \right)$$

Damped Newton Iteration: In a damped Newton step we choose

$$x_{+} = x + \frac{1}{1 + \sigma_x(\Delta x_{\rm nt})} \Delta x_{\rm nt}$$



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#### Theorem:

In a damped Newton step the cost decreases by at least

 $\lambda_t(x) - \log(1 + \lambda_t(x))$ 

**Proof:** The decrease in cost is

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 $\geq \lambda_t(x) - \log(1 + \lambda_t(x))$  $\geq 0.09$ 

#### for $\lambda_t(x) \ge 0.5$

**Centering Algorithm:** Input: precision  $\delta$ ; starting point x**1.** compute  $\Delta x_{nt}$  and  $\lambda_t(x)$ **2.** if  $\lambda_t(x) \le \delta$  return x**3.** set  $x := x + \alpha \Delta x_{nt}$  with

$$\alpha = \begin{cases} \frac{1}{1 + \sigma_x(\Delta x_{\mathsf{nt}})} & \lambda_t \ge 1/2\\ 1 & \text{otw.} \end{cases}$$



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## Centering

### Lemma 56

## The centering algorithm starting at $x_0$ reaches a point with $\lambda_t(x) \le \delta$ after

$$\frac{f_t(x_0) - \min_{\mathcal{Y}} f_t(\mathcal{Y})}{0.09} + \mathcal{O}(\log \log(1/\delta))$$

iterations.

This can be very, very slow...



# Let $P = \{Ax \le b\}$ be our (feasible) polyhedron, and $x_0$ a feasible point.

We change  $b \to b + \frac{1}{\lambda} \cdot \vec{1}$ , where  $L = \langle A \rangle + \langle b \rangle + \langle c \rangle$  (encoding length) and  $\lambda = 2^{2L}$ . Recall that a basis is feasible in the old LP iff it is feasible in the new LP.



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### **Lemma** [without proof] The inverse of a matrix M can be represented with rational numbers that have denominators $z_{ij} = det(M)$ .

For two basis solutions  $x_B$ ,  $x_{\bar{B}}$ , the cost-difference  $c^T x_B - c^T x_{\bar{B}}$ can be represented by a rational number that has denominator  $z = \det(A_B) \cdot \det(A_{\bar{B}}) \cdot \lambda$ .

This means that in the perturbed LP it is sufficient to decrease the duality gap to  $1/2^{4L}$  (i.e.,  $t \approx 2^{4L}$ ). This means the previous analysis essentially also works for the perturbed LP.

For a point x from the polytope (not necessarily BFS) the objective value  $\bar{c}^T x$  is at most  $n2^M 2^L$ , where  $M \leq L$  is the encoding length of the largest entry in  $\bar{c}$ .



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Start at  $x_0$ .

Choose  $\hat{c} := -\nabla \phi(x)$ .

 $x_0 = x^*(1)$  is point on central path for  $\hat{c}$  and t = 1.

You can travel the central path in both directions. Go towards 0 until  $t \approx 1/2^{\Omega(L)}$ . This requires  $O(\sqrt{m}L)$  outer iterations.

Let  $x_{\hat{c}}$  denote this point.

Let  $x_c$  denote the point that minimizes

 $t \cdot c^T x + \phi(x)$ 

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Clearly,

$$t \cdot \hat{c}^T \boldsymbol{x}_{\hat{c}} + \phi(\boldsymbol{x}_{\hat{c}}) \leq t \cdot \hat{c}^T \boldsymbol{x}_{\boldsymbol{c}} + \phi(\boldsymbol{x}_{\boldsymbol{c}})$$

### The different between $f_t(x_{\hat{c}})$ and $f_t(x_c)$ is

 $\begin{aligned} tc^T \boldsymbol{x}_{\hat{c}} + \boldsymbol{\phi}(\boldsymbol{x}_{\hat{c}}) - tc^T \boldsymbol{x}_c - \boldsymbol{\phi}(\boldsymbol{x}_c) \\ &\leq t(c^T \boldsymbol{x}_{\hat{c}} + \hat{c}^T \boldsymbol{x}_c - \hat{c}^T \boldsymbol{x}_{\hat{c}} - c^T \boldsymbol{x}_c) \\ &\leq 4tn 2^{3L} \end{aligned}$ 

For  $t = 1/2^{\Omega(L)}$ ) the last term becomes constant. Hence, using damped Newton we can move from  $x_{\hat{c}}$  to  $x_c$  quickly.

In total for this analysis we require  $\mathcal{O}(\sqrt{mL})$  outer iterations for the whole algorithm.

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## Part III

## **Approximation Algorithms**



- Heuristics.
- Exploit special structure of instances occurring in practise.
- Consider algorithms that do not compute the optimal solution but provide solutions that are close to optimum.



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### **Definition 57**

An  $\alpha$ -approximation for an optimization problem is a polynomial-time algorithm that for all instances of the problem produces a solution whose value is within a factor of  $\alpha$  of the value of an optimal solution.



#### We need algorithms for hard problems.

- It gives a rigorous mathematical base for studying heunstics.
- It provides a metric to compare the difficulty of various optimization problems.
- Proving theorems may give a deeper theoretical understanding which in turn leads to new algorithmic approaches.

### Why not?



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### **Definition 58**

An optimization problem  $P = (\mathcal{I}, \text{sol}, m, \text{goal})$  is in **NPO** if

- $x \in \mathcal{I}$  can be decided in polynomial time
- $y \in sol(\mathcal{I})$  can be verified in polynomial time
- *m* can be computed in polynomial time
- goal  $\in \{\min, \max\}$

In other words: the decision problem is there a solution y with m(x, y) at most/at least z is in NP.



- x is problem instance
- y is candidate solution
- $m^*(x)$  cost/profit of an optimal solution

### **Definition 59 (Performance Ratio)**

$$R(x, y) := \max\left\{\frac{m(x, y)}{m^*(x)}, \frac{m^*(x)}{m(x, y)}\right\}$$



## **Definition 60 (***r***-approximation)**

An algorithm A is an r-approximation algorithm iff

 $\forall x \in \mathcal{I}: R(x, A(x)) \leq r$ ,

and A runs in polynomial time.



## **Definition 61 (PTAS)**

A PTAS for a problem *P* from NPO is an algorithm that takes as input  $x \in \mathcal{I}$  and  $\epsilon > 0$  and produces a solution  $\mathcal{Y}$  for x with

 $R(x,y) \leq 1 + \epsilon$  .

The running time is polynomial in |x|.

approximation with arbitrary good factor... fast?



## Problems that have a PTAS

**Scheduling.** Given m jobs with known processing times; schedule the jobs on n machines such that the MAKESPAN is minimized.



## **Definition 62 (FPTAS)**

An FPTAS for a problem *P* from NPO is an algorithm that takes as input  $x \in \mathcal{I}$  and  $\epsilon > 0$  and produces a solution  $\mathcal{Y}$  for x with

 $R(x,y) \leq 1 + \epsilon$  .

The running time is polynomial in |x| and  $1/\epsilon$ .

approximation with arbitrary good factor... fast!



## Problems that have an FPTAS

**KNAPSACK.** Given a set of items with profits and weights choose a subset of total weight at most W s.t. the profit is maximized.



## **Definition 63 (APX - approximable)**

A problem *P* from NPO is in APX if there exist a constant  $r \ge 1$  and an *r*-approximation algorithm for *P*.

constant factor approximation...



## Problems that are in APX

**MAXCUT.** Given a graph G = (V, E); partition V into two disjoint pieces A and B s.t. the number of edges between both pieces is maximized.

**MAX-3SAT**. Given a 3CNF-formula. Find an assignment to the variables that satisfies the maximum number of clauses.



## Problems with polylogarithmic approximation guarantees

- Set Cover
- Minimum Multicut
- Sparsest Cut
- Minimum Bisection

There is an *r*-approximation with  $r \leq O(\log^{c}(|x|))$  for some constant *c*.

Note that only for some of the above problem a matching lower bound is known.



## There are really difficult problems!

## Theorem 64

For any constant  $\epsilon > 0$  there does not exist an  $\Omega(n^{1-\epsilon})$ -approximation algorithm for the maximum clique problem on a given graph G with n nodes unless P = NP.

Note that an *n*-approximation is trivial.



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## There are weird problems!

Asymmetric *k*-Center admits an  $O(\log^* n)$ -approximation.

There is no  $o(\log^* n)$ -approximation to Asymmetric *k*-Center unless  $NP \subseteq DTIME(n^{\log \log \log n})$ .



Class APX not important in practise.

Instead of saying problem P is in APX one says problem P admits a 4-approximation.

One only says that a problem is APX-hard.



A crucial ingredient for the design and analysis of approximation algorithms is a technique to obtain an upper bound (for maximization problems) or a lower bound (for minimization problems).

Therefore Linear Programs or Integer Linear Programs play a vital role in the design of many approximation algorithms.



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Therefore Linear Programs or Integer Linear Programs play a vital role in the design of many approximation algorithms.



### **Definition 65**

An Integer Linear Program or Integer Program is a Linear Program in which all variables are required to be integral.

## **Definition 66**

A Mixed Integer Program is a Linear Program in which a subset of the variables are required to be integral.



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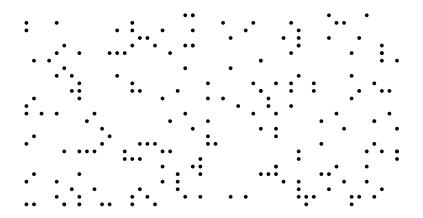
Given a ground set U, a collection of subsets  $S_1, \ldots, S_k \subseteq U$ , where the *i*-th subset  $S_i$  has weight/cost  $w_i$ . Find a collection  $I \subseteq \{1, \ldots, k\}$  such that

 $\forall u \in U \exists i \in I : u \in S_i$  (every element is covered)

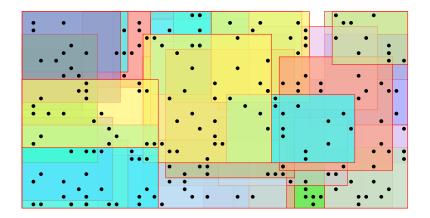
and

$$\sum_{i\in I} w_i$$
 is minimized.

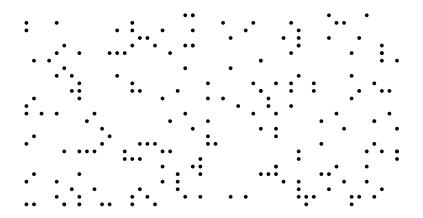




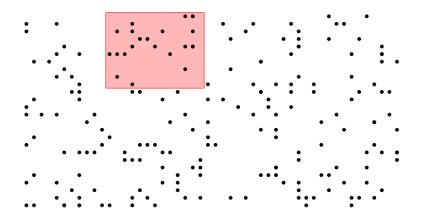




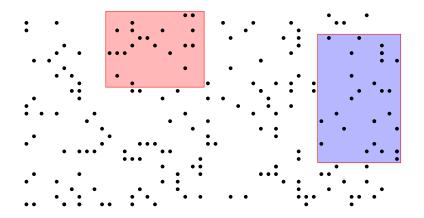




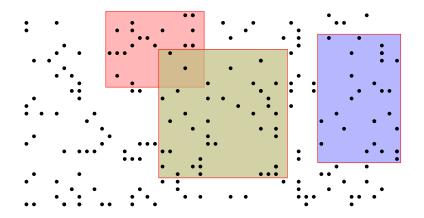




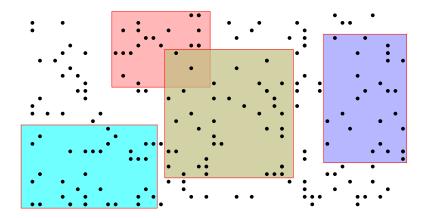




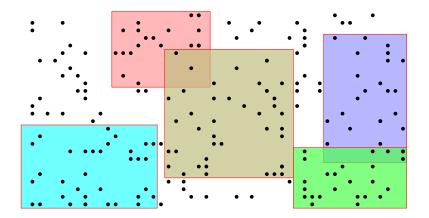




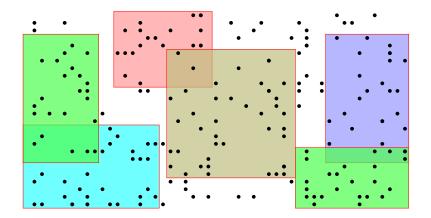




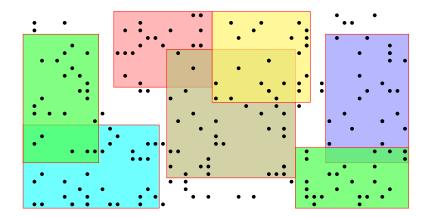




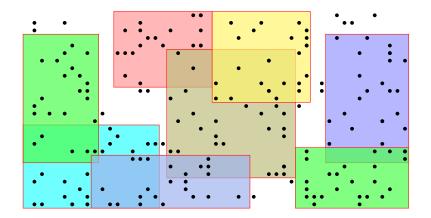




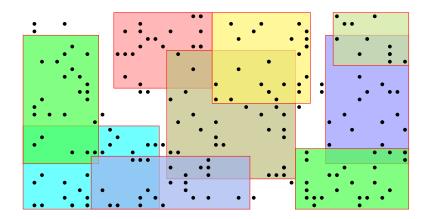




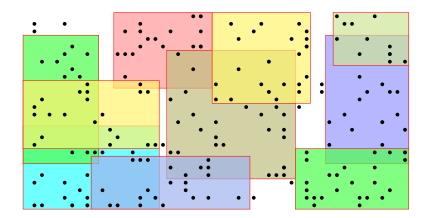




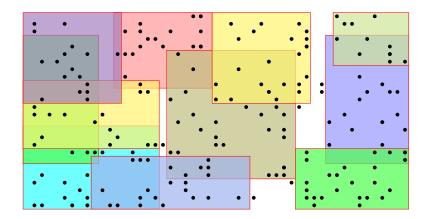




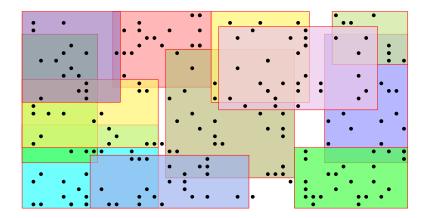




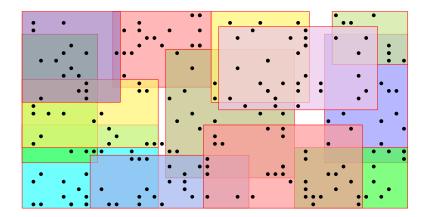






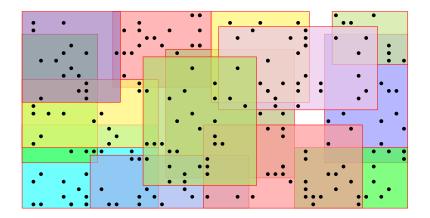






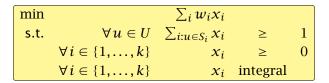


# **Set Cover**





## **IP-Formulation of Set Cover**





### **Vertex Cover**

Given a graph G = (V, E) and a weight  $w_v$  for every node. Find a vertex subset  $S \subseteq V$  of minimum weight such that every edge is incident to at least one vertex in S.



# **IP-Formulation of Vertex Cover**

$$\begin{array}{c|cccc} \min & & \sum_{v \in V} w_v x_v \\ \text{s.t.} & \forall e = (i, j) \in E & & x_i + x_j & \geq & 1 \\ & \forall v \in V & & x_v & \in & \{0, 1\} \end{array}$$



# **Maximum Weighted Matching**

Given a graph G = (V, E), and a weight  $w_e$  for every edge  $e \in E$ . Find a subset of edges of maximum weight such that no vertex is incident to more than one edge.





# **Maximum Weighted Matching**

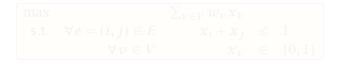
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max	$\sum_{e\in E} w_e x_e$			
s.t.	$\forall v \in V$	$\sum_{e:v \in e} x_e$	$\leq$	1
	$\forall e \in E$	$x_e$	$\in$	$\{0, 1\}$



# **Maximum Independent Set**

Given a graph G = (V, E), and a weight  $w_v$  for every node  $v \in V$ . Find a subset  $S \subseteq V$  of nodes of maximum weight such that no two vertices in S are adjacent.





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# Knapsack

Given a set of items  $\{1, ..., n\}$ , where the *i*-th item has weight  $w_i$  and profit  $p_i$ , and given a threshold K. Find a subset  $I \subseteq \{1, ..., n\}$  of items of total weight at most K such that the profit is maximized.





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$$\begin{array}{|c|c|c|c|c|} \max & & \sum_{i=1}^{n} p_i x_i \\ \text{s.t.} & & \sum_{i=1}^{n} w_i x_i &\leq K \\ & \forall i \in \{1, \dots, n\} & & x_i &\in \{0, 1\} \end{array}$$



### Relaxations

#### **Definition 67**

A linear program LP is a relaxation of an integer program IP if any feasible solution for IP is also feasible for LP and if the objective values of these solutions are identical in both programs.

We obtain a relaxation for all examples by writing  $x_i \in [0, 1]$ instead of  $x_i \in \{0, 1\}$ .



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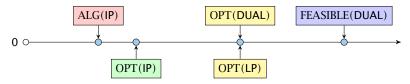


By solving a relaxation we obtain an upper bound for a maximization problem and a lower bound for a minimization problem.

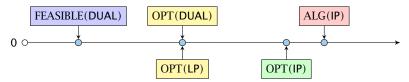


### Relations

#### **Maximization Problems:**



#### **Minimization Problems:**





We first solve the LP-relaxation and then we round the fractional values so that we obtain an integral solution.

Set Cover relaxation:

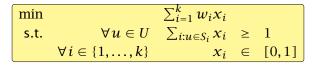


Let  $f_u$  be the number of sets that the element u is contained in (the frequency of u). Let  $f = \max_u \{f_u\}$  be the maximum frequency.



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#### Rounding Algorithm:

Set all  $x_i$ -values with  $x_i \ge \frac{1}{f}$  to 1. Set all other  $x_i$ -values to 0.



### Lemma 68

The rounding algorithm gives an f-approximation.

### **Proof:** Every $u \in U$ is covered.

- We know that Spread 20 201
- The sum contains at most (i) < () elements...
- Therefore one of the sets that contain or must have  $x_i \approx 0/2$
- This set will be selected. Hence, at is covered.



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- We know that  $\sum_{i:u \in S_i} x_i \ge 1$ .
- The sum contains at most  $f_u \leq f$  elements.
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$$\sum_{i\in I} w_i$$



$$\sum_{i \in I} w_i \le \sum_{i=1}^k w_i (f \cdot x_i)$$



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### **Relaxation for Set Cover**

#### Primal:

 $\begin{array}{c|c} \min & \sum_{i \in I} w_i x_i \\ \text{s.t. } \forall u & \sum_{i: u \in S_i} x_i \ge 1 \\ & x_i \ge 0 \end{array}$ 

Dual:





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#### Dual:

$$\begin{array}{c|c}
\max & \sum_{u \in U} \mathcal{Y}_{u} \\
\text{s.t. } \forall i & \sum_{u:u \in S_{i}} \mathcal{Y}_{u} \leq w_{i} \\
\mathcal{Y}_{u} \geq 0
\end{array}$$



#### **Rounding Algorithm:**

Let I denote the index set of sets for which the dual constraint is tight. This means for all  $i \in I$ 

$$\sum_{u:u\in S_i} y_u = w_i$$



**Lemma 69** *The resulting index set is an f-approximation.* 

**Proof:** Every  $u \in U$  is covered.

- Suppose there is a u that is not covered.
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 $I\subseteq I'$  .

- Suppose that we take S<sub>i</sub> in the first algorithm. Let,
   This means on a single sector.
- Because of Complementary Slackness Conditions the corresponding constraint in the dual must be tight.
- Hence, the second algorithm will also choose Similar



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- Suppose that we take  $S_i$  in the first algorithm. I.e.,  $i \in I$ .
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The previous two rounding algorithms have the disadvantage that it is necessary to solve the LP. The following method also gives an f-approximation without solving the LP.

For estimating the cost of the solution we only required two properties.

The solution is dual feasible and, hence,

where x is an optimum solution to the primal LP.

The set 4 contains only sets for which the dual inequality is tight.



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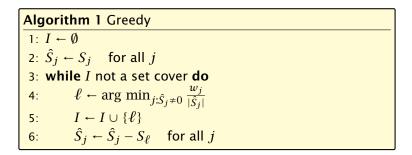
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2. The set *I* contains only sets for which the dual inequality is tight.

Algorithm 1 PrimalDual
$1: \ \mathcal{Y} \leftarrow 0$ $2: \ I \leftarrow \emptyset$
2: $I \leftarrow \emptyset$
3: while exists $u \notin \bigcup_{i \in I} S_i$ do
4: increase dual variable $y_u$ until constraint for some
new set $S_\ell$ becomes tight
5: $I \leftarrow I \cup \{\ell\}$





In every round the Greedy algorithm takes the set that covers remaining elements in the most cost-effective way.

We choose a set such that the ratio between cost and still uncovered elements in the set is minimized.



#### Lemma 70

Given positive numbers  $a_1, \ldots, a_k$  and  $b_1, \ldots, b_k$ , and  $S \subseteq \{1, \ldots, k\}$  then

$$\min_{i} \frac{a_i}{b_i} \le \frac{\sum_{i \in S} a_i}{\sum_{i \in S} b_i} \le \max_{i} \frac{a_i}{b_i}$$



Let  $n_{\ell}$  denote the number of elements that remain at the beginning of iteration  $\ell$ .  $n_1 = n = |U|$  and  $n_{s+1} = 0$  if we need s iterations.

In the  $\ell$ -th iteration

since an optimal algorithm can cover the remaining  $n_\ell$  elements with cost <code>OPT</code>.

Let  $\hat{S}_j$  be a subset that minimizes this ratio. Hence,  $w_j/|\hat{S}_j| \leq \frac{\text{OPT}}{n_\ell}$ .



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 $\min_{j} \frac{w_{j}}{|\hat{S}_{j}|} \leq \frac{\sum_{j \in \text{OPT}} w_{j}}{\sum_{j \in \text{OPT}} |\hat{S}_{j}|} = \frac{\text{OPT}}{\sum_{j \in \text{OPT}} |\hat{S}_{j}|} \leq \frac{\text{OPT}}{n_{\ell}}$ 

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Let  $\hat{S}_j$  be a subset that minimizes this ratio. Hence,  $w_j/|\hat{S}_j| \leq \frac{\text{OPT}}{n_\ell}$ .



Adding this set to our solution means  $n_{\ell+1} = n_{\ell} - |\hat{S}_j|$ .

$$w_j \le \frac{|\hat{S}_j| \text{OPT}}{n_\ell} = \frac{n_\ell - n_{\ell+1}}{n_\ell} \cdot \text{OPT}$$



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$$\sum_{j \in I} w_j \le \sum_{\ell=1}^s \frac{n_\ell - n_{\ell+1}}{n_\ell} \cdot \text{OPT}$$



$$\sum_{j \in I} w_j \le \sum_{\ell=1}^{s} \frac{n_{\ell} - n_{\ell+1}}{n_{\ell}} \cdot \text{OPT}$$
$$\le \text{OPT} \sum_{\ell=1}^{s} \left( \frac{1}{n_{\ell}} + \frac{1}{n_{\ell} - 1} + \dots + \frac{1}{n_{\ell+1} + 1} \right)$$



$$\sum_{j \in I} w_j \le \sum_{\ell=1}^s \frac{n_\ell - n_{\ell+1}}{n_\ell} \cdot \text{OPT}$$
$$\le \text{OPT} \sum_{\ell=1}^s \left( \frac{1}{n_\ell} + \frac{1}{n_\ell - 1} + \dots + \frac{1}{n_{\ell+1} + 1} \right)$$
$$= \text{OPT} \sum_{i=1}^k \frac{1}{i}$$

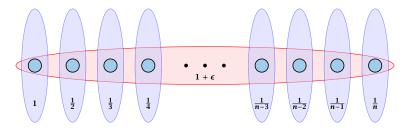


$$\sum_{j \in I} w_j \le \sum_{\ell=1}^s \frac{n_\ell - n_{\ell+1}}{n_\ell} \cdot \text{OPT}$$
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$$= \text{OPT} \sum_{i=1}^k \frac{1}{i}$$
$$= H_n \cdot \text{OPT} \le \text{OPT}(\ln n + 1) \quad .$$



# **Technique 4: The Greedy Algorithm**

#### A tight example:





13.4 Greedy

# **Technique 5: Randomized Rounding**

#### One round of randomized rounding: Pick set $S_j$ uniformly at random with probability $1 - x_j$ (for all j).

**Version A:** Repeat rounds until you nearly have a cover. Cover remaining elements by some simple heuristic.

**Version B:** Repeat for *s* rounds. If you have a cover STOP. Otherwise, repeat the whole algorithm.



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Pr[*u* not covered in one round]



 $\Pr[u \text{ not covered in one round}]$ 

$$= \prod_{j:u\in S_j} (1-x_j)$$



 $\Pr[u \text{ not covered in one round}]$ 

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$$= \prod_{j:u\in S_j} (1-x_j) \le \prod_{j:u\in S_j} e^{-x_j}$$
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Pr[*u* not covered in one round]

$$= \prod_{j:u\in S_j} (1-x_j) \le \prod_{j:u\in S_j} e^{-x_j}$$
$$= e^{-\sum_{j:u\in S_j} x_j} \le e^{-1} .$$

Probability that  $u \in U$  is not covered (after  $\ell$  rounds):

$$\Pr[u \text{ not covered after } \ell \text{ round}] \leq \frac{1}{e^{\ell}}$$
.







=  $\Pr[u_1 \text{ not covered} \lor u_2 \text{ not covered} \lor \ldots \lor u_n \text{ not covered}]$ 



=  $\Pr[u_1 \text{ not covered} \lor u_2 \text{ not covered} \lor \ldots \lor u_n \text{ not covered}]$ 

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**Lemma 71** With high probability  $O(\log n)$  rounds suffice.



$$= \Pr[u_1 \text{ not covered } \lor u_2 \text{ not covered } \lor \dots \lor u_n \text{ not covered}]$$
  
$$\leq \sum_i \Pr[u_i \text{ not covered after } \ell \text{ rounds}] \leq ne^{-\ell} .$$

#### **Lemma 71** With high probability $O(\log n)$ rounds suffice.

#### With high probability:

For any constant  $\alpha$  the number of rounds is at most  $O(\log n)$  with probability at least  $1 - n^{-\alpha}$ .



Proof: We have

 $\Pr[\#\mathsf{rounds} \ge (\alpha + 1) \ln n] \le n e^{-(\alpha + 1) \ln n} = n^{-\alpha} .$ 



Version A.

Repeat for  $s = (\alpha + 1) \ln n$  rounds. If you don't have a cover simply take for each element u the cheapest set that contains u.



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 $E[\cos t] \le (\alpha + 1) \ln n \cdot \cos(LP) + (n \cdot OPT) n^{-\alpha}$ 



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 $E[\text{cost}] \le (\alpha+1) \ln n \cdot \text{cost}(LP) + (n \cdot \text{OPT})n^{-\alpha} = \mathcal{O}(\ln n) \cdot \text{OPT}$ 



Version B.

Repeat for  $s = (\alpha + 1) \ln n$  rounds. If you don't have a cover simply repeat the whole process.

E[cost] =



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Repeat for  $s = (\alpha + 1) \ln n$  rounds. If you don't have a cover simply repeat the whole process.

```
E[cost] = Pr[success] \cdot E[cost | success] 
+ Pr[no success] \cdot E[cost | no success]
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This means

*E*[cost | success]



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E[cost] = Pr[success] \cdot E[cost | success] + Pr[no success] \cdot E[cost | no success]
```

This means

```
E[\cos t | \text{success}] = \frac{1}{\Pr[\text{succ.}]} \Big( E[\cos t] - \Pr[\text{no success}] \cdot E[\cos t | \text{no success}] \Big)
```



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*E*[cost | success]

$$= \frac{1}{\Pr[\mathsf{succ.}]} \Big( E[\cos t] - \Pr[\mathsf{no \ success}] \cdot E[\cos t | \mathsf{no \ success}] \Big)$$
  
$$\leq \frac{1}{\Pr[\mathsf{succ.}]} E[\cos t] \leq \frac{1}{1 - n^{-\alpha}} (\alpha + 1) \ln n \cdot \operatorname{cost}(\operatorname{LP})$$



Version B.

Repeat for  $s = (\alpha + 1) \ln n$  rounds. If you don't have a cover simply repeat the whole process.

 $E[cost] = Pr[success] \cdot E[cost | success] + Pr[no success] \cdot E[cost | no success]$ 

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$$= \frac{1}{\Pr[\mathsf{succ.}]} \left( E[\cos t] - \Pr[\mathsf{no \ success}] \cdot E[\cos t \mid \mathsf{no \ success}] \right)$$
  
$$\leq \frac{1}{\Pr[\mathsf{succ.}]} E[\cos t] \leq \frac{1}{1 - n^{-\alpha}} (\alpha + 1) \ln n \cdot \operatorname{cost}(\operatorname{LP})$$
  
$$\leq 2(\alpha + 1) \ln n \cdot \operatorname{OPT}$$



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for  $n \ge 2$  and  $\alpha \ge 1$ .



# Randomized rounding gives an $O(\log n)$ approximation. The running time is polynomial with high probability.

#### Theorem 72 (without proof)

There is no approximation algorithm for set cover with approximation guarantee better than  $\frac{1}{2}\log n$  unless NP has quasi-polynomial time algorithms (algorithms with running time  $2poly(\log n)$ ).



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# **Integrality Gap**

The integrality gap of the SetCover LP is  $\Omega(\log n)$ .

▶  $n = 2^k - 1$ 

- Elements are all vectors  $\vec{x}$  over GF[2] of length k (excluding zero vector).
- Every vector  $\vec{y}$  defines a set as follows

$$S_{\vec{y}} := \{ \vec{x} \mid \vec{x}^T \vec{y} = 1 \}$$

• each set contains  $2^{k-1}$  vectors; each vector is contained in  $2^{k-1}$  sets

• 
$$x_i = \frac{1}{2^{k-1}} = \frac{2}{n+1}$$
 is fractional solution.



# **Integrality Gap**

#### Every collection of p < k sets does not cover all elements.

Hence, we get a gap of  $\Omega(\log n)$ .



#### **Techniques:**

- Deterministic Rounding
- Rounding of the Dual
- Primal Dual
- Greedy
- Randomized Rounding
- Local Search
- Rounding Data + Dynamic Programming



# Scheduling Jobs on Identical Parallel Machines

Given n jobs, where job  $j \in \{1, ..., n\}$  has processing time  $p_j$ . Schedule the jobs on m identical parallel machines such that the Makespan (finishing time of the last job) is minimized.

Here the variable  $x_{j,i}$  is the decision variable that describes whether job j is assigned to machine i.



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min		L		
s.t.	$\forall$ machines $i$	$\sum_j p_j \cdot x_{j,i}$	$\leq$	L
	$\forall jobs \ j$	$\sum_{i} x_{j,i} \ge 1$		
	$\forall i, j$	$x_{j,i}$	$\in$	$\{0, 1\}$

Here the variable  $x_{j,i}$  is the decision variable that describes whether job j is assigned to machine i.



Let for a given schedule  $C_j$  denote the finishing time of machine j, and let  $C_{\text{max}}$  be the makespan.

Let  $C^*_{max}$  denote the makespan of an optimal solution.

Clearly

 $C^*_{\max} \ge \max_j p_j$ 

as the longest job needs to be scheduled somewhere.



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## The average work performed by a machine is $\frac{1}{m} \sum_j p_j$ . Therefore,





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$$C_{\max}^* \ge \frac{1}{m} \sum_j p_j$$



A local search algorithm successively makes certain small (cost/profit improving) changes to a solution until it does not find such changes anymore.

It is conceptionally very different from a Greedy algorithm as a feasible solution is always maintained.



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## Local Search for Scheduling

**Local Search Strategy:** Take the job that finishes last and try to move it to another machine. If there is such a move that reduces the makespan, perform the switch.

REPEAT



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REPEAT



Let  $S_{\ell}$  be its start time, and let  $C_{\ell}$  be its completion time.



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The interval  $[S_{\ell}, C_{\ell}]$  is of length  $p_{\ell} \leq C^*_{\max}$ .

During the first interval  $[0, S_{\ell}]$  all processors are busy, and, hence, the total work performed in this interval is

$$m \cdot S_{\ell} \leq \sum_{j \neq \ell} p_j$$
.



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$$m \cdot S_{\ell} \leq \sum_{j \neq \ell} p_j$$
.

$$p_{\ell} + \frac{1}{m} \sum_{j \neq \ell} p_j = (1 - \frac{1}{m}) p_{\ell} + \frac{1}{m} \sum_j p_j \le (2 - \frac{1}{m}) C_{\max}^*$$



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Hence, the length of the schedule is at most

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14.1 Local Search

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.

Hence, the length of the schedule is at most

$$p_{\ell} + \frac{1}{m} \sum_{j \neq \ell} p_j = (1 - \frac{1}{m}) p_{\ell} + \frac{1}{m} \sum_j p_j \le (2 - \frac{1}{m}) C_{\max}^*$$



14.1 Local Search

## A Tight Example

$$p_{\ell} \approx S_{\ell} + \frac{S_{\ell}}{m-1}$$

$$\frac{\text{ALG}}{\text{OPT}} = \frac{S_{\ell} + p_{\ell}}{p_{\ell}} \approx \frac{2 + \frac{1}{m-1}}{1 + \frac{1}{m-1}} = 2 - \frac{1}{m}$$

$$p_{\ell}$$

$$S_{\ell}$$

**List Scheduling:** 

Order all processes in a list. When a machine runs empty assign the next yet unprocessed job to it.

Alternatively:

Consider processes in some order. Assign the *i*-th process to the least loaded machine.



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#### Lemma 73

If we order the list according to non-increasing processing times the approximation guarantee of the list scheduling strategy improves to 4/3.



- Let  $p_1 \ge \cdots \ge p_n$  denote the processing times of a set of jobs that form a counter-example.
- Wlog. the last job to finish is n (otw. deleting this job gives another counter-example with fewer jobs).
- If p<sub>n</sub> ≤ C<sup>\*</sup><sub>max</sub>/3 the previous analysis gives us a schedule length of at most

$$C_{\max}^* + p_n \le \frac{4}{3}C_{\max}^* \ .$$

- Hence,  $p_n > C_{\max}^n/3$ .
- This means that all jobs must have a processing time of the second second
- But then any machine in the optimum schedule can handle at most two jobs.
- For such instances Longest-Processing-Time-First is optimal.



- Let  $p_1 \ge \cdots \ge p_n$  denote the processing times of a set of jobs that form a counter-example.
- Wlog. the last job to finish is n (otw. deleting this job gives another counter-example with fewer jobs).
- If  $p_n \le C^*_{\max}/3$  the previous analysis gives us a schedule length of at most

$$C_{\max}^* + p_n \le \frac{4}{3}C_{\max}^* \ .$$

- Hence,  $p_n \ge C_{\max}^*/3$ .
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14.2 Greedy

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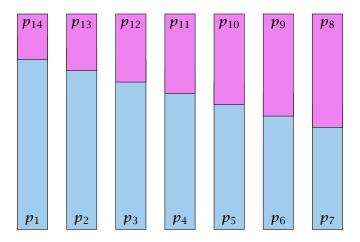
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When in an optimal solution a machine can have at most 2 jobs the optimal solution looks as follows.





- We can assume that one machine schedules p<sub>1</sub> and p<sub>n</sub> (the largest and smallest job).
- If not assume wlog, that p<sub>1</sub> is scheduled on machine A and p<sub>n</sub> on machine B.
- ► Let *p*<sub>A</sub> and *p*<sub>B</sub> be the other job scheduled on A and B, respectively.
- ▶  $p_1 + p_n \le p_1 + p_A$  and  $p_A + p_B \le p_1 + p_A$ , hence scheduling  $p_1$  and  $p_n$  on one machine and  $p_A$  and  $p_B$  on the other, cannot increase the Makespan.
- Repeat the above argument for the remaining machines.



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- ▶ 2 jobs with length 2m, 2m 1, 2m 2, ..., m + 1 (2m 2 jobs in total)



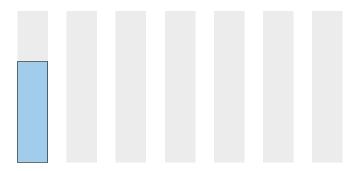


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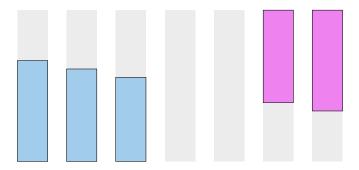


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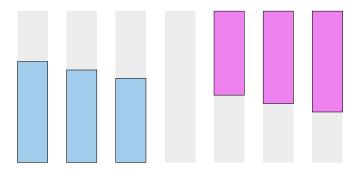


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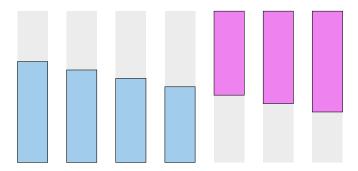


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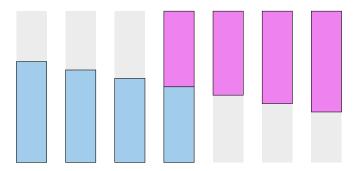


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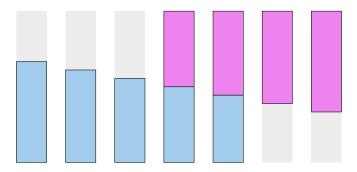


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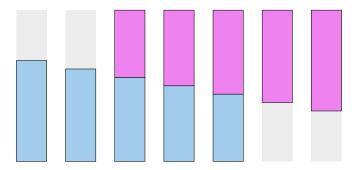


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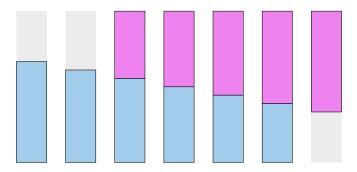


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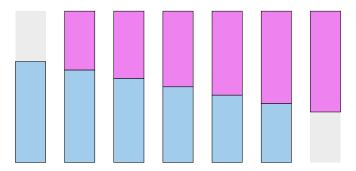


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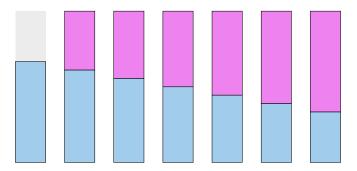


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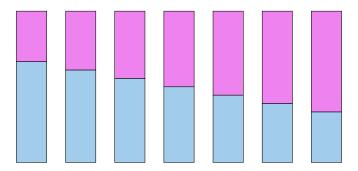


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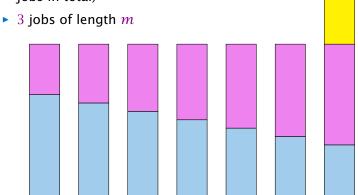


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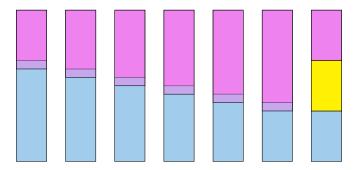


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Given a set of cities  $(\{1, ..., n\})$  and a symmetric matrix  $C = (c_{ij}), c_{ij} \ge 0$  that specifies for every pair  $(i, j) \in [n] \times [n]$  the cost for travelling from city *i* to city *j*. Find a permutation  $\pi$  of the cities such that the round-trip cost

$$C_{\pi(1)\pi(n)} + \sum_{i=1}^{n-1} C_{\pi(i)\pi(i+1)}$$

is minimized.



## Theorem 74

There does not exist an  $O(2^n)$ -approximation algorithm for TSP.

## Hamiltonian Cycle:

- Given an instance to HAMPATH we create an instance for TSP.
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# Metric Traveling Salesman

# In the metric version we assume for every triple $i, j, k \in \{1, \dots, n\}$

 $c_{ij} \leq c_{ij} + c_{jk}$  .

It is convenient to view the input as a complete undirected graph G = (V, E), where  $c_{ij}$  for an edge (i, j) defines the distance between nodes i and j.



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#### Lemma 75

The cost  $OPT_{TSP}(G)$  of an optimum traveling salesman tour is at least as large as the weight  $OPT_{MST}(G)$  of a minimum spanning tree in G.

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#### Start with a tour on a subset *S* containing a single node.

- Take the node v closest to S. Add it S and expand the existing tour on S to include v.
- Repeat until all nodes have been processed.

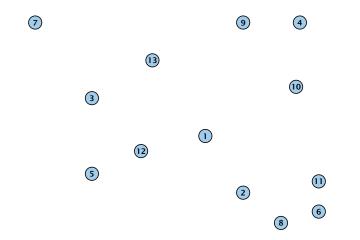


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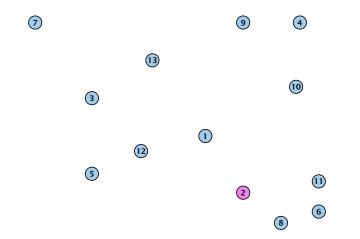


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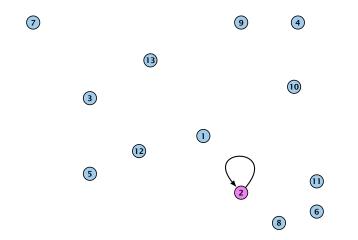




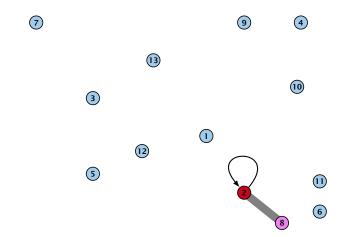




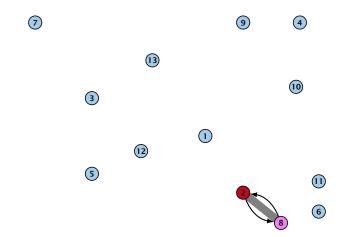




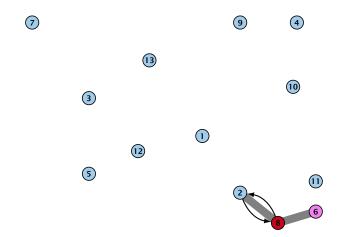




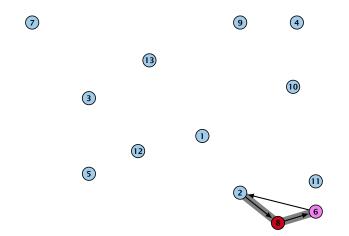




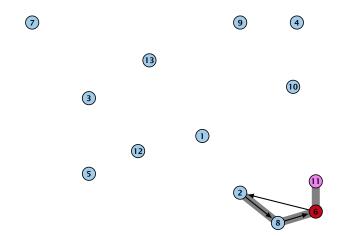




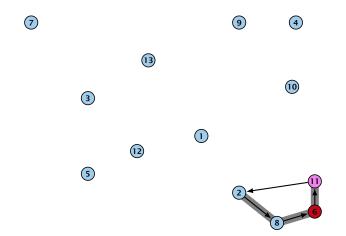




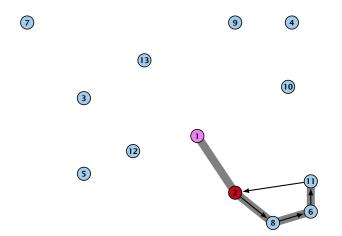




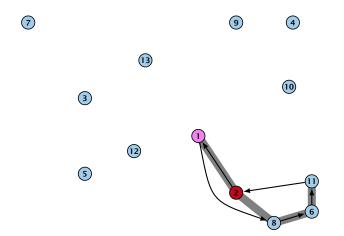




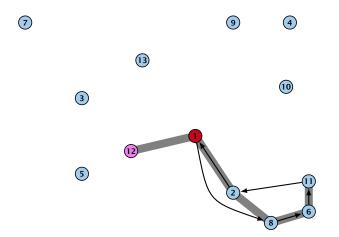




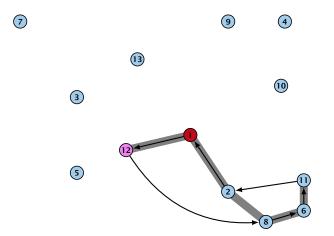




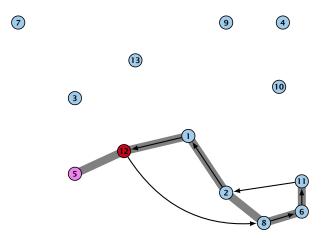




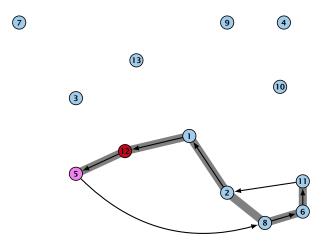




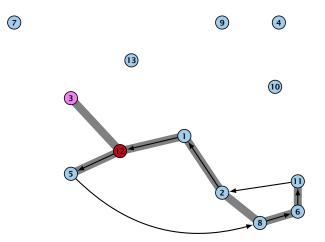




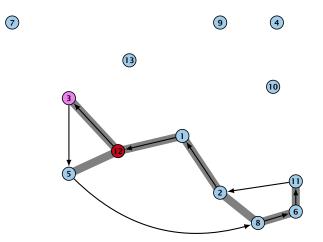




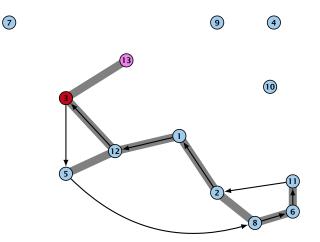




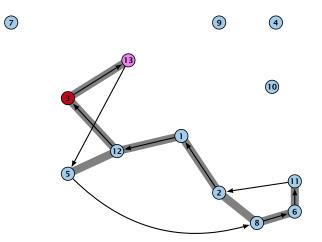




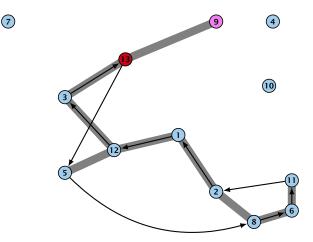






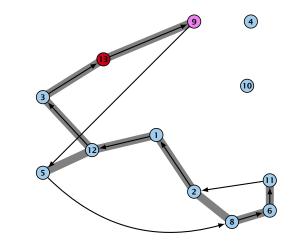






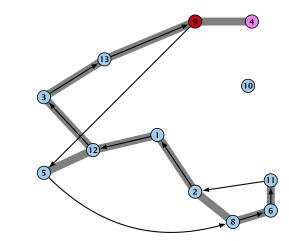


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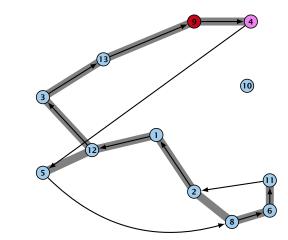


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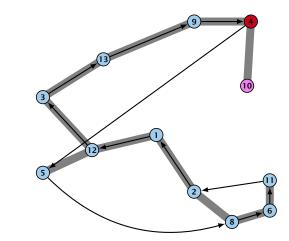


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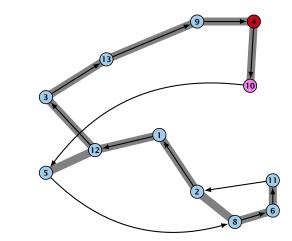


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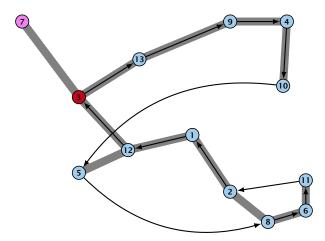




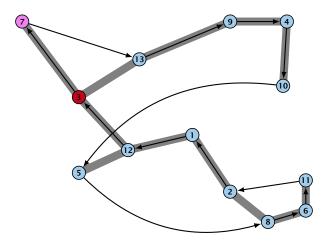
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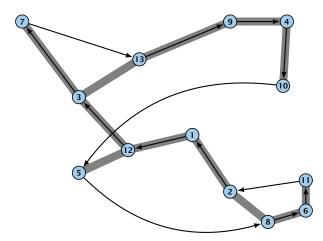




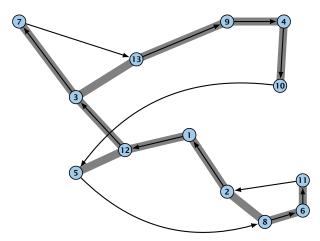












Lemma 76

The Greedy algorithm is a 2-approximation algorithm.

Let  $S_i$  be the set at the start of the *i*-th iteration, and let  $v_i$  denote the node added during the iteration.

Further let  $s_i \in S_i$  be the node closest to  $v_i \in S_i$ .

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Suppose that we are given an Eulerian graph G' = (V, E', c') of G = (V, E, c) such that for any edge  $(i, j) \in E' c'(i, j) \ge c(i, j)$ .

Then we can find a TSP-tour of cost at most

$$\sum_{e\in E'}c'(e)$$

- Find an Euler tour of Gal.
- Fix a permutation of the cities (i.e., a TSP-tour) by traversing the Euler tour and only note the first occurrence of a city.
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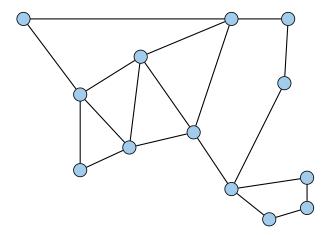


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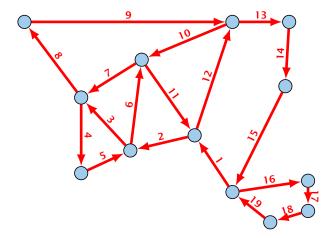
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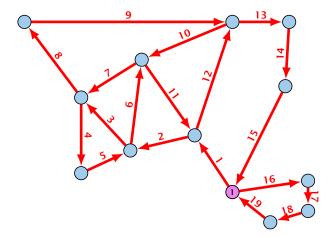
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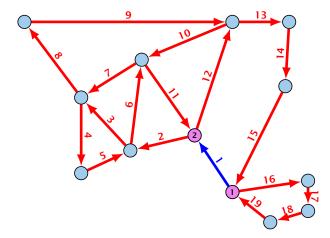




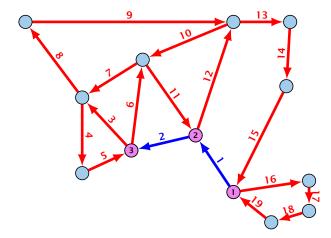




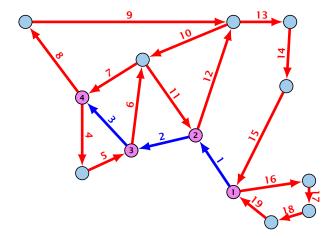




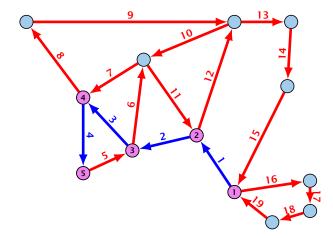




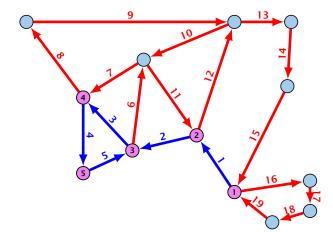




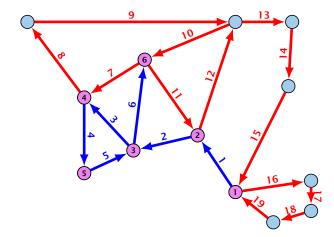




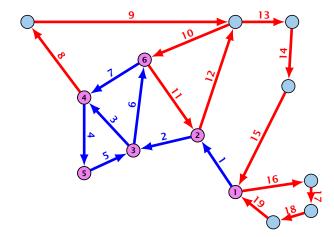




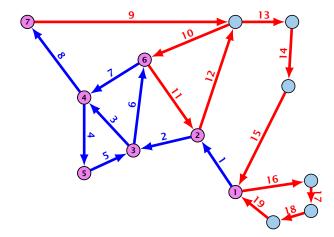




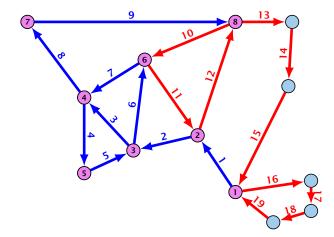




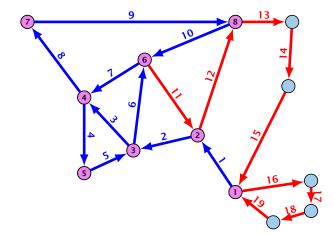




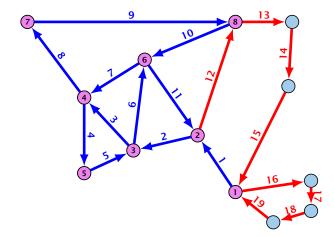




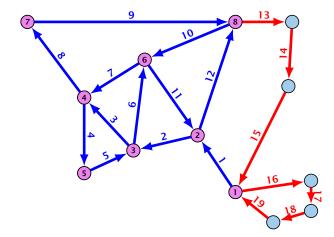




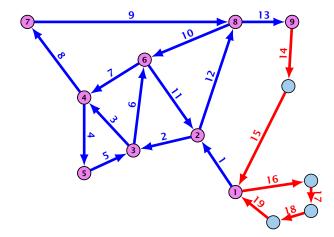




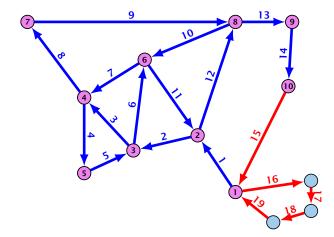




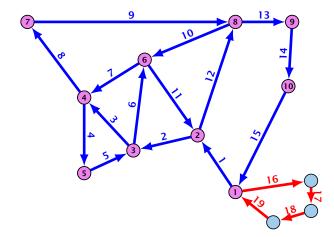




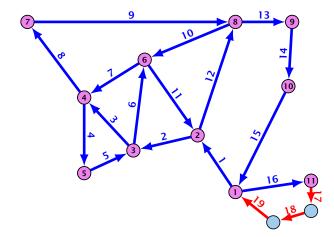




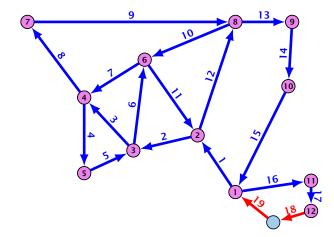




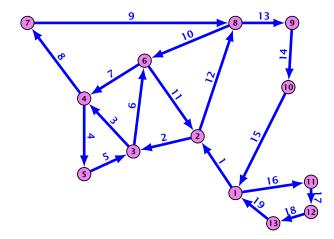




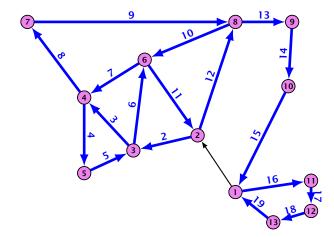




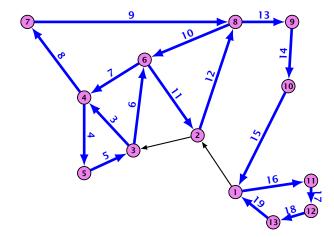




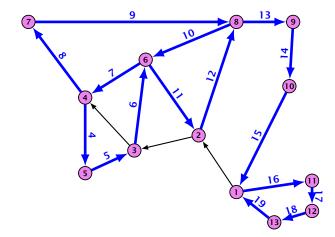




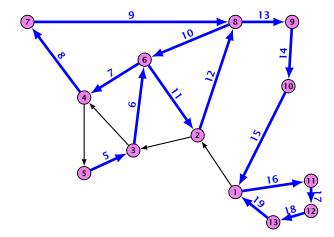




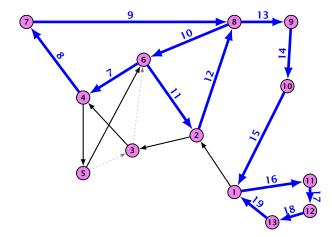




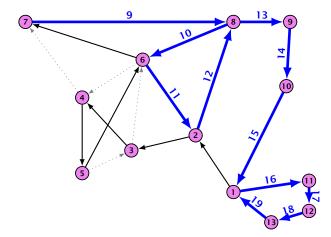




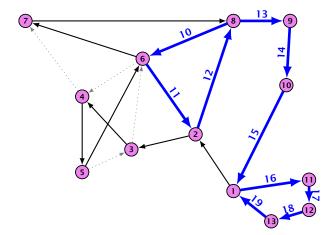




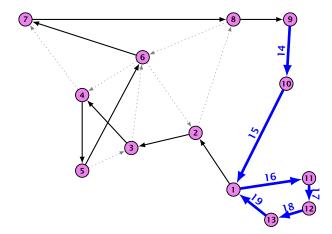




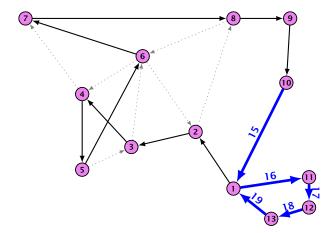




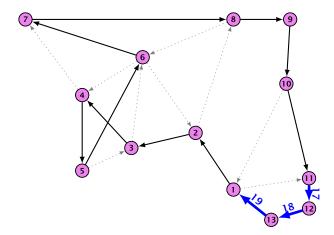




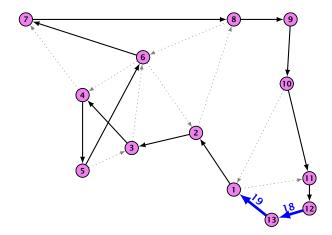




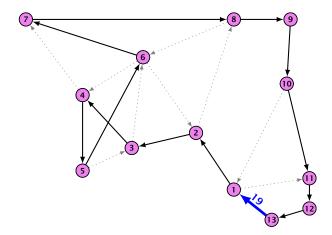




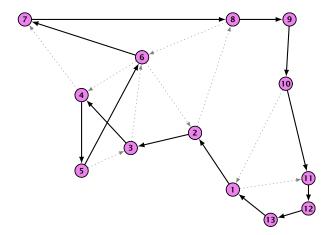




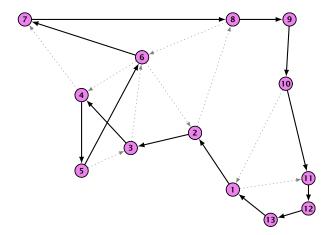














Consider the following graph:

- Compute an MST of *G*.
- Duplicate all edges.

This graph is Eulerian, and the total cost of all edges is at most  $2 \cdot OPT_{MST}(G)$ .

Hence, short-cutting gives a tour of cost no more than  $2 \cdot OPT_{MST}(G)$  which means we have a 2-approximation.



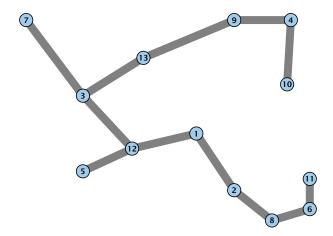
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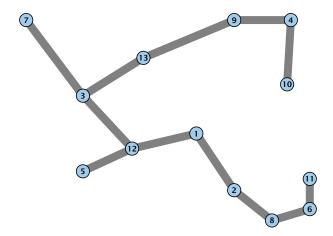
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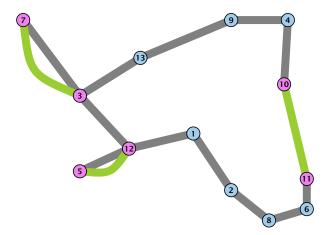




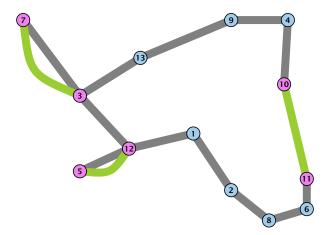




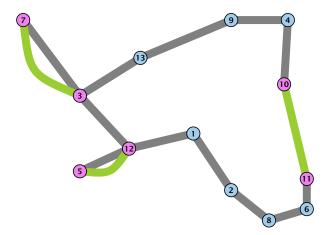




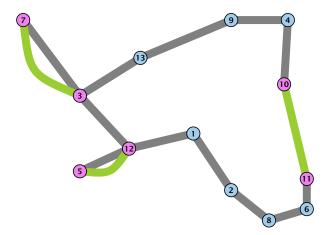














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An optimal tour on the odd-degree vertices has cost at most  $OPT_{TSP}(G)$ .

However, the edges of this tour give rise to two disjoint matchings. One of these matchings must have weight less than  $OPT_{TSP}(G)/2$ .

Adding this matching to the MST gives an Eulerian graph with edge weight at most

 $OPT_{MST}(G) + OPT_{TSP}(G)/2 \le \frac{3}{2}OPT_{TSP}(G)$ ,

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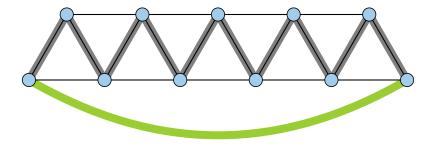
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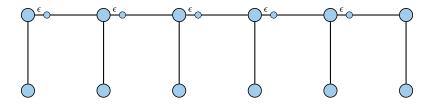
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# **Christofides. Tight Example**



- optimal tour: n edges.
- ▶ MST: *n* − 1 edges.
- weight of matching (n + 1)/2 1
- MST+matching  $\approx 3/2 \cdot n$

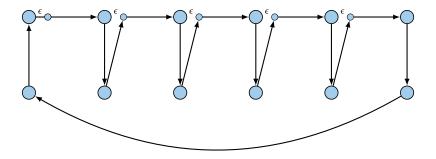
# Tree shortcutting. Tight Example



#### edges have Euclidean distance.



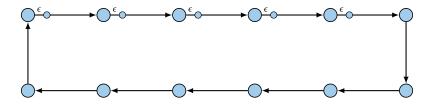
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# 16 Rounding Data + Dynamic Programming

#### Knapsack:

Given a set of items  $\{1, ..., n\}$ , where the *i*-th item has weight  $w_i \in \mathbb{N}$  and profit  $p_i \in \mathbb{N}$ , and given a threshold W. Find a subset  $I \subseteq \{1, ..., n\}$  of items of total weight at most W such that the profit is maximized (we can assume each  $w_i \leq W$ ).





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max		$\sum_{i=1}^{n} p_i x_i$		
s.t.		$\sum_{i=1}^{n} w_i x_i$	$\leq$	W
	$\forall i \in \{1, \dots, n\}$	$x_i$	∈	$\{0, 1\}$



16.1 Knapsack

Algorithm 1 Knapsack1:  $A(1) \leftarrow [(0,0), (p_1, w_1)]$ 2: for  $j \leftarrow 2$  to n do3:  $A(j) \leftarrow A(j-1)$ 4: for each  $(p, w) \in A(j-1)$  do5: if  $w + w_j \le W$  then6: add  $(p + p_j, w + w_j)$  to A(j)7: remove dominated pairs from A(j)8: return  $\max_{(p,w)\in A(n)} p$ 

The running time is  $O(n \cdot \min\{W, P\})$ , where  $P = \sum_i p_i$  is the total profit of all items. This is only pseudo-polynomial.



#### **Definition 77**

An algorithm is said to have pseudo-polynomial running time if the running time is polynomial when the numerical part of the input is encoded in unary.



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Running time is at most

 $\mathcal{O}(nP')$ 



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- Set  $\mu := \epsilon M/n$ .
- Set  $p'_i := \lfloor p_i / \mu \rfloor$  for all *i*.
- Run the dynamic programming algorithm on this revised instance.

Running time is at most

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Together with the obervation that if each  $p_i \ge \frac{1}{3}C_{\max}^*$  then LPT is optimal this gave a 4/3-approximation.



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Idea:

1. Find the optimum Makespan for the long jobs by brute force.



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#### Idea:

- 1. Find the optimum Makespan for the long jobs by brute force.
- 2. Then use the list scheduling algorithm for the short jobs, always assigning the next job to the least loaded machine.



We still have a cost of

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If  $\ell$  is a short job its length is at most

 $p_\ell \leq \sum_j p_j / (mk)$ 

which is at most  $C^*_{\max}/k$ .



#### Hence we get a schedule of length at most

 $\left(1+\frac{1}{k}\right)C_{\max}^*$ 

There are at most km long jobs. Hence, the number of possibilities of scheduling these jobs on m machines is at most  $m^{km}$ , which is constant if m is constant. Hence, it is easy to implement the algorithm in polynomial time.

#### Theorem 78

The above algorithm gives a polynomial time approximation scheme (PTAS) for the problem of scheduling n jobs on m identical machines if m is constant.

We choose  $k = \lceil \frac{1}{\epsilon} \rceil$ .



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We first design an algorithm that works as follows: On input of *T* it either finds a schedule of length  $(1 + \frac{1}{k})T$  or certifies that no schedule of length at most *T* exists (assume  $T \ge \frac{1}{m}\sum_j p_j$ ).

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#### • We round all long jobs down to multiples of $T/k^2$ .

- For these rounded sizes we first find an optimal schedule.
- If this schedule does not have length at most T we conclude that also the original sizes don't allow such a schedule.
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# After the first phase the rounded sizes of the long jobs assigned to a machine add up to at most T.

There can be at most k (long) jobs assigned to a machine as otw. their rounded sizes would add up to more than T (note that the rounded size of a long job is at least T/k).

Since, jobs had been rounded to multiples of  $T/k^2$  going from rounded sizes to original sizes gives that the Makespan is at most

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$$T + \frac{T}{k} \le \left(1 + \frac{1}{k}\right)T \; .$$



Hence, any large job has rounded size of  $\frac{i}{k^2}T$  for  $i \in \{k, ..., k^2\}$ . Therefore the number of different inputs is at most  $n^{k^2}$ (described by a vector of length  $k^2$  where, the *i*-th entry describes the number of jobs of size  $\frac{i}{k^2}T$ ). This is polynomial.

The schedule/configuration of a particular machine x can be described by a vector of length  $k^2$  where the *i*-th entry describes the number of jobs of rounded size  $\frac{i}{k^2}T$  assigned to x. There are only  $(k + 1)^{k^2}$  different vectors.



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If  $OPT(n_1, \ldots, n_{k^2}) \leq m$  we can schedule the input.

We have

 $OPT(n_1, \dots, n_{k^2}) = \begin{cases} 0 & (n_1, \dots, n_{k^2}) = 0\\ 1 + \min_{(s_1, \dots, s_{k^2}) \in C} OPT(n_1 - s_1, \dots, n_{k^2} - s_{k^2}) & (n_1, \dots, n_{k^2}) \ge 0\\ \infty & \text{otw.} \end{cases}$ 

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Scheduling on identical machines with the goal of minimizing Makespan is a strongly NP-complete problem.

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- Suppose we have an instance with polynomially bounded processing times p<sub>i</sub> ≤ q(n)
- We set  $k := \lceil 2nq(n) \rceil \ge 2 \text{ OPT}$

$$ALG \le \left(1 + \frac{1}{k}\right) OPT \le OPT + \frac{1}{2}$$

- But this means that the algorithm computes the optimal solution as the optimum is integral.
- This means we can solve problem instances if processing times are polynomially bounded
- Running time is  $\mathcal{O}(\operatorname{poly}(n,k)) = \mathcal{O}(\operatorname{poly}(n))$
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# **More General**

Let  $OPT(n_1, ..., n_A)$  be the number of machines that are required to schedule input vector  $(n_1, ..., n_A)$  with Makespan at most T (*A*: number of different sizes).

If  $OPT(n_1, \ldots, n_A) \le m$  we can schedule the input.

$$OPT(n_1, \dots, n_A) = 0$$

$$= \begin{cases} 0 & (n_1, \dots, n_A) = 0\\ 1 + \min_{(s_1, \dots, s_A) \in C} OPT(n_1 - s_1, \dots, n_A - s_A) & (n_1, \dots, n_A) \ge 0\\ \infty & \text{otw.} \end{cases}$$

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 $|C| \le (B+1)^A$ , where B is the number of jobs that possibly can fit on the same machine.

The running time is then  $O((B+1)^A n^A)$  because the dynamic programming table has just  $n^A$  entries.

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Given *n* items with sizes  $s_1, \ldots, s_n$  where

 $1 > s_1 \ge \cdots \ge s_n > 0$ .

Pack items into a minimum number of bins where each bin can hold items of total size at most 1.

**Theorem 80** There is no  $\rho$ -approximation for Bin Packing with  $\rho < 3/2$  unless P = NP.



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#### **Theorem 80**

There is no  $\rho$ -approximation for Bin Packing with  $\rho < 3/2$  unless P = NP.



### Proof

▶ In the partition problem we are given positive integers  $b_1, \ldots, b_n$  with  $B = \sum_i b_i$  even. Can we partition the integers into two sets *S* and *T* s.t.

$$\sum_{i\in S} b_i = \sum_{i\in T} b_i \quad ?$$

- We can solve this problem by setting  $s_i := 2b_i/B$  and asking whether we can pack the resulting items into 2 bins or not.
- A ρ-approximation algorithm with ρ < 3/2 cannot output 3 or more bins when 2 are optimal.
- Hence, such an algorithm can solve Partition.



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### **Definition 81**

An asymptotic polynomial-time approximation scheme (APTAS) is a family of algorithms  $\{A_{\epsilon}\}$  along with a constant c such that  $A_{\epsilon}$  returns a solution of value at most  $(1 + \epsilon)$ OPT + c for minimization problems.

- Note that for Set Cover or for Knapsack it makes no sense to differentiate between the notion of a PTAS or an APTAS because of scaling.
- However, we will develop an APTAS for 8in Packing.



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- Note that for Set Cover or for Knapsack it makes no sense to differentiate between the notion of a PTAS or an APTAS because of scaling.
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Again we can differentiate between small and large items.

Lemma 82

- If after Greedy we use more than 6 bins, all bins (apart from the last) must be full to at least 1 = 1;
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Choose  $\gamma = \epsilon/2$ . Then we either use  $\ell$  bins or at most

$$\frac{1}{1 - \epsilon/2} \cdot \text{OPT} + 1 \le (1 + \epsilon) \cdot \text{OPT} + 1$$

bins.

It remains to find an algorithm for the large items.



#### Linear Grouping:

- Order large items according to size.
- Let the first k items belong to group 1; the following k items belong to group 2; etc.
- Delete items in the first group;
- Round items in the remaining groups to the size of the largest item in the group.



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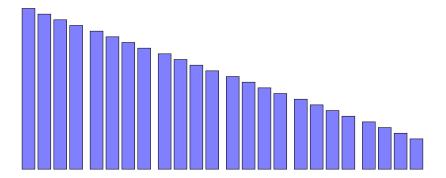
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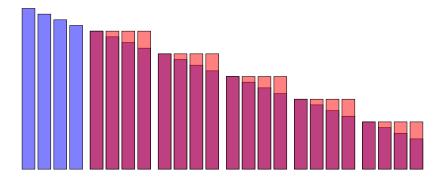
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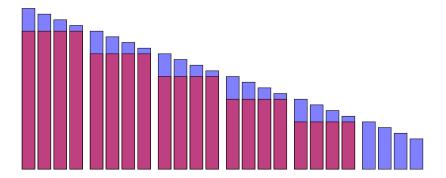




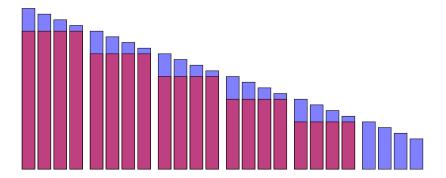














- Any bin packing for / gives a bin packing for / as follows.
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## Proof 2:

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- Pack the items of group 1 into k new bins;
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We set  $k = \lfloor \epsilon \text{SIZE}(I) \rfloor$ .

Then  $n/k \le n/\lfloor \epsilon^2 n/2 \rfloor \le 4/\epsilon^2$  (here we used  $\lfloor \alpha \rfloor \ge \alpha/2$  for  $\alpha \ge 1$ ).

Hence, after grouping we have a constant number of piece sizes  $(4/\epsilon^2)$  and at most a constant number  $(2/\epsilon)$  can fit into any bin.

We can find an optimal packing for such instances by the previous Dynamic Programming approach.

cost (for large items) at most

 $OPT(I') + k \le OPT(I) + \epsilon SIZE(I) \le (1 + \epsilon)OPT(I)$ 

• running time  $\mathcal{O}((\frac{2}{\epsilon}n)^{4/\epsilon^2})$ .

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#### Can we do better?

In the following we show how to obtain a solution where the number of bins is only

 $OPT(I) + \mathcal{O}(\log^2(SIZE(I)))$ .

Note that this is usually better than a guarantee of

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16.4 Advanced Rounding for Bin Packing

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- Let s<sub>1</sub> denote the largest size, and let b<sub>1</sub> denote the number of pieces of size s<sub>1</sub>.
- $s_2$  is second largest size and  $b_2$  number of pieces of size  $s_2$ ;
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## **Change of Notation:**

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# A possible packing of a bin can be described by an *m*-tuple $(t_1, \ldots, t_m)$ , where $t_i$ describes the number of pieces of size $s_i$ . Clearly,



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Let N be the number of configurations (exponential).

Let  $T_1, \ldots, T_N$  be the sequence of all possible configurations (a configuration  $T_j$  has  $T_{ji}$  pieces of size  $s_i$ ).



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$$\begin{array}{c|cccc} \min & & \sum_{j=1}^{N} x_j \\ \text{s.t.} & \forall i \in \{1 \dots m\} & \sum_{j=1}^{N} T_{ji} x_j & \geq & b_i \\ & \forall j \in \{1, \dots, N\} & & x_j & \geq & 0 \\ & \forall j \in \{1, \dots, N\} & & x_j & \text{integral} \end{array}$$



## How to solve this LP?

later...



We can assume that each item has size at least 1/SIZE(I).



## Sort items according to size (monotonically decreasing).

- Process items in this order; close the current group if size of items in the group is at least 2 (or larger). Then open new group.
- I.e.,  $G_1$  is the smallest cardinality set of largest items s.t. total size sums up to at least 2. Similarly, for  $G_2, \ldots, G_{r-1}$ .
- Only the size of items in the last group  $G_r$  may sum up to less than 2.



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- Round all items in a group to the size of the largest group member.
- Delete all items from group  $G_1$  and  $G_r$ .
- For groups  $G_2, \ldots, G_{r-1}$  delete  $n_i n_{i-1}$  items.
- Observe that  $n_i \ge n_{i-1}$ .



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## **Lemma 85** The number of different sizes in I' is at most SIZE(I)/2.

- Each group that survives (recall that Gy and Gy are deleted) has total size at least Gy
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- Consider a group (), that has strictly more items than (),
- It discards  $m_1 m_2$  , pieces of total size at most



- since the smallest piece has size at most 3/mar
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(note that is, self-2000) since we assume that the size of each item is at least ((2000)).

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- ► The total size of items in G<sub>1</sub> and G<sub>r</sub> is at most 6 as a group has total size at most 3.
- Consider a group  $G_i$  that has strictly more items than  $G_{i-1}$ .
- ► It discards n<sub>i</sub> n<sub>i-1</sub> pieces of total size at most

$$3\frac{n_i - n_{i-1}}{n_i} \le \sum_{j=n_{i-1}+1}^{n_i} \frac{3}{j}$$

since the smallest piece has size at most  $3/n_i$ .

Summing over all *i* that have n<sub>i</sub> > n<sub>i-1</sub> gives a bound of at most

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## Algorithm 1 BinPack

- 1: **if** SIZE(I) < 10 **then**
- 2: pack remaining items greedily
- 3: Apply harmonic grouping to create instance I'; pack discarded items in at most  $O(\log(\text{SIZE}(I)))$  bins.
- 4: Let x be optimal solution to configuration LP
- 5: Pack  $\lfloor x_j \rfloor$  bins in configuration  $T_j$  for all j; call the packed instance  $I_1$ .
- 6: Let  $I_2$  be remaining pieces from I'
- 7: Pack  $I_2$  via BinPack $(I_2)$



## $OPT_{LP}(I_1) + OPT_{LP}(I_2) \le OPT_{LP}(I') \le OPT_{LP}(I)$

**Proof:** 

- Each piece surviving in () can be mapped to a piece in () of no lesser size. Hence, (0935) (1015) (0935) (2)
- $|||_{\mathcal{O}}$  is feasible solution for  $\mathcal{O}_{\mathcal{O}}$  (even integral).
- Solution for loss



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Each level of the recursion partitions pieces into three types

- 1. Pieces discarded at this level.
- **2.** Pieces scheduled because they are in  $I_1$ .
- **3.** Pieces in *I*<sup>2</sup> are handed down to the next level.

Pieces of type 2 summed over all recursion levels are packed into at most  $\mathrm{OPT}_{\mathrm{LP}}$  many bins.

Pieces of type 1 are packed into at most

 $\mathcal{O}(\log(\text{SIZE}(I))) \cdot L$ 



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- 1. Pieces discarded at this level.
- **2.** Pieces scheduled because they are in  $I_1$ .
- **3.** Pieces in  $I_2$  are handed down to the next level.

Pieces of type 2 summed over all recursion levels are packed into at most  $\mathrm{OPT}_{\mathrm{LP}}$  many bins.

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many bins where L is the number of recursion levels.



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- ▶ The total size of items in  $I_2$  can be at most  $\sum_{j=1}^{N} x_j \lfloor x_j \rfloor$  which is at most the number of non-zero entries in the solution to the configuration LP.



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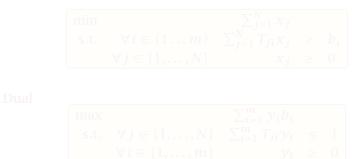


## How to solve the LP?

Let  $T_1, \ldots, T_N$  be the sequence of all possible configurations (a configuration  $T_j$  has  $T_{ji}$  pieces of size  $s_i$ ).

In total we have  $b_i$  pieces of size  $s_i$ .

Primal

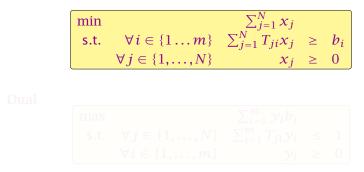




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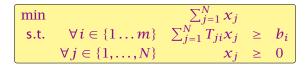




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#### Primal





$$\begin{array}{ll} \max & \sum_{i=1}^{m} y_i b_i \\ \text{s.t.} & \forall j \in \{1, \dots, N\} \quad \sum_{i=1}^{m} T_{ji} y_i \leq 1 \\ & \forall i \in \{1, \dots, m\} \quad y_i \geq 0 \end{array}$$



Suppose that I am given variable assignment y for the dual.

How do I find a violated constraint?

I have to find a configuration  $T_j = (T_{j1}, ..., T_{jm})$  that is feasible, i.e.,

and has a large profit

But this is the Knapsack problem.



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We have FPTAS for Knapsack. This means if a constraint is violated with  $1 + \epsilon' = 1 + \frac{\epsilon}{1-\epsilon}$  we find it, since we can obtain at least  $(1 - \epsilon)$  of the optimal profit.

The solution we get is feasible for:

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**Dual** 

 $\begin{array}{|c|c|c|c|c|} \max & & \sum_{i=1}^{m} y_i b_i \\ \text{s.t.} & \forall j \in \{1, \dots, N\} & \sum_{i=1}^{m} T_{ji} y_i & \leq & 1 + \epsilon' \\ & \forall i \in \{1, \dots, m\} & y_i & \geq & 0 \end{array}$ 

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min		$(1+\epsilon')\sum_{j=1}^N x_j$		
s.t.	$\forall i \in \{1 \dots m\}$	$\sum_{j=1}^{N} T_{ji} x_j$	$\geq$	$b_i$
	$\forall j \in \{1, \dots, N\}$	$x_j$	$\geq$	0

# If the value of the computed dual solution (which may be infeasible) is $\boldsymbol{z}$ then

### $OPT \le z \le (1 + \epsilon')OPT$

- The constraints used when computing a certify that the solution is feasible for 000001.
- Suppose that we drop all unused constraints in 00040. We will compute the same solution feasible for 000001.
- Let DUAL be DUAL without unused constraints.
- The dual to 000000 is 000000 where we ignore variables for which the corresponding dual constraint has not been used.
- The optimum value for  $PRDMAL^{2}$  is at most  $PP = e^{2}DPL$
- We can compute the corresponding solution in polytime.

If the value of the computed dual solution (which may be infeasible) is z then

#### $OPT \le z \le (1 + \epsilon')OPT$

- The constraints used when computing z certify that the solution is feasible for DUAL'.
- Suppose that we drop all unused constraints in DUAL. We will compute the same solution feasible for DUAL'.
- ► Let DUAL'' be DUAL without unused constraints.
- The dual to DUAL" is PRIMAL where we ignore variables for which the corresponding dual constraint has not been used.
- The optimum value for PRIMAL'' is at most  $(1 + \epsilon')$ OPT.
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#### This gives that overall we need at most

## $(1 + \epsilon')$ OPT<sub>LP</sub> $(I) + O(\log^2(SIZE(I)))$

#### bins.

We can choose  $\epsilon' = \frac{1}{OPT}$  as  $OPT \le \#$ items and since we have a fully polynomial time approximation scheme (FPTAS) for knapsack.



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#### Lemma 87 (Chernoff Bounds)

Let  $X_1, ..., X_n$  be *n* independent 0-1 random variables, not necessarily identically distributed. Then for  $X = \sum_{i=1}^n X_i$  and  $\mu = E[X], L \le \mu \le U$ , and  $\delta > 0$ 

$$\Pr[X \ge (1+\delta)U] < \left(\frac{e^{\delta}}{(1+\delta)^{1+\delta}}\right)^U$$
,

and

$$\Pr[X \le (1-\delta)L] < \left(\frac{e^{-\delta}}{(1-\delta)^{1-\delta}}\right)^L ,$$



### **Lemma 88** *For* $0 \le \delta \le 1$ *we have that*

$$\left(\frac{e^{\delta}}{(1+\delta)^{1+\delta}}\right)^U \le e^{-U\delta^2/3}$$

and

$$\left(\frac{e^{-\delta}}{(1-\delta)^{1-\delta}}\right)^L \le e^{-L\delta^2/2}$$



#### Markovs Inequality:

# Let X be random variable taking non-negative values. Then

#### $\Pr[X \ge a] \le \mathbb{E}[X]/a$

Trivial!



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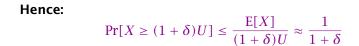
#### $\Pr[X \ge a] \le \mathbb{E}[X]/a$

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# Hence: $\Pr[X \ge (1 + \delta)U] \le \frac{\mathbb{E}[X]}{(1 + \delta)U}$







# Hence: $\Pr[X \ge (1+\delta)U] \le \frac{\mathbb{E}[X]}{(1+\delta)U} \approx \frac{1}{1+\delta}$

That's awfully weak :(



Set  $p_i = \Pr[X_i = 1]$ . Assume  $p_i > 0$  for all *i*.



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**Cool Trick:** 

 $\Pr[X \ge (1+\delta)U] = \Pr[e^{tX} \ge e^{t(1+\delta)U}]$ 



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**Cool Trick:** 

 $\Pr[X \ge (1+\delta)U] = \Pr[e^{tX} \ge e^{t(1+\delta)U}]$ 

Now, we apply Markov:

$$\Pr[e^{tX} \ge e^{t(1+\delta)U}] \le \frac{\mathrm{E}[e^{tX}]}{e^{t(1+\delta)U}} .$$



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#### This may be a lot better (!?)



 $\mathbf{E}\left[e^{tX}\right]$ 



$$\mathbf{E}\left[e^{tX}\right] = \mathbf{E}\left[e^{t\sum_{i}X_{i}}\right]$$



$$\mathbf{E}\left[e^{tX}\right] = \mathbf{E}\left[e^{t\sum_{i}X_{i}}\right] = \mathbf{E}\left[\prod_{i}e^{tX_{i}}\right]$$



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$$\mathbf{E}\left[e^{tX_i}\right] = (1 - p_i) + p_i e^t$$



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17.1 Chernoff Bounds

$$\mathbf{E}\left[e^{tX}\right] = \mathbf{E}\left[e^{t\sum_{i}X_{i}}\right] = \mathbf{E}\left[\prod_{i}e^{tX_{i}}\right] = \prod_{i}\mathbf{E}\left[e^{tX_{i}}\right]$$

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$$\prod_{i} \mathbb{E}\left[e^{tX_{i}}\right] \leq \prod_{i} e^{p_{i}(e^{t}-1)} = e^{\sum p_{i}(e^{t}-1)} = e^{(e^{t}-1)U}$$



17.1 Chernoff Bounds

 $\Pr[X \ge (1 + \delta)U] = \Pr[e^{tX} \ge e^{t(1+\delta)U}]$  $\le \frac{\operatorname{E}[e^{tX}]}{e^{t(1+\delta)U}}$ 



$$\Pr[X \ge (1+\delta)U] = \Pr[e^{tX} \ge e^{t(1+\delta)U}]$$
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We choose  $t = \ln(1 + \delta)$ .



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$$\le \frac{\operatorname{E}[e^{tX}]}{e^{t(1+\delta)U}} \le \frac{e^{(e^t-1)U}}{e^{t(1+\delta)U}} \le \left(\frac{e^{\delta}}{(1+\delta)^{1+\delta}}\right)^U$$

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#### **Lemma 89** *For* $0 \le \delta \le 1$ *we have that*

$$\left(\frac{e^{\delta}}{(1+\delta)^{1+\delta}}\right)^U \le e^{-U\delta^2/3}$$

and

$$\left(\frac{e^{-\delta}}{(1-\delta)^{1-\delta}}\right)^L \le e^{-L\delta^2/2}$$



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$$\left(\frac{e^{\delta}}{(1+\delta)^{1+\delta}}\right)^U \le e^{-U\delta^2/3}$$

Take logarithms:

$$U(\delta - (1+\delta)\ln(1+\delta)) \le -U\delta^2/3$$



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True for  $\delta = 0$ .



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Take logarithms:

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True for  $\delta = 0$ . Divide by U and take derivatives:

 $-\ln(1+\delta) \leq -2\delta/3$ 

#### Reason:

As long as derivative of left side is smaller than derivative of right side the inequality holds.



#### $f(\delta) := -\ln(1+\delta) + 2\delta/3 \le 0$



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A convex function ( $f''(\delta) \ge 0$ ) on an interval takes maximum at the boundaries.



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$$f'(\delta) = -\frac{1}{1+\delta} + 2/3$$
  $f''(\delta) = \frac{1}{(1+\delta)^2}$ 

f(0) = 0 and  $f(1) = -\ln(2) + 2/3 < 0$ 



For  $\delta \geq 1$  we show

$$\left(\frac{e^{\delta}}{(1+\delta)^{1+\delta}}\right)^U \le e^{-U\delta/3}$$



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For  $\delta \ge 1$  we show

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Take logarithms:

$$U(\delta - (1 + \delta)\ln(1 + \delta)) \le -U\delta/3$$

True for  $\delta = 0$ . Divide by *U* and take derivatives:

 $-\ln(1+\delta) \le -1/3 \iff \ln(1+\delta) \ge 1/3$  (true)

#### Reason:

As long as derivative of left side is smaller than derivative of right side the inequality holds.



$$\left(\frac{e^{-\delta}}{(1-\delta)^{1-\delta}}\right)^L \le e^{-L\delta^2/2}$$



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$$L(-\delta - (1 - \delta)\ln(1 - \delta)) \le -L\delta^2/2$$

True for  $\delta = 0$ . Divide by *L* and take derivatives:

 $\ln(1-\delta) \leq -\delta$ 

#### Reason:

As long as derivative of left side is smaller than derivative of right side the inequality holds.



$$\ln(1-\delta) \le -\delta$$



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True for  $\delta = 0$ .



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 $\ln(1-\delta) \leq -\delta$ 

True for  $\delta = 0$ . Take derivatives:

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This holds for  $0 \le \delta < 1$ .



- Given  $s_i$ - $t_i$  pairs in a graph.
- Connect each pair by a path such that not too many path use any given edge.



#### **Randomized Rounding:**

For each *i* choose one path from the set  $\mathcal{P}_i$  at random according to the probability distribution given by the Linear Programming solution.



#### **Theorem 90**

If  $W^* \ge c \ln n$  for some constant c, then with probability at least  $n^{-c/3}$  the total number of paths using any edge is at most  $W^* + \sqrt{cW^* \ln n}$ .

#### Theorem 91

With probability at least  $n^{-c/3}$  the total number of paths using any edge is at most  $W^* + c \ln n$ .



Let  $X_e^i$  be a random variable that indicates whether the path for  $s_i \cdot t_i$  uses edge e.

Then the number of paths using edge e is  $Y_e = \sum_i X_e^i$ .



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$$E[Y_e] = \sum_{i \ p \in P_i e \in P} x_p^* = \sum_{p:e \in P} x_p^* \le W^*$$



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Choose  $\delta = \sqrt{(c \ln n)/W^*}$ .

Then

 $\Pr[Y_e \ge (1+\delta)W^*] < e^{-W^*\delta^2/3} = \frac{1}{n^{c/3}}$ 



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Then

$$\Pr[Y_e \ge (1+\delta)W^*] < e^{-W^*\delta^2/3} = \frac{1}{n^{c/3}}$$



#### **Problem definition:**

- n Boolean variables
- *m* clauses  $C_1, \ldots, C_m$ . For example

 $C_7 = x_3 \vee \bar{x}_5 \vee \bar{x}_9$ 

- Non-negative weight  $w_j$  for each clause  $C_j$ .
- Find an assignment of true/false to the variables sucht that the total weight of clauses that are satisfied is maximum.



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- A variable  $x_i$  and its negation  $\bar{x}_i$  are called literals.
- ► Hence, each clause consists of a set of literals (i.e., no duplications: x<sub>i</sub> ∨ x<sub>i</sub> ∨ x̄<sub>i</sub> is not a clause).
- We assume a clause does not contain  $x_i$  and  $\bar{x}_i$  for any i.
- x<sub>i</sub> is called a positive literal while the negation x
  <sub>i</sub> is called a negative literal.
- For a given clause  $C_j$  the number of its literals is called its length or size and denoted with  $\ell_j$ .
- Clauses of length one are called unit clauses.



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- ► For a given clause  $C_j$  the number of its literals is called its length or size and denoted with  $\ell_j$ .
- Clauses of length one are called unit clauses.



# **MAXSAT: Flipping Coins**

# Set each $x_i$ independently to true with probability $\frac{1}{2}$ (and, hence, to false with probability $\frac{1}{2}$ , as well).



#### Define random variable $X_j$ with

$$X_j = \begin{cases} 1 & \text{if } C_j \text{ satisfied} \\ 0 & \text{otw.} \end{cases}$$

Then the total weight W of satisfied clauses is given by

 $W = \sum_{j} w_{j} X_{j}$ 



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#### E[W]



$$E[W] = \sum_{j} w_{j} E[X_{j}]$$



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$$E[W] = \sum_{j} w_{j} E[X_{j}]$$
$$= \sum_{j} w_{j} \Pr[C_{j} \text{ is satisified}]$$



$$E[W] = \sum_{j} w_{j} E[X_{j}]$$
  
=  $\sum_{j} w_{j} \Pr[C_{j} \text{ is satisified}]$   
=  $\sum_{j} w_{j} \left(1 - \left(\frac{1}{2}\right)^{\ell_{j}}\right)$ 



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 $\ge \frac{1}{2} \operatorname{OPT}$ 



### **MAXSAT: LP formulation**

Let for a clause C<sub>j</sub>, P<sub>j</sub> be the set of positive literals and N<sub>j</sub> the set of negative literals.

$$C_j = \bigvee_{j \in P_j} x_i \lor \bigvee_{j \in N_j} \bar{x}_i$$





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max		$\sum_j w_j z_j$		
s.t.	$\forall j$	$\sum_{i \in P_i} y_i + \sum_{i \in N_i} (1 - y_i)$	$\geq$	$z_j$
	$\forall i$	<i>v v</i>		$\{0, 1\}$
	$\forall j$	$Z_j$	$\leq$	1



# **MAXSAT: Randomized Rounding**

Set each  $x_i$  independently to true with probability  $y_i$  (and, hence, to false with probability  $(1 - y_i)$ ).



#### **Lemma 92 (Geometric Mean** $\leq$ **Arithmetic Mean)** For any nonnegative $a_1, \ldots, a_k$

$$\left(\prod_{i=1}^k a_i\right)^{1/k} \le \frac{1}{k} \sum_{i=1}^k a_i$$



#### **Definition 93**

# A function f on an interval I is concave if for any two points s and r from I and any $\lambda \in [0, 1]$ we have

#### $f(\lambda s + (1-\lambda) r) \geq \lambda f(s) + (1-\lambda) f(r)$

#### Lemma 94

Let f be a concave function on the interval [0,1], with f(0) = aand f(1) = a + b. Then

$$f(\lambda)$$

#### *for* $\lambda \in [0, 1]$ .



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**Lemma 94** Let f be a concave function on the interval [0,1], with f(0) = aand f(1) = a + b. Then

> $f(\lambda) = f((1 - \lambda)0 + \lambda 1)$   $\geq (1 - \lambda)f(0) + \lambda f(1)$  $= a + \lambda b$

for  $\lambda \in [0,1]$ .



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 $\Pr[C_j \text{ not satisfied}]$ 



 $\Pr[C_j \text{ not satisfied}] = \prod_{i \in P_j} (1 - y_i) \prod_{i \in N_j} y_i$ 



$$\Pr[C_j \text{ not satisfied}] = \prod_{i \in P_j} (1 - y_i) \prod_{i \in N_j} y_i$$
$$\leq \left[ \frac{1}{\ell_j} \left( \sum_{i \in P_j} (1 - y_i) + \sum_{i \in N_j} y_i \right) \right]^{\ell_j}$$



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$$= \left[ 1 - \frac{1}{\ell_j} \left( \sum_{i \in P_j} y_i + \sum_{i \in N_j} (1 - y_i) \right) \right]^{\ell_j}$$
$$\leq \left( 1 - \frac{z_j}{\ell_j} \right)^{\ell_j}.$$



The function  $f(z) = 1 - (1 - \frac{z}{\ell})^\ell$  is concave. Hence,

 $\Pr[C_j \text{ satisfied}]$ 



The function  $f(z) = 1 - (1 - \frac{z}{\ell})^{\ell}$  is concave. Hence,

$$\Pr[C_j \text{ satisfied}] \ge 1 - \left(1 - \frac{z_j}{\ell_j}\right)^{\ell_j}$$



The function  $f(z) = 1 - (1 - \frac{z}{\ell})^{\ell}$  is concave. Hence,

$$\Pr[C_j \text{ satisfied}] \ge 1 - \left(1 - \frac{z_j}{\ell_j}\right)^{\ell_j}$$
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The function  $f(z) = 1 - (1 - \frac{z}{\ell})^{\ell}$  is concave. Hence,

$$\Pr[C_j \text{ satisfied}] \ge 1 - \left(1 - \frac{z_j}{\ell_j}\right)^{\ell_j}$$
$$\ge \left[1 - \left(1 - \frac{1}{\ell_j}\right)^{\ell_j}\right] \cdot z_j .$$

$$f''(z) = -\frac{\ell-1}{\ell} \Big[ 1 - \frac{z}{\ell} \Big]^{\ell-2} \le 0$$
 for  $z \in [0,1]$ . Therefore,  $f$  is concave.



## E[W]



$$E[W] = \sum_{j} w_{j} \Pr[C_{j} \text{ is satisfied}]$$



$$E[W] = \sum_{j} w_{j} \Pr[C_{j} \text{ is satisfied}]$$
  
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$$\geq \left(1 - \frac{1}{e}\right) \text{ OPT }.$$



## MAXSAT: The better of two

### **Theorem 95**

# Choosing the better of the two solutions given by randomized rounding and coin flipping yields a $\frac{3}{4}$ -approximation.



 $E[\max\{W_1, W_2\}]$ 

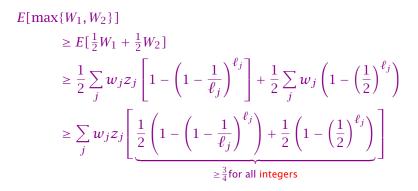


```
E[\max\{W_1, W_2\}] \\ \ge E[\frac{1}{2}W_1 + \frac{1}{2}W_2]
```

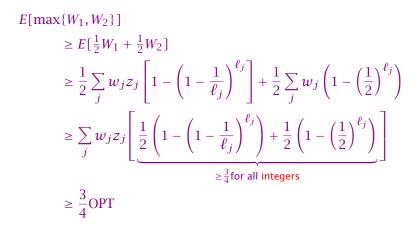


$$E[\max\{W_1, W_2\}] \\ \ge E[\frac{1}{2}W_1 + \frac{1}{2}W_2] \\ \ge \frac{1}{2}\sum_j w_j z_j \left[1 - \left(1 - \frac{1}{\ell_j}\right)^{\ell_j}\right] + \frac{1}{2}\sum_j w_j \left(1 - \left(\frac{1}{2}\right)^{\ell_j}\right)$$

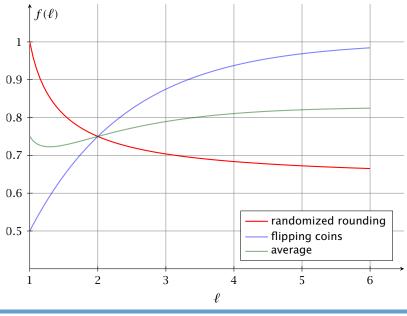












EADS II Harald Räcke

## **MAXSAT: Nonlinear Randomized Rounding**

So far we used linear randomized rounding, i.e., the probability that a variable is set to 1/true was exactly the value of the corresponding variable in the linear program.

We could define a function  $f : [0,1] \rightarrow [0,1]$  and set  $x_i$  to true with probability  $f(y_i)$ .



So far we used linear randomized rounding, i.e., the probability that a variable is set to 1/true was exactly the value of the corresponding variable in the linear program.

We could define a function  $f : [0,1] \rightarrow [0,1]$  and set  $x_i$  to true with probability  $f(y_i)$ .



## **MAXSAT: Nonlinear Randomized Rounding**

Let  $f : [0,1] \rightarrow [0,1]$  be a function with

 $1 - 4^{-x} \le f(x) \le 4^{x-1}$ 

#### Theorem 96

Rounding the LP-solution with a function f of the above form gives a  $\frac{3}{4}$ -approximation.



## **MAXSAT: Nonlinear Randomized Rounding**

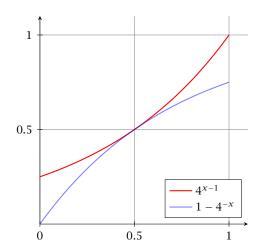
Let  $f : [0,1] \rightarrow [0,1]$  be a function with

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#### **Theorem 96**

Rounding the LP-solution with a function f of the above form gives a  $\frac{3}{4}$ -approximation.







## $\Pr[C_j \text{ not satisfied}]$



# $\Pr[C_j \text{ not satisfied}] = \prod_{i \in P_j} (1 - f(y_i)) \prod_{i \in N_j} f(y_i)$



$$\Pr[C_j \text{ not satisfied}] = \prod_{i \in P_j} (1 - f(\gamma_i)) \prod_{i \in N_j} f(\gamma_i)$$
$$\leq \prod_{i \in P_j} 4^{-\gamma_i} \prod_{i \in N_j} 4^{\gamma_i - 1}$$



$$\Pr[C_j \text{ not satisfied}] = \prod_{i \in P_j} (1 - f(y_i)) \prod_{i \in N_j} f(y_i)$$
$$\leq \prod_{i \in P_j} 4^{-y_i} \prod_{i \in N_j} 4^{y_i - 1}$$
$$= 4^{-(\sum_{i \in P_j} y_i + \sum_{i \in N_j} (1 - y_i))}$$



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$$\leq 4^{-z_j}$$



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Therefore,

E[W]



$$\Pr[C_j \text{ satisfied}] \ge 1 - 4^{-z_j} \ge \frac{3}{4} z_j$$
.

Therefore,

 $E[W] = \sum_{j} w_{j} \Pr[C_{j} \text{ satisfied}]$ 



$$\Pr[C_j \text{ satisfied}] \ge 1 - 4^{-z_j} \ge \frac{3}{4} z_j$$
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Therefore,

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$$\Pr[C_j \text{ satisfied}] \ge 1 - 4^{-z_j} \ge \frac{3}{4} z_j$$
.

Therefore,

$$E[W] = \sum_{j} w_{j} \Pr[C_{j} \text{ satisfied}] \ge \frac{3}{4} \sum_{j} w_{j} z_{j} \ge \frac{3}{4} \operatorname{OPT}$$



Not if we compare ourselves to the value of an optimum LP-solution.

## **Definition 97 (Integrality Gap)**

The integrality gap for an ILP is the worst-case ratio over all instances of the problem of the value of an optimal IP-solution to the value of an optimal solution to its linear programming relaxation.

Note that the integrality is less than one for maximization problems and larger than one for minimization problems (of course, equality is possible).

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Note that the integrality is less than one for maximization problems and larger than one for minimization problems (of course, equality is possible).

#### Lemma 98

# Our ILP-formulation for the MAXSAT problem has integrality gap at most $\frac{3}{4}$ .

max		$\sum_j w_j z_j$		
s.t.	$\forall j$	$\sum_{i \in P_i} y_i + \sum_{i \in N_i} (1 - y_i)$	$\geq$	$z_j$
	$\forall i$	$\mathcal{Y}_i$	$\in$	$\{0, 1\}$
	$\forall j$	$z_j$	$\leq$	1

Consider:  $(x_1 \lor x_2) \land (\bar{x}_1 \lor x_2) \land (x_1 \lor \bar{x}_2) \land (\bar{x}_1 \lor \bar{x}_2)$ 

- any solution can satisfy at most 3 clauses
- ▶ we can set y<sub>1</sub> = y<sub>2</sub> = 1/2 in the LP; this allows to set z<sub>1</sub> = z<sub>2</sub> = z<sub>3</sub> = z<sub>4</sub> = 1
- ▶ hence, the LP has value 4.

#### Lemma 98

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## MaxCut

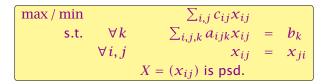
#### MaxCut

Given a weighted graph G = (V, E, w),  $w(v) \ge 0$ , partition the vertices into two parts. Maximize the weight of edges between the parts.

**Trivial 2-approximation** 



# Semidefinite Programming



- linear objective, linear contraints
- we can constrain a square matrix of variables to be symmetric positive definite

## **Vector Programming**

 $\begin{array}{cccc} \max / \min & & \sum_{i,j} c_{ij}(v_i^t v_j) \\ \text{s.t.} & \forall k & \sum_{i,j,k} a_{ijk}(v_i^t v_j) &= b_k \\ & \forall i,j & & x_{ij} &= x_{ji} \\ & & v_i \in \mathbb{R}^n \end{array}$ 

- variables are vectors in n-dimensional space
- objective functions and contraints are linear in inner products of the vectors

#### This is equivalent!



## Fact [without proof]

We (essentially) can solve Semidefinite Programs in polynomial time...



## **Quadratic Programs**

### **Quadratic Program for MaxCut:**

$$\begin{array}{ccc} \max & \frac{1}{2} \sum_{i,j} w_{ij} (1 - y_i y_j) \\ \forall i & y_i \in \{-1,1\} \end{array}$$

#### This is exactly MaxCut!



# **Semidefinite Relaxation**

max		$\frac{1}{2}\sum_{i,j}w_{ij}(1-v_i^t v_j)$		
	∀i	$v_i^t v_i$	=	1
	$\forall i$	$v_i$	$\in$	$\mathbb{R}^{n}$

- this is clearly a relaxation
- the solution will be vectors on the unit sphere



- Choose a random vector r such that r/||r|| is uniformly distributed on the unit sphere.
- If  $r^t v_i > 0$  set  $y_i = 1$  else set  $y_i = -1$



Choose the *i*-th coordinate  $r_i$  as a Gaussian with mean 0 and variance 1, i.e.,  $r_i \sim \mathcal{N}(0, 1)$ .

Density function:

$$\varphi(x) = \frac{1}{\sqrt{2\pi}} e^{x^2/2}$$

Then

$$\Pr[r = (x_1, \dots, x_n)]$$
  
=  $\frac{1}{(\sqrt{2\pi})^n} e^{x_1^2/2} \cdot e^{x_2^2/2} \cdot \dots \cdot e^{x_n^2/2} dx_1 \cdot \dots \cdot dx_n$   
=  $\frac{1}{(\sqrt{2\pi})^n} e^{\frac{1}{2}(x_1^2 + \dots + x_n^2)} dx_1 \cdot \dots \cdot dx_n$ 

Hence the probability for a point only depends on its distance to the origin.

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Hence the probability for a point only depends on its distance to the origin.

#### Fact

The projection of r onto two unit vectors  $e_1$  and  $e_2$  are independent and are normally distributed with mean 0 and variance 1 iff  $e_1$  and  $e_2$  are orthogonal.

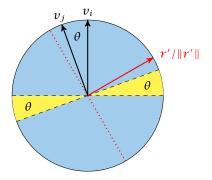
Note that this is clear if  $e_1$  and  $e_2$  are standard basis vectors.



#### Corollary

If we project r onto a hyperplane its normalized projection (r'/||r'||) is uniformly distributed on the unit circle within the hyperplane.





- if the normalized projection falls into the shaded region, v<sub>i</sub> and v<sub>j</sub> are rounded to different values
- this happens with probability  $heta/\pi$

► contribution of edge (*i*, *j*) to the SDP-relaxation:

$$\frac{1}{2}w_{ij}\left(1-v_i^t v_j\right)$$

- ▶ (expected) contribution of edge (*i*, *j*) to the rounded instance w<sub>ij</sub> arccos(v<sup>t</sup><sub>i</sub>v<sub>j</sub>)/π
- ratio is at most

 $\min_{x \in [-1,1]} \frac{2 \arccos(x)}{\pi (1-x)} \ge 0.878$ 



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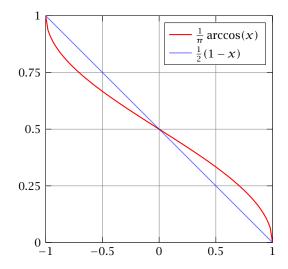
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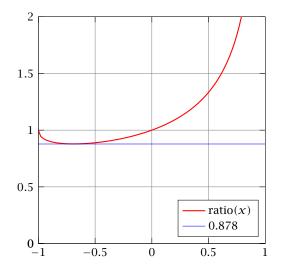
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#### Theorem 99

Given the unique games conjecture, there is no  $\alpha$ -approximation for the maximum cut problem with constant

 $\alpha > \min_{x \in [-1,1]} \frac{2 \arccos(x)}{\pi(1-x)}$ 

unless P = NP.



#### **Primal Relaxation:**

min		$\sum_{i=1}^k w_i x_i$		
s.t.	$\forall u \in U$	$\sum_{i:u\in S_i} x_i$	$\geq$	1
	$\forall i \in \{1, \dots, k\}$	$x_i$	$\geq$	0

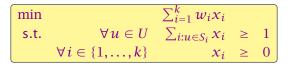
**Dual Formulation:** 

 $\begin{array}{ll} \max & \sum_{u \in U} y_u \\ \text{s.t.} \quad \forall i \in \{1, \dots, k\} \quad \sum_{u:u \in S_i} y_u \leq w_i \\ y_u \geq 0 \end{array}$ 



18.1 Primal Dual Revisited

#### **Primal Relaxation:**



#### **Dual Formulation:**



18.1 Primal Dual Revisited

## Algorithm:

- Start with y = 0 (feasible dual solution).
   Start with x = 0 (integral primal solution that may be infeasible).
- While x not feasible
  - Identify an element of that is not covered in current primal integral solution.
  - Increase dual variable (c) until a dual constraint becomes tight (maybe increase by U).
  - If this is the constraint for set 5, set  $\phi_1 = 0$  (add this set to your solution).



## Algorithm:

- Start with y = 0 (feasible dual solution).
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- While x not feasible
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Hence our cost is

$$\sum_{j} w_{j} x_{j} = \sum_{j} \sum_{e \in S_{j}} y_{e} = \sum_{e} |\{j : e \in S_{j}\}| \cdot y_{e}$$
$$\leq f \cdot \sum_{e} y_{e} \leq f \cdot \text{OPT}$$



18.1 Primal Dual Revisited

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This means

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This means

$$x_j > 0 \Rightarrow \sum_{e \in S_j} y_e = w_j$$

If we would also fulfill dual slackness conditions

$$y_e > 0 \Rightarrow \sum_{j:e \in S_j} x_j = 1$$

then the solution would be optimal!!!



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This is sufficient to show that the solution is an f-approximation.



Suppose we have a primal/dual pair



Suppose we have a primal/dual pair

and solutions that fulfill approximate slackness conditions:

$$x_j > 0 \Rightarrow \sum_i a_{ij} y_i \ge \frac{1}{\alpha} c_j$$
$$y_i > 0 \Rightarrow \sum_j a_{ij} x_j \le \beta b_i$$



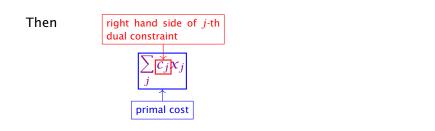
18.1 Primal Dual Revisited













$$\frac{\sum_{j} c_{j} x_{j}}{\uparrow} \leq \alpha \sum_{j} \left( \sum_{i} a_{ij} y_{i} \right) x_{j}$$
primal cost



$$\sum_{j} c_{j} x_{j} \leq \alpha \sum_{j} \left( \sum_{i} a_{ij} y_{i} \right) x_{j}$$

$$\uparrow$$

$$primal cost} \alpha \sum_{i} \left( \sum_{j} a_{ij} x_{j} \right) y_{i}$$



$$\frac{\sum_{j} c_{j} x_{j}}{\uparrow} \leq \alpha \sum_{j} \left( \sum_{i} a_{ij} y_{i} \right) x_{j}$$

$$\xrightarrow{\text{primal cost}} \alpha \sum_{i} \left( \sum_{j} a_{ij} x_{j} \right) y_{i}$$

$$\leq \alpha \beta \cdot \sum_{i} b_{i} y_{i}$$



$$\frac{\sum_{j} c_{j} x_{j}}{\uparrow} \leq \alpha \sum_{j} \left( \sum_{i} a_{ij} y_{i} \right) x_{j}$$

$$\stackrel{\text{primal cost}}{=} \alpha \sum_{i} \left( \sum_{j} a_{ij} x_{j} \right) y_{i}$$

$$\leq \alpha \beta \cdot \sum_{i} b_{i} y_{i}$$

$$\stackrel{\uparrow}{=}$$

$$\text{dual objective}$$



### Feedback Vertex Set for Undirected Graphs

• Given a graph G = (V, E) and non-negative weights  $w_v \ge 0$  for vertex  $v \in V$ .



## Feedback Vertex Set for Undirected Graphs

- Given a graph G = (V, E) and non-negative weights  $w_v \ge 0$  for vertex  $v \in V$ .
- Choose a minimum cost subset of vertices s.t. every cycle contains at least one vertex.



We can encode this as an instance of Set Cover

 Each vertex can be viewed as a set that contains some cycles.



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- However, this encoding gives a Set Cover instance of non-polynomial size.



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- Each vertex can be viewed as a set that contains some cycles.
- However, this encoding gives a Set Cover instance of non-polynomial size.
- The O(log n)-approximation for Set Cover does not help us to get a good solution.



Let  $\mathbb C$  denote the set of all cycles (where a cycle is identified by its set of vertices)



18.2 Feedback Vertex Set for Undirected Graphs

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**Primal Relaxation:** 

$$\begin{array}{|c|c|c|c|c|} \min & & \sum_{v} w_{v} x_{v} \\ \text{s.t.} & \forall C \in \mathfrak{C} & \sum_{v \in C} x_{v} \geq 1 \\ & \forall v & x_{v} \geq 0 \end{array}$$

**Dual Formulation:** 



18.2 Feedback Vertex Set for Undirected Graphs

• Start with x = 0 and y = 0



- Start with x = 0 and y = 0
- While there is a cycle C that is not covered (does not contain a chosen vertex).



- Start with x = 0 and y = 0
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- Start with x = 0 and y = 0
- While there is a cycle C that is not covered (does not contain a chosen vertex).
  - Increase y<sub>C</sub> until dual constraint for some vertex v becomes tight.
  - set  $x_v = 1$ .



 $\sum_{v} w_{v} x_{v}$ 



18.2 Feedback Vertex Set for Undirected Graphs

$$\sum_{v} w_{v} x_{v} = \sum_{v} \sum_{C: v \in C} y_{C} x_{v}$$



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where S is the set of vertices we choose.



18.2 Feedback Vertex Set for Undirected Graphs

$$\sum_{v} w_{v} x_{v} = \sum_{v} \sum_{C:v \in C} y_{C} x_{v}$$
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where S is the set of vertices we choose.



18.2 Feedback Vertex Set for Undirected Graphs

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$$= \sum_{C} |S \cap C| \cdot y_{C}$$

where S is the set of vertices we choose.

If every cycle is short we get a good approximation ratio, but this is unrealistic.



# Algorithm 1 FeedbackVertexSet

- 1:  $y \leftarrow 0$
- 2:  $x \leftarrow 0$
- 3: while exists cycle C in G do
- 4: increase  $y_C$  until there is  $v \in C$  s.t.  $\sum_{C:v \in C} y_C = w_v$

5: 
$$x_v = 1$$

- 6: remove v from G
- 7: repeatedly remove vertices of degree 1 from G



#### Idea:

Always choose a short cycle that is not covered. If we always find a cycle of length at most  $\alpha$  we get an  $\alpha$ -approximation.



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Always choose a short cycle that is not covered. If we always find a cycle of length at most  $\alpha$  we get an  $\alpha$ -approximation.

#### **Observation:**

For any path P of vertices of degree 2 in G the algorithm chooses at most one vertex from P.



# **Observation:**

If we always choose a cycle for which the number of vertices of degree at least 3 is at most  $\alpha$  we get a  $2\alpha$ -approximation.



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If we always choose a cycle for which the number of vertices of degree at least 3 is at most  $\alpha$  we get a  $2\alpha$ -approximation.

# Theorem 100

In any graph with no vertices of degree 1, there always exists a cycle that has at most  $O(\log n)$  vertices of degree 3 or more. We can find such a cycle in linear time.

This means we have

 $\mathcal{Y}_C > 0 \Rightarrow |S \cap C| \leq \mathcal{O}(\log n)$  .



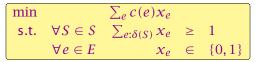
Given a graph G = (V, E) with two nodes  $s, t \in V$  and edge-weights  $c : E \to \mathbb{R}^+$  find a shortest path between s and tw.r.t. edge-weights c.



Here  $\delta(S)$  denotes the set of edges with exactly one end-point in S, and  $S = \{S \subseteq V : s \in S, t \notin S\}$ .



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# The Dual:

max		$\sum_{S} \gamma_{S}$		
s.t.	$\forall e \in E$	$\sum_{S:e\in\delta(S)} \mathcal{Y}_S$	$\leq$	c(e)
	$\forall S \in S$	$\mathcal{Y}S$	$\geq$	0

Here  $\delta(S)$  denotes the set of edges with exactly one end-point in S, and  $S = \{S \subseteq V : s \in S, t \notin S\}$ .



18.3 Primal Dual for Shortest Path

The Dual:

Here  $\delta(S)$  denotes the set of edges with exactly one end-point in S, and  $S = \{S \subseteq V : s \in S, t \notin S\}$ .



- We can interpret the value  $y_S$  as the width of a moat surounding the set S.
- Each set can have its own moat but all moats must be disjoint.
- An edge cannot be shorter than all the moats that it has to cross.



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# Algorithm 1 PrimalDualShortestPath

- 1:  $\gamma \leftarrow 0$
- 2:  $F \leftarrow \emptyset$
- 3: while there is no *s*-*t* path in (V, F) do
- Let C be the connected component of (V, F) con-4: taining s
- 5: Increase  $\gamma_C$  until there is an edge  $e' \in \delta(C)$  such that  $\sum_{S:e' \in \delta(S)} y_S = c(e')$ . 6

$$F \leftarrow F \cup \{e'\}$$

7: Let *P* be an *s*-*t* path in (V, F)

```
8: return P
```



# **Lemma 101** At each point in time the set F forms a tree.

Proof:

- In each iteration we take the current connected component from (2007) that contains 2 (call this component C) and add some edge from (2007) to 2007.
- Since, at most one end-point of the new edge is in 12 the edge cannot close a cycle.



### Lemma 101

At each point in time the set F forms a tree.

# Proof:

- In each iteration we take the current connected component from (V, F) that contains *s* (call this component *C*) and add some edge from  $\delta(C)$  to *F*.
- Since, at most one end-point of the new edge is in C the edge cannot close a cycle.



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At each point in time the set F forms a tree.

# Proof:

- ► In each iteration we take the current connected component from (V, F) that contains *s* (call this component *C*) and add some edge from  $\delta(C)$  to *F*.
- Since, at most one end-point of the new edge is in C the edge cannot close a cycle.







$$\sum_{e \in P} c(e) = \sum_{e \in P} \sum_{S: e \in \delta(S)} \mathcal{Y}_S$$



$$\sum_{e \in P} c(e) = \sum_{e \in P} \sum_{S: e \in \delta(S)} \gamma_S$$
$$= \sum_{S: s \in S, t \notin S} |P \cap \delta(S)| \cdot \gamma_S .$$



$$\sum_{e \in P} c(e) = \sum_{e \in P} \sum_{S: e \in \delta(S)} y_S$$
$$= \sum_{S: s \in S, t \notin S} |P \cap \delta(S)| \cdot y_S .$$

If we can show that  $y_S > 0$  implies  $|P \cap \delta(S)| = 1$  gives

$$\sum_{e \in P} c(e) = \sum_{S} y_{S} \le \text{OPT}$$

by weak duality.



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$$\sum_{e \in P} c(e) = \sum_{S} y_{S} \le \text{OPT}$$

by weak duality.

Hence, we find a shortest path.



When we increased  $y_S$ , S was a connected component of the set of edges F' that we had chosen till this point.

 $F' \cup P'$  contains a cycle. Hence, also the final set of edges contains a cycle.



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#### **Steiner Forest Problem:**

Given a graph G = (V, E), together with source-target pairs  $s_i, t_i$ , i = 1, ..., k, and a cost function  $c : E \to \mathbb{R}^+$  on the edges. Find a subset  $F \subseteq E$  of the edges such that for every  $i \in \{1, ..., k\}$  there is a path between  $s_i$  and  $t_i$  only using edges in F.



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$$\begin{array}{cccc} \max & \sum_{S: \exists i \text{ s.t. } S \in S_i} \mathcal{Y}S \\ \text{s.t.} & \forall e \in E & \sum_{S:e \in \delta(S)} \mathcal{Y}S &\leq c(e) \\ & & \mathcal{Y}S &\geq 0 \end{array}$$

The difference to the dual of the shortest path problem is that we have many more variables (sets for which we can generate a moat of non-zero width).



# Algorithm 1 FirstTry

1: 
$$y \leftarrow 0$$
  
2:  $F \leftarrow \emptyset$   
3: while not all  $s_i \cdot t_i$  pairs connected in  $F$  do  
4: Let  $C$  be some connected component of  $(V, F)$   
such that  $|C \cap \{s_i, t_i\}| = 1$  for some  $i$ .  
5: Increase  $y_C$  until there is an edge  $e' \in \delta(C)$  s.t.  
 $\sum_{S \in S_i: e' \in \delta(S)} y_S = C_{e'}$   
6:  $F \leftarrow F \cup \{e'\}$   
7: return  $\bigcup_i P_i$ 







$$\sum_{e \in F} c(e) = \sum_{e \in F} \sum_{S: e \in \delta(S)} \gamma_S$$



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However, this is not true:

• Take a complete graph on k + 1 vertices  $v_0, v_1, \ldots, v_k$ .



$$\sum_{e \in F} c(e) = \sum_{e \in F} \sum_{S: e \in \delta(S)} \gamma_S = \sum_S |\delta(S) \cap F| \cdot \gamma_S .$$

- Take a complete graph on k + 1 vertices  $v_0, v_1, \ldots, v_k$ .
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- The first component *C* could be  $\{v_0\}$ .
- We only set  $y_{\{v_0\}} = 1$ . All other dual variables stay 0.
- The final set *F* contains all edges  $\{v_0, v_i\}$ , i = 1, ..., k.
- $y_{\{v_0\}} > 0$  but  $|\delta(\{v_0\}) \cap F| = k$ .



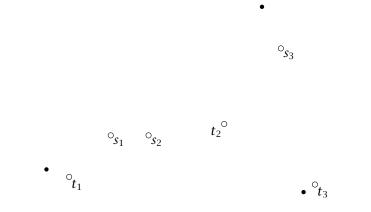
#### Algorithm 1 SecondTry

1: 
$$y \leftarrow 0$$
;  $F \leftarrow \emptyset$ ;  $\ell \leftarrow 0$   
2: while not all  $s_i \cdot t_i$  pairs connected in  $F$  do  
3:  $\ell \leftarrow \ell + 1$   
4: Let  $\mathfrak{C}$  be set of all connected components  $C$  of  $(V, F)$   
such that  $|C \cap \{s_i, t_i\}| = 1$  for some  $i$ .  
5: Increase  $y_C$  for all  $C \in \mathfrak{C}$  uniformly until for some edge  
 $e_\ell \in \delta(C'), C' \in \mathfrak{C}$  s.t.  $\sum_{S:e_\ell \in \delta(S)} y_S = c_{e_\ell}$   
6:  $F \leftarrow F \cup \{e_\ell\}$   
7:  $F' \leftarrow F$   
8: for  $k \leftarrow \ell$  downto 1 do // reverse deletion  
9: if  $F' - e_k$  is feasible solution then  
10: remove  $e_k$  from  $F'$   
11: return  $F'$ 

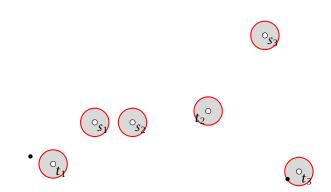


The reverse deletion step is not strictly necessary this way. It would also be sufficient to simply delete all unnecessary edges in any order.





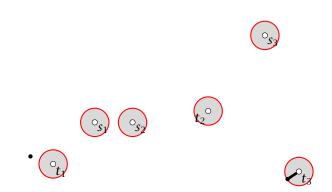








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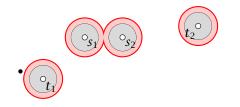






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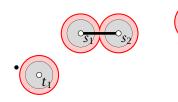










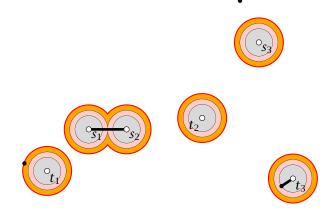






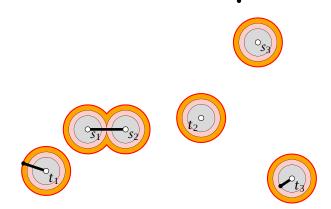


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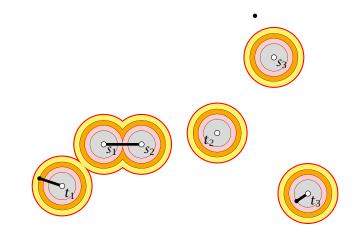






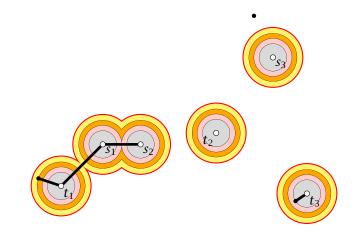






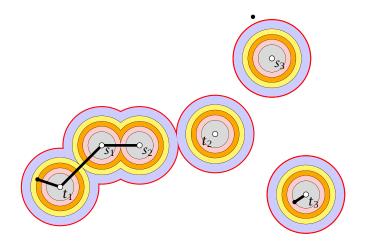




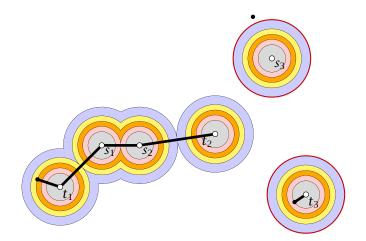




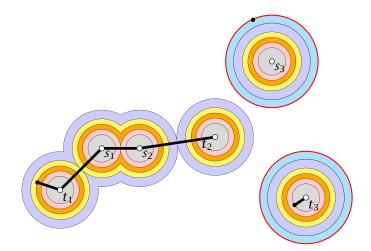




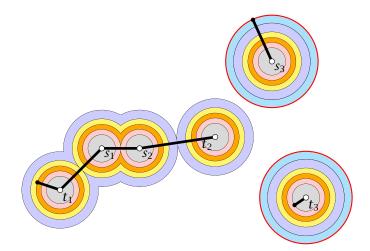




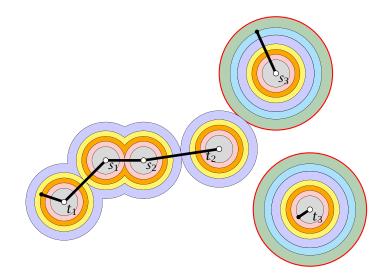






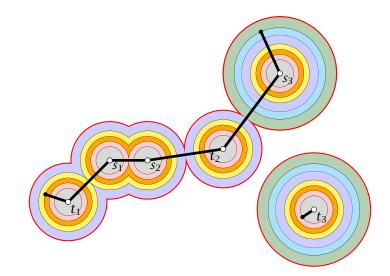




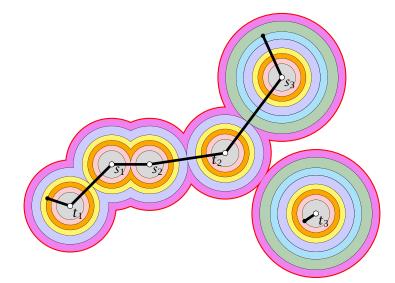






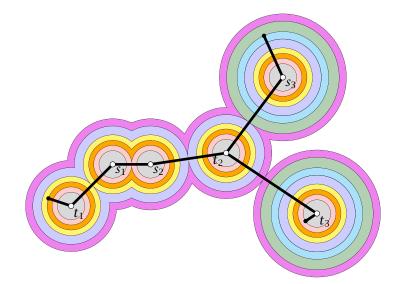




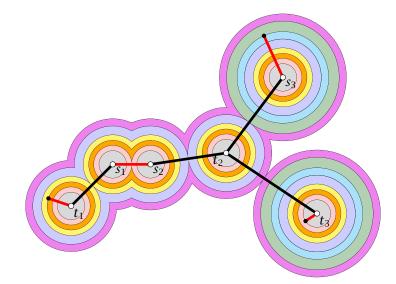














#### **Lemma 102** For any C in any iteration of the algorithm

 $\sum_{C \in \mathfrak{C}} |\delta(C) \cap F'| \le 2|\mathfrak{C}|$ 

This means that the number of times a moat from  $\mathbb{C}$  is crossed in the final solution is at most twice the number of moats.

Proof: later...



 $\sum_{e \in F'} c_e = \sum_{e \in F'} \sum_{S:e \in \delta(S)} y_S = \sum_{S} |F' \cap \delta(S)| \cdot y_S.$ 

$$\sum_{S} |F' \cap \delta(S)| \cdot y_{S} \le 2 \sum_{S} y_{S}$$

In the 1-th iteration the increase of the left-hand side is

and the increase of the right hand side is 2e(0).



$$\sum_{e \in F'} c_e = \sum_{e \in F'} \sum_{S: e \in \delta(S)} \mathcal{Y}_S = \sum_{S} |F' \cap \delta(S)| \cdot \gamma_S .$$

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In the 3-th iteration the increase of the left-hand side is

and the increase of the right hand side is 2010.



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In the *i*-th iteration the increase of the left-hand side is

 $\epsilon \sum_{C \in \mathfrak{C}} |F' \cap \delta(C)|$ 

#### and the increase of the right hand side is $2\epsilon |C|$ .



$$\sum_{e \in F'} c_e = \sum_{e \in F'} \sum_{S: e \in \delta(S)} y_S = \sum_S |F' \cap \delta(S)| \cdot y_S .$$

We want to show that

$$\sum_{S} |F' \cap \delta(S)| \cdot \gamma_{S} \le 2 \sum_{S} \gamma_{S}$$

► In the *i*-th iteration the increase of the left-hand side is

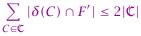
$$\epsilon \sum_{C \in \mathfrak{C}} |F' \cap \delta(C)|$$

and the increase of the right hand side is  $2\epsilon |C|$ .

Hence, by the previous lemma the inequality holds after the iteration if it holds in the beginning of the iteration.



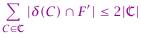
# For any set of connected components $\ensuremath{\mathbb{C}}$ in any iteration of the algorithm



- At any point during the algorithm the set of edges forms a forest (why?).
- Focileration is Let () be the set of edges in () at the beginning of the iteration.
- $\mathbf{Let} \ B = B' B_{\mathbf{h}}$
- All edges in *U* are necessary for the solution.



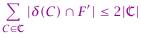
For any set of connected components  $\mathbb C$  in any iteration of the algorithm



- At any point during the algorithm the set of edges forms a forest (why?).
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- Let  $H = F' F_i$ .
- ▶ All edges in *H* are necessary for the solution.



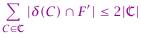
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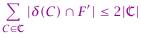
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- ► Contract all edges in *F<sub>i</sub>* into single vertices *V*′.
- ▶ We can consider the forest *H* on the set of vertices *V*′.
- Let deg(v) be the degree of a vertex  $v \in V'$  within this forest.
- Color a vertex  $v \in V'$  red if it corresponds to a component from  $\mathbb{C}$  (an active component). Otw. color it blue. (Let *B* the set of blue vertices (with non-zero degree) and *R* the set of red vertices)
- We have

$$\sum_{v \in R} \deg(v) \ge \sum_{C \in \mathbb{C}} |\delta(C) \cap F'| \stackrel{?}{\le} 2|\mathbb{C}| = 2|R|$$



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- Suppose that no node in *B* has degree one.
- Then

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  - But this means that the cluster corresponding to b must separate a source-target pair.



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  - But this means that the cluster corresponding to b must separate a source-target pair.
  - But then it must be a red node.



### **Shortest Path**

S is the set of subsets that separate s from t.

The Dual:



The Separation Problem for the Shortest Path LP is the Minimum Cut Problem.



### **Shortest Path**

$$\begin{array}{c|cccc} \min & & \sum_{e} c(e) x_{e} \\ \text{s.t.} & \forall S \in S & \sum_{e \in \delta(S)} x_{e} & \geq & 1 \\ & \forall e \in E & x_{e} & \geq & 0 \end{array}$$

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EADS II Harald Räcke

### **Minimum Cut**

### $\mathcal{P}$ is the set of path that connect s and t.

The Dual:



The Separation Problem for the Minimum Cut LP is the Shortest Path Problem.



### **Minimum Cut**

$$\begin{array}{c|cccc} \min & \sum_{e} c(e) x_{e} \\ \text{s.t.} & \forall P \in \mathcal{P} & \sum_{e \in P} x_{e} \geq 1 \\ & \forall e \in E & x_{e} \geq 0 \end{array}$$

 $\mathcal{P}$  is the set of path that connect *s* and *t*.

The Dual:

The Separation Problem for the Minimum Cut LP is the Shortest Path Problem.



### **Minimum Cut**

$$\begin{array}{c|cccc} \min & & \sum_{e} c(e) x_{e} \\ \text{s.t.} & \forall P \in \mathcal{P} & \sum_{e \in P} x_{e} & \geq & 1 \\ & \forall e \in E & & x_{e} & \geq & 0 \end{array}$$

 $\mathcal{P}$  is the set of path that connect s and t.

The Dual:

$$\begin{array}{c|ccc} \max & & \sum_{P} y_{P} \\ \text{s.t.} & \forall e \in E & \sum_{P:e \in P} y_{P} &\leq & c(e) \\ & \forall P \in \mathcal{P} & y_{P} &\geq & 0 \end{array}$$

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### Suppose that $\ell_e$ -values are solution to Minimum Cut LP.

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Let B(s, r) be the ball of radius r around s (w.r.t. metric d). Formally:

 $B = \{ v \in V \mid d(s, v) \le r \}$ 

For  $0 \le r < 1$ , B(s, r) is an *s*-*t*-cut.

Which value of r should we choose? choose randomly!!!

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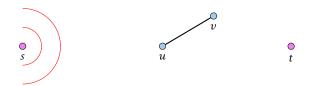
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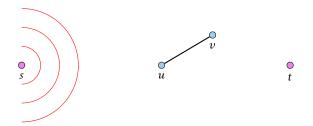




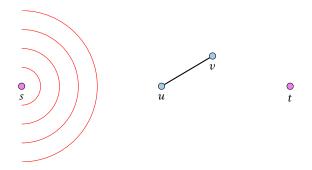




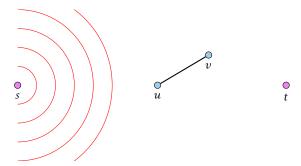




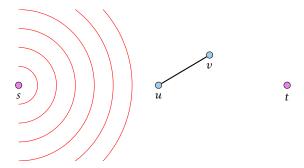




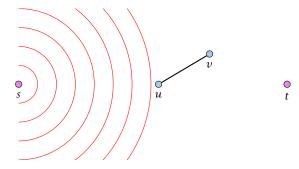




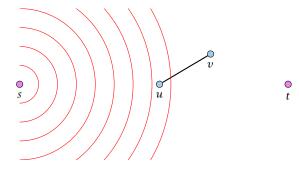




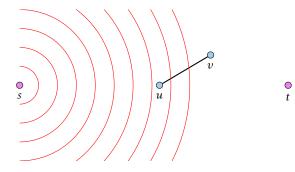




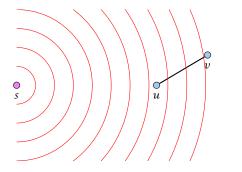








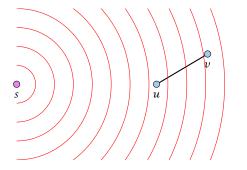






19 Cuts & Metrics

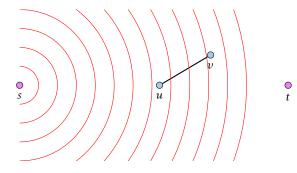
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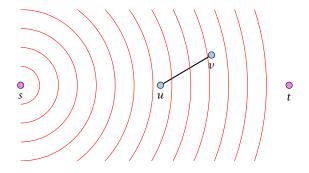


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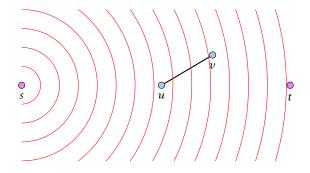
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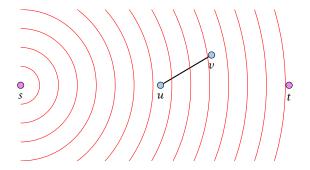








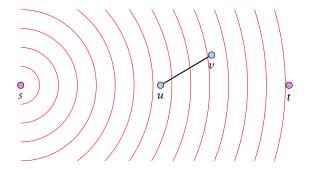




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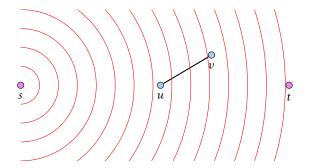




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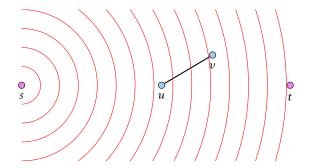




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$$E[\text{size of cut}] = E[\sum_{e} c(e) \Pr[e \text{ is cut}]]$$
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On the other hand:

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### **Minimum Multicut:**

Given a graph G = (V, E), together with source-target pairs  $s_i, t_i$ , i = 1, ..., k, and a capacity function  $c : E \to \mathbb{R}^+$  on the edges. Find a subset  $F \subseteq E$  of the edges such that all  $s_i$ - $t_i$  pairs lie in different components in  $G = (V, E \setminus F)$ .

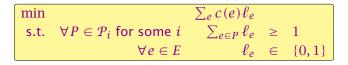


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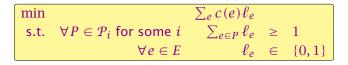


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# Re-using the analysis for the single-commodity case is difficult.

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19 Cuts & Metrics

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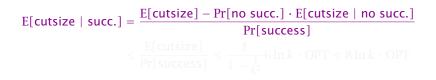


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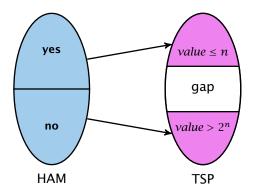
If we are not successful we simply perform a trivial *k*-approximation.

This only increases the expected cost by at most  $\frac{1}{k^2} \cdot k\text{OPT} \leq \text{OPT}/k$ .

Hence, our final cost is  $\mathcal{O}(\ln k) \cdot \text{OPT}$  in expectation.



## **Gap Introducing Reduction**



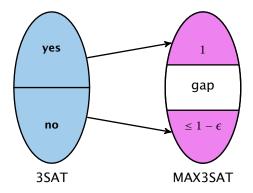
#### **Reduction from Hamiltonian cycle to TSP**

- instance that has Hamiltonian cycle is mapped to TSP instance with small cost
- otherwise it is mapped to instance with large cost
- $\Rightarrow$  there is no  $2^n/n$ -approximation for TSP

### **PCP theorem: Approximation View**

#### Theorem 104 (PCP Theorem A)

There exists  $\epsilon > 0$  for which there is gap introducing reduction between 3SAT and MAX3SAT.



## **PCP theorem: Proof System View**

#### Definition 105 (NP)

A language  $L \in NP$  if there exists a polynomial time, deterministic verifier V (a Turing machine), s.t.

- [x ∈ L] completeness $There exists a proof string <math>\mathcal{Y}$ ,  $|\mathcal{Y}| = poly(|x|)$ , s.t.  $V(x, \mathcal{Y}) =$  "accept".
- [*x* ∉ *L*] soundness For any proof string y, V(x, y) = "reject".

Note that requiring |y| = poly(|x|) for  $x \notin L$  does not make a difference (why?).



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#### Definition 106 (NP)

A language  $L \in \mathbb{NP}$  if there exists a polynomial time, deterministic verifier V (a Turing machine), s.t.

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[ $x \notin L$ ] For any proof string y, V(x, y) = "reject".

Note that requiring |y| = poly(|x|) for  $x \notin L$  does not make a difference (why?).



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An Oracle Turing Machine M is a Turing machine that has access to an oracle.

Such an oracle allows M to solve some problem in a single step.

For example having access to a TSP-oracle  $\pi_{TSP}$  would allow M to write a TSP-instance x on a special oracle tape and obtain the answer (yes or no) in a single step.

For such TMs one looks in addition to running time also at query complexity, i.e., how often the machine queries the oracle.

For a proof string y,  $\pi_y$  is an oracle that upon given an index i returns the *i*-th character  $y_i$  of y.



#### Definition 107 (PCP)

A language  $L \in PCP_{c(n),s(n)}(r(n),q(n))$  if there exists a polynomial time, non-adaptive, randomized verifier V, s.t.

- $[x \in L]$  There exists a proof string y, s.t.  $V^{\pi_y}(x) =$  "accept" with proability  $\geq c(n)$ .
- [*x* ∉ *L*] For any proof string *y*,  $V^{\pi_y}(x) =$  "accept" with probability ≤ *s*(*n*).

The verifier uses at most O(r(n)) random bits and makes at most O(q(n)) oracle queries.

c(n) is called the completeness. If not specified otw. c(n) = 1. Probability of accepting a correct proof.

s(n) < c(n) is called the soundness. If not specified otw. s(n) = 1/2. Probability of accepting a wrong proof.

r(n) is called the randomness complexity, i.e., how many random bits the (randomized) verifier uses.

q(n) is the query complexity of the verifier.



#### $\blacktriangleright P = PCP(0, 0)$

verifier without randomness and proof access is deterministic algorithm

▶  $PCP(\log n, 0) \subseteq P$ 

we can simulate collog corrandom bits in deterministic, polynomial time

▶  $PCP(0, \log n) \subseteq P$ 

we can simulate short proofs in polynomial time

►  $PCP(poly(n), 0) = coRP \stackrel{?!}{=} P$ 

by definition; collinis randomized polytime with one sided error (positive probability of accepting NO-instance)



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#### • PCP(0, poly(n)) = NP

by definition; NP-verifier does not use randomness and asks polynomially many queries

- ▶ PCP( $\log n$ , poly(n))  $\subseteq$  NP NP-verifier can simulate  $O(\log n)$  random bits
- $PCP(poly(n), 0) = coRP \stackrel{?!}{\subseteq} NP$
- NP ⊆ PCP(log n, 1) hard part of the PCP-theorem



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# **Probabilistic Checkable Proofs**

#### • PCP(0, poly(n)) = NP

by definition; NP-verifier does not use randomness and asks polynomially many queries

•  $PCP(\log n, poly(n)) \subseteq NP$ 

NP-verifier can simulate  $\mathcal{O}(\log n)$  random bits

- $PCP(poly(n), 0) = coRP \stackrel{?!}{\subseteq} NP$
- ▶ NP  $\subseteq$  PCP(log n, 1)

hard part of the PCP-theorem



## **PCP theorem: Proof System View**

## **Theorem 108 (PCP Theorem B)** NP = PCP(log n, 1)



GNI is the language of pairs of non-isomorphic graphs



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Verifier gets input  $(G_0, G_1)$  (two graphs with *n*-nodes)



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Verifier gets input  $(G_0, G_1)$  (two graphs with *n*-nodes)

It expects a proof of the following form:

► For any labeled *n*-node graph *H* the *H*'s bit *P*[*H*] of the proof fulfills

 $G_0 \equiv H \implies P[H] = 0$   $G_1 \equiv H \implies P[H] = 1$  $G_0, G_1 \neq H \implies P[H] = \text{arbitrary}$ 



- choose  $b \in \{0, 1\}$  at random
- ► take graph *G*<sup>*b*</sup> and apply a random permutation to obtain a labeled graph *H*
- check whether P[H] = b



## Verifier:

- choose  $b \in \{0, 1\}$  at random
- take graph G<sub>b</sub> and apply a random permutation to obtain a labeled graph H
- check whether P[H] = b

If  $G_0 \neq G_1$  then by using the obvious proof the verifier will always accept.



## Verifier:

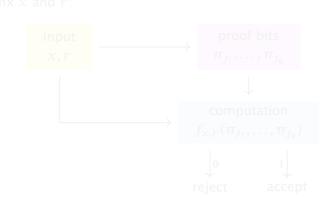
- choose  $b \in \{0, 1\}$  at random
- take graph G<sub>b</sub> and apply a random permutation to obtain a labeled graph H
- check whether P[H] = b

If  $G_0 \neq G_1$  then by using the obvious proof the verifier will always accept.

If  $G_0 \equiv G_1$  a proof only accepts with probability 1/2.

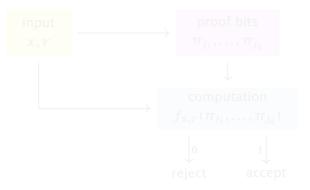
- suppose  $\pi(G_0) = G_1$
- if we accept for b = 1 and permutation  $\pi_{rand}$  we reject for b = 0 and permutation  $\pi_{rand} \circ \pi$

► For 3SAT there exists a verifier that uses c log n random bits, reads q = O(1) bits from the proof, has completeness 1 and soundness 1/2.



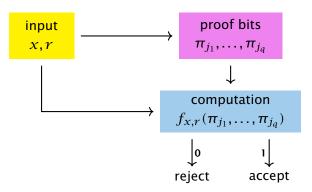


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- ▶ fix *x* and *r*:





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- fix x and r:





- transform Boolean formula f<sub>x,r</sub> into 3SAT formula C<sub>x,r</sub> (constant size, variables are proof bits)
- consider 3SAT formula  $C_x = \bigwedge_r C_{x,r}$
- $[x \in L]$  There exists proof string y, s.t. all formulas  $C_{x,r}$  evaluate to 1. Hence, all clauses in  $C_x$  satisfied.
- [ $x \notin L$ ] For any proof string  $\gamma$ , at most 50% of formulas  $C_{x,r}$  evaluate to 1. Since each contains only a constant number of clauses, a constant fraction of clauses in  $C_x$  are not satisfied.



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  - this means we have gap introducing reduction



#### We show: Version A $\implies$ NP $\subseteq$ PCP<sub>1,1- $\epsilon$ </sub>(log *n*, 1).

given  $L \in NP$  we build a PCP-verifier for L

- SSAT is NP-complete; map instance of for 0 into SSAT instance (c) s.t. (c) satisfiable iff (c) e 0
- map 1, to MAX3SAT instance (... (Comparison))
- interpret proof as assignment to variables in  $\mathbb{C}_{\mathbb{N}}$
- choose random clause 20 from (Cause 20 from Cause 20 from
- query variable assignment of for 3;
- accept if  $\mathcal{K}(m) = \operatorname{true}$  otw. reject

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- ▶ 3SAT is NP-complete; map instance x for L into 3SAT instance  $I_x$ , s.t.  $I_x$  satisfiable iff  $x \in L$
- map  $I_x$  to MAX3SAT instance  $C_x$  (PCP Thm. Version A)
- interpret proof as assignment to variables in C<sub>x</sub>
- choose random clause X from  $C_X$
- query variable assignment  $\sigma$  for X;
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- $[x \in L]$  There exists proof string  $\gamma$ , s.t. all clauses in  $C_{\chi}$  evaluate to 1. In this case the verifier returns 1.
- $[x \notin L]$  For any proof string  $\gamma$ , at most a  $(1 \epsilon)$ -fraction of clauses in  $C_x$  evaluate to 1. The verifier will reject with probability at least  $\epsilon$ .

To show Theorem B we only need to run this verifier a constant number of times to push rejection probability above 1/2.



# PCP(poly(n), 1) means we have a potentially exponentially long proof but we only read a constant number of bits from it.

The idea is to encode an NP-witness (e.g. a satisfying assignment (say n bits)) by a code whose code-words have  $2^n$  bits.

#### A wrong proof is either

- a code-word whose pre-image does not correspond to a satisfying assignment
- or, a sequence of bits that does not correspond to a code-word

We can detect both cases by querying a few positions.



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We can detect both cases by querying a few positions.



 $u \in \{0,1\}^n$  (satisfying assignment)

Walsh-Hadamard Code: WH<sub>u</sub> :  $\{0, 1\}^n \rightarrow \{0, 1\}, x \mapsto x^T u$  (over GF(2))

The code-word for u is  $WH_u$ . We identify this function by a bit-vector of length  $2^n$ .



#### **Lemma 109** If $u \neq u'$ then $WH_u$ and $WH_{u'}$ differ in at least $2^{n-1}$ bits.

**Proof:** Suppose that  $u - u' \neq 0$ . Then

#### $WH_u(x) \neq WH_{u'}(x) \iff (u - u')^T x \neq 0$

This holds for  $2^{n-1}$  different vectors x.



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This holds for  $2^{n-1}$  different vectors x.



## Suppose we are given access to a function $f : \{0, 1\}^n \rightarrow \{0, 1\}$ and want to check whether it is a codeword.

Since the set of codewords is the set of all linear functions  $\{0,1\}^n$  to  $\{0,1\}$  we can check

$$f(x + y) = f(x) + f(y)$$

for all  $2^{2n}$  pairs x, y. But that's not very efficient.



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$$f(x + y) = f(x) + f(y)$$

for all  $2^{2n}$  pairs x, y. But that's not very efficient.



Can we just check a constant number of positions?



#### **Definition 110**

Let  $\rho \in [0,1]$ . We say that  $f, g : \{0,1\}^n \to \{0,1\}$  are  $\rho$ -close if

 $\Pr_{x \in \{0,1\}^n} [f(x) = g(x)] \ge \rho \ .$ 

#### **Theorem 111 (proof deferred)** Let $f: \{0, 1\}^n \rightarrow \{0, 1\}$ with

$$\Pr_{x,y \in \{0,1\}^n} \left[ f(x) + f(y) = f(x+y) \right] \ge \rho > \frac{1}{2} \ .$$

Then there is a linear function  $ilde{f}$  such that f and  $ilde{f}$  are ho-close.



#### **Definition 110**

Let  $\rho \in [0,1]$ . We say that  $f, g : \{0,1\}^n \to \{0,1\}$  are  $\rho$ -close if

$$\Pr_{x \in \{0,1\}^n} [f(x) = g(x)] \ge \rho \; .$$

#### Theorem 111 (proof deferred)

Let  $f: \{0, 1\}^n \to \{0, 1\}$  with

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Then there is a linear function  $\tilde{f}$  such that f and  $\tilde{f}$  are  $\rho$ -close.



# We need $\mathcal{O}(1/\delta)$ trials to be sure that f is $(1 - \delta)$ -close to a linear function with (arbitrary) constant probability.



#### Suppose for $\delta < 1/4 f$ is $(1 - \delta)$ -close to some linear function $\tilde{f}$ .

 $\widehat{f}$  is uniquely defined by f , since linear functions differ on at least half their inputs.

Suppose we are given  $x \in \{0, 1\}^n$  and access to f. Can we compute  $\tilde{f}(x)$  using only constant number of queries?



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- **1.** Choose  $x' \in \{0, 1\}^n$  u.a.r.
- **2.** Set x'' := x + x'.
- **3.** Let y' = f(x') and y'' = f(x'').
- **4.** Output y' + y''.

x' and x'' are uniformly distributed (albeit dependent). With probability at least  $1 - 2\delta$  we have  $f(x') = \tilde{f}(x')$  and  $f(x'') = \tilde{f}(x'')$ .

Then the above routine returns  $\tilde{f}(x)$ .

# This technique is known as local decoding of the Walsh-Hadamard code.

Suppose we are given  $x \in \{0, 1\}^n$  and access to f. Can we compute  $\tilde{f}(x)$  using only constant number of queries?

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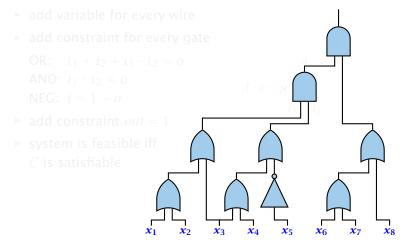
This technique is known as local decoding of the Walsh-Hadamard code.

We show that  $QUADEQ \in PCP(poly(n), 1)$ . The theorem follows since any PCP-class is closed under polynomial time reductions.

 $\ensuremath{\textbf{QUADEQ}}$  Given a system of quadratic equations over GF(2). Is there a solution?

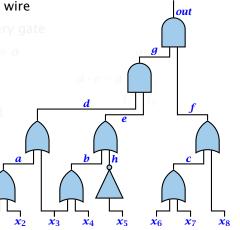


▶ given 3SAT instance *C* represent it as Boolean circuit e.g.  $C = (x_1 \lor x_2 \lor x_3) \land (x_3 \lor x_4 \lor \bar{x}_5) \land (x_6 \lor x_7 \lor x_8)$ 





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- add variable for every wire
- ► add constraint for every gat OR:  $i_1 + i_2 + i_1 \cdot i_2 = o$ AND:  $i_1 \cdot i_2 = o$ NEC: i = 1 - o
- add constraint out = 1
- system is feasible iff
   C is satisfiable

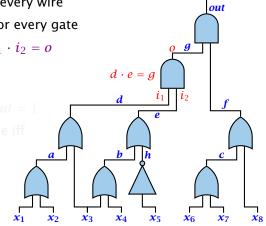




20 Hardness of Approximation

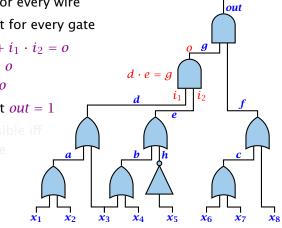
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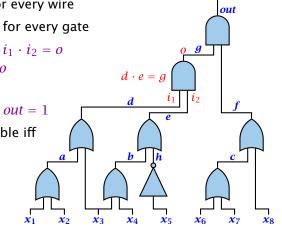




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We encode an instance of QUADEQ by a matrix A that has  $n^2$  columns; one for every pair *i*, *j*; and a right hand side vector *b*.

For an *n*-dimensional vector x we use  $x \otimes x$  to denote the  $n^2$ -dimensional vector whose i, j-th entry is  $x_i x_j$ .

Then we are asked whether

 $A(x \otimes x) = b$ 

has a solution.

Let A, b be an instance of QUADEQ. Let u be a satisfying assignment.

The correct PCP-proof will be the Walsh-Hadamard encodings of u and  $u \otimes u$ . The verifier will accept such a proof with probability 1.

We have to make sure that we reject proofs that do not correspond to codewords for vectors of the form u, and  $u \otimes u$ .

We also have to reject proofs that correspond to codewords for vectors of the form z, and  $z \otimes z$ , where z is not a satisfying assignment.



Step 1. Linearity Test.

The proof contains  $2^n + 2^{n^2}$  bits. This is interpreted as a pair of functions  $f: \{0,1\}^n \to \{0,1\}$  and  $g: \{0,1\}^{n^2} \to \{0,1\}$ .

We do a 0.999-linearity test for both functions (requires a constant number of queries).

We also assume that for the remaining constant number of accesses WH-decoding succeeds and we recover  $\tilde{f}(x)$ .

Hence, our proof will only ever see  $\tilde{f}$ . To simplify notation we use f for  $\tilde{f}$ , in the following (similar for g,  $\tilde{g}$ ).

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# Step 2. Verify that g encodes $u \otimes u$ where u is string encoded by f.

 $f(r) = u^T r$  and  $g(z) = w^T z$  since f, g are linear.

- choose r, r' independently, u.a.r. from  $\{0, 1\}^n$
- if  $f(r)f(r') \neq g(r \otimes r')$  reject
- repeat 3 times



A correct proof survives the test

 $f(\mathbf{r}) \cdot f(\mathbf{r}')$ 



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 $f(r) \cdot f(r') = u^T r \cdot u^T r'$ 



A correct proof survives the test

$$f(r) \cdot f(r') = u^T r \cdot u^T r'$$
$$= \left(\sum_i u_i r_i\right) \cdot \left(\sum_j u_j r'_j\right)$$



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$$f(r) \cdot f(r') = u^{T}r \cdot u^{T}r'$$
$$= \left(\sum_{i} u_{i}r_{i}\right) \cdot \left(\sum_{j} u_{j}r'_{j}\right)$$
$$= \sum_{ij} u_{i}u_{j}r_{i}r'_{j}$$



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If  $U \neq W$  then  $Wr' \neq Ur'$  with probability at least 1/2. Then  $r^TWr' \neq r^TUr'$  with probability at least 1/4.

## Step 3. Verify that f encodes satisfying assignment.

We need to check

 $A_k(u \otimes u) = b_k$ 

where  $A_k$  is the *k*-th row of the constraint matrix. But the left hand side is just  $g(A_k^T)$ .

We can handle this by a single query but checking all constraints would take  $\mathcal{O}(m)$  steps.

We compute  $r^T A$ , where  $r \in_R \{0,1\}^m$ . If u is not a satisfying assignment then with probability 1/2 the vector r will hit an odd number of violated constraints.

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We compute  $r^T A$ , where  $r \in_R \{0, 1\}^m$ . If u is not a satisfying assignment then with probability 1/2 the vector r will hit an odd number of violated constraints.

We used the following theorem for the linearity test:

**Theorem 111** Let  $f : \{0, 1\}^n \to \{0, 1\}$  with

$$\Pr_{x,y \in \{0,1\}^n} \left[ f(x) + f(y) = f(x+y) \right] \ge \rho > \frac{1}{2} \ .$$

Then there is a linear function  $\tilde{f}$  such that f and  $\tilde{f}$  are  $\rho$ -close.



## Fourier Transform over GF(2)

In the following we use  $\{-1,1\}$  instead of  $\{0,1\}$ . We map  $b \in \{0,1\}$  to  $(-1)^b$ .

This turns summation into multiplication.

The set of function  $f: \{-1,1\}^n \to \mathbb{R}$  form a  $2^n$ -dimensional Hilbert space.



### Hilbert space

- addition (f + g)(x) = f(x) + g(x)
- scalar multiplication  $(\alpha f)(x) = \alpha f(x)$
- ▶ inner product  $\langle f, g \rangle = E_{x \in \{-1,1\}^n}[f(x)g(x)]$ (bilinear,  $\langle f, f \rangle \ge 0$ , and  $\langle f, f \rangle = 0 \Rightarrow f = 0$ )
- completeness: any sequence x<sub>k</sub> of vectors for which

$$\sum_{k=1}^{\infty} \|x_k\| < \infty \text{ fulfills } \left\| L - \sum_{k=1}^{N} x_k \right\| \to 0$$

for some vector L.



#### standard basis

$$e_{x}(y) = \begin{cases} 1 & x = y \\ 0 & \text{otw.} \end{cases}$$

Then,  $f(x) = \sum_i \alpha_i e_i(x)$  where  $\alpha_x = f(x)$ , this means the functions  $e_i$  form a basis. This basis is orthonormal.



fourier basis

For  $\alpha \subseteq [n]$  define

 $\chi_{\alpha}(x) = \prod_{i \in \alpha} x_i$ 



fourier basis

For  $\alpha \subseteq [n]$  define

 $\chi_{\alpha}(x) = \prod_{i \in \alpha} x_i$ 

Note that

 $\langle \chi_{\alpha}, \chi_{\beta} \rangle$ 



20 Hardness of Approximation

fourier basis

For  $\alpha \subseteq [n]$  define

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Note that

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For  $\alpha \subseteq [n]$  define

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This means the  $\chi_{\alpha}$ 's also define an orthonormal basis. (since we have  $2^n$  orthonormal vectors...)



A function  $\chi_{\alpha}$  multiplies a set of  $x_i$ 's. Back in the GF(2)-world this means summing a set of  $z_i$ 's where  $x_i = (-1)^{z_i}$ .

This means the function  $\chi_{\alpha}$  correspond to linear functions in the GF(2) world.



We can write any function  $f : \{-1, 1\}^n \to \mathbb{R}$  as

$$f = \sum_{\alpha} \hat{f}_{\alpha} \chi_{\alpha}$$

We call  $\hat{f}_{\alpha}$  the  $\alpha^{th}$  Fourier coefficient.

## Lemma 112

1.  $\langle f, g \rangle = \sum_{\alpha} f_{\alpha} g_{\alpha}$ 2.  $\langle f, f \rangle = \sum_{\alpha} f_{\alpha}^2$ 

Note that for Boolean functions  $f : \{-1, 1\}^n \rightarrow \{-1, 1\}$ ,  $\langle f, f \rangle = 1$ .



# **Linearity Test**

in GF(2): We want to show that if  $Pr_{x,y}[f(x) + f(y) = f(x + y)]$  is large than f has a large agreement with a linear function.



# **Linearity Test**

in GF(2):

We want to show that if  $Pr_{x,y}[f(x) + f(y) = f(x + y)]$  is large than f has a large agreement with a linear function.

in Hilbert space: (we will prove) Suppose  $f : \{\pm 1\}^n \rightarrow \{-1, 1\}$  fulfills

$$\Pr_{x,y}[f(x)f(y) = f(x \circ y)] \ge \frac{1}{2} + \epsilon .$$

Then there is some  $\alpha \subseteq [n]$ , s.t.  $\hat{f}_{\alpha} \ge 2\epsilon$ .









 $2\epsilon \leq \hat{f}_{\alpha} = \langle f, \chi_{\alpha} \rangle$ 



 $2\epsilon \leq \hat{f}_{\alpha} = \langle f, \chi_{\alpha} \rangle = \text{agree} - \text{disagree}$ 



 $2\epsilon \leq \hat{f}_{\alpha} = \langle f, \chi_{\alpha} \rangle = \text{agree} - \text{disagree} = 2\text{agree} - 1$ 



 $2\epsilon \leq \hat{f}_{\alpha} = \langle f, \chi_{\alpha} \rangle = \text{agree} - \text{disagree} = 2\text{agree} - 1$ 

This gives that the agreement between f and  $\chi_{\alpha}$  is at least  $\frac{1}{2} + \epsilon$ .



## **Linearity Test**

$$\Pr_{x,y}[f(x \circ y) = f(x)f(y)] \ge \frac{1}{2} + \epsilon$$

means that the fraction of inputs x, y on which  $f(x \circ y)$  and f(x)f(y) agree is at least  $1/2 + \epsilon$ .

This gives

 $E_{x,y}[f(x \circ y)f(x)f(y)] = \text{agreement} - \text{disagreement}$ = 2agreement - 1  $\geq 2\epsilon$ 



# $2\epsilon \leq E_{x,y}\left[f(x\circ y)f(x)f(y)\right]$



$$\begin{aligned} 2\epsilon &\leq E_{x,y} \left[ f(x \circ y) f(x) f(y) \right] \\ &= E_{x,y} \left[ \left( \sum_{\alpha} \hat{f}_{\alpha} \chi_{\alpha}(x \circ y) \right) \cdot \left( \sum_{\beta} \hat{f}_{\beta} \chi_{\beta}(x) \right) \cdot \left( \sum_{\gamma} \hat{f}_{\gamma} \chi_{\gamma}(y) \right) \right] \end{aligned}$$



20 Hardness of Approximation

$$\begin{aligned} 2\epsilon &\leq E_{x,y} \left[ f(x \circ y) f(x) f(y) \right] \\ &= E_{x,y} \left[ \left( \sum_{\alpha} \hat{f}_{\alpha} \chi_{\alpha}(x \circ y) \right) \cdot \left( \sum_{\beta} \hat{f}_{\beta} \chi_{\beta}(x) \right) \cdot \left( \sum_{\gamma} \hat{f}_{\gamma} \chi_{\gamma}(y) \right) \right] \\ &= E_{x,y} \left[ \sum_{\alpha,\beta,\gamma} \hat{f}_{\alpha} \hat{f}_{\beta} \hat{f}_{\gamma} \chi_{\alpha}(x) \chi_{\alpha}(y) \chi_{\beta}(x) \chi_{\gamma}(y) \right] \end{aligned}$$



$$\begin{aligned} 2\epsilon &\leq E_{x,y} \left[ f(x \circ y) f(x) f(y) \right] \\ &= E_{x,y} \left[ \left( \sum_{\alpha} \hat{f}_{\alpha} \chi_{\alpha}(x \circ y) \right) \cdot \left( \sum_{\beta} \hat{f}_{\beta} \chi_{\beta}(x) \right) \cdot \left( \sum_{\gamma} \hat{f}_{\gamma} \chi_{\gamma}(y) \right) \right] \\ &= E_{x,y} \left[ \sum_{\alpha,\beta,\gamma} \hat{f}_{\alpha} \hat{f}_{\beta} \hat{f}_{\gamma} \chi_{\alpha}(x) \chi_{\alpha}(y) \chi_{\beta}(x) \chi_{\gamma}(y) \right] \\ &= \sum_{\alpha,\beta,\gamma} \hat{f}_{\alpha} \hat{f}_{\beta} \hat{f}_{\gamma} \cdot E_{x} \left[ \chi_{\alpha}(x) \chi_{\beta}(x) \right] E_{y} \left[ \chi_{\alpha}(y) \chi_{\gamma}(y) \right] \end{aligned}$$



$$\begin{aligned} 2\epsilon &\leq E_{x,y} \left[ f(x \circ y) f(x) f(y) \right] \\ &= E_{x,y} \left[ \left( \sum_{\alpha} \hat{f}_{\alpha} \chi_{\alpha}(x \circ y) \right) \cdot \left( \sum_{\beta} \hat{f}_{\beta} \chi_{\beta}(x) \right) \cdot \left( \sum_{\gamma} \hat{f}_{\gamma} \chi_{\gamma}(y) \right) \right] \\ &= E_{x,y} \left[ \sum_{\alpha,\beta,\gamma} \hat{f}_{\alpha} \hat{f}_{\beta} \hat{f}_{\gamma} \chi_{\alpha}(x) \chi_{\alpha}(y) \chi_{\beta}(x) \chi_{\gamma}(y) \right] \\ &= \sum_{\alpha,\beta,\gamma} \hat{f}_{\alpha} \hat{f}_{\beta} \hat{f}_{\gamma} \cdot E_{x} \left[ \chi_{\alpha}(x) \chi_{\beta}(x) \right] E_{y} \left[ \chi_{\alpha}(y) \chi_{\gamma}(y) \right] \\ &= \sum_{\alpha} \hat{f}_{\alpha}^{3} \end{aligned}$$



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20 Hardness of Approximation

# **Approximation Preserving Reductions**

#### **AP-reduction**

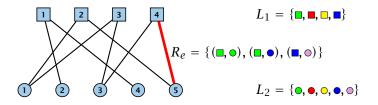
- $x \in I_1 \Rightarrow f(x,r) \in I_2$
- $SOL_1(x) \neq \emptyset \Rightarrow SOL_1(f(x,r)) \neq \emptyset$
- $y \in SOL_2(f(x,r)) \Rightarrow g(x,y,r) \in SOL_1(x)$
- f, g are polynomial time computable
- $\blacktriangleright R_2(f(x,r),y) \le r \Rightarrow R_1(x,g(x,y,r)) \le 1 + \alpha(r-1)$



## **Label Cover**

#### Input:

- bipartite graph  $G = (V_1, V_2, E)$
- ▶ label sets *L*<sub>1</sub>, *L*<sub>2</sub>
- ▶ for every edge  $(u, v) \in E$  a relation  $R_{u,v} \subseteq L_1 \times L_2$  that describe assignments that make the edge happy.
- maximize number of happy edges



## **Label Cover**

- ▶ an instance of label cover is (d<sub>1</sub>, d<sub>2</sub>)-regular if every vertex in L<sub>1</sub> has degree d<sub>1</sub> and every vertex in L<sub>2</sub> has degree d<sub>2</sub>.
- if every vertex has the same degree d the instance is called d-regular

#### Minimization version:

- assign a set L<sub>x</sub> ⊆ L<sub>1</sub> of labels to every node x ∈ L<sub>1</sub> and a set L<sub>y</sub> ⊆ L<sub>2</sub> to every node y ∈ L<sub>2</sub>
- make sure that for every edge (x, y) there is  $\ell_x \in L_x$  and  $\ell_y \in L_y$  s.t.  $(\ell_x, \ell_y) \in R_{x,y}$
- minimize  $\sum_{x \in L_1} |L_x| + \sum_{y \in L_2} |L_y|$  (total labels used)



## MAX E3SAT via Label Cover

#### instance:

 $\Phi(x) = (x_1 \vee \bar{x}_2 \vee x_3) \land (x_4 \vee x_2 \vee \bar{x}_3) \land (\bar{x}_1 \vee x_2 \vee \bar{x}_4)$ 

#### corresponding graph:



label sets:  $L_1 = \{T, F\}^3, L_2 = \{T, F\}$  (*T*=true, *F*=false)

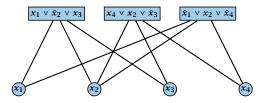
relation:  $R_{C,x_i} = \{((u_i, u_j, u_k), u_i)\}$ , where the clause C is over variables  $x_i, x_j, x_k$  and assignment  $(u_i, u_j, u_k)$  satisfies C

 $R = \{ ((F,F,F),F), ((F,T,F),F), ((F,F,T),T), ((F,T,T),T), ((T,T,T),T), ((T,T,F),F), ((T,F,F),F) \}$ 

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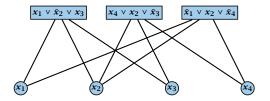
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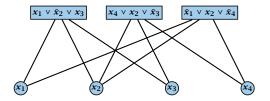
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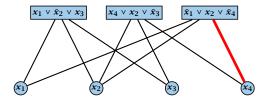
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## Lemma 113

If we can satisfy k out of m clauses in  $\phi$  we can make at least 3k + 2(m - k) edges happy.

- for 35-use the setting of the assignment that satisfies to clauses
- of for satisfied clauses in 30 use the corresponding assignment to the clause-variables (gives 30 happy edges)
- for unsatisfied clauses flip assignment of one of the variables; this makes one incident edge unhappy (gives disconsisting and the provident edge on the provident of the disconsisting and the provident edges (the second sec



## Lemma 113

If we can satisfy k out of m clauses in  $\phi$  we can make at least 3k + 2(m - k) edges happy.

- ▶ for V<sub>2</sub> use the setting of the assignment that satisfies k clauses
- for satisfied clauses in V<sub>1</sub> use the corresponding assignment to the clause-variables (gives 3k happy edges)
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#### Lemma 114

If we can satisfy at most k clauses in  $\Phi$  we can make at most 3k + 2(m - k) = 2m + k edges happy.

- the labeling of nodes in 1% gives an assignment.
- every unsatisfied clause in this assignment cannot be assigned a label that satisfies all 3 incident edges
- hence at most (im = (im = k) = (im = k)edges are happy



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# **Hardness for Label Cover**

We cannot distinguish between the following two cases

- all 3m edges can be made happy
- at most  $2m + (1 \epsilon)m = (3 \epsilon)m$  out of the 3m edges can be made happy

Hence, we cannot obtain an approximation constant  $lpha > rac{3-\epsilon}{3}$ .



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Hence, we cannot obtain an approximation constant  $\alpha > \frac{3-\epsilon}{3}$ .



# (3, 5)-regular instances

## Theorem 115

There is a constant  $\rho$  s.t. MAXE3SAT is hard to approximate with a factor of  $\rho$  even if restricted to instances where a variable appears in exactly 5 clauses.

Then our reduction has the following properties:

- the resulting Label Cover instance is (3,5)-regular
- it is hard to approximate for a constant  $\alpha < 1$
- ► given a label l<sub>1</sub> for x there is at most one label l<sub>2</sub> for y that makes edge (x, y) happy (uniqueness property)



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# (3, 5)-regular instances

The previous theorem can be obtained with a series of gap-preserving reductions:

- MAX3SAT  $\leq$  MAX3SAT( $\leq$  29)
- $MAX3SAT(\leq 29) \leq MAX3SAT(\leq 5)$
- $MAX3SAT(\leq 5) \leq MAX3SAT(= 5)$
- $MAX3SAT(= 5) \le MAXE3SAT(= 5)$

Here MAX3SAT( $\leq 29$ ) is the variant of MAX3SAT in which a variable appears in at most 29 clauses. Similar for the other problems.



# **Regular instances**

## Theorem 116

There is a constant  $\alpha < 1$  such if there is an  $\alpha$ -approximation algorithm for Label Cover on 15-regular instances than P=NP.

Given a label  $\ell_1$  for  $x \in V_1$  there is at most one label  $\ell_2$  for y that makes (x, y) happy. (uniqueness property)



We would like to increase the inapproximability for Label Cover.

In the verifier view, in order to decrease the acceptance probability of a wrong proof (or as here: a pair of wrong proofs) one could repeat the verification several times.

Unfortunately, we have a 2P1R-system, i.e., we are stuck with a single round and cannot simply repeat.

The idea is to use parallel repetition, i.e., we simply play several rounds in parallel and hope that the acceptance probability of wrong proofs goes down.



Given Label Cover instance I with  $G = (V_1, V_2, E)$ , label sets  $L_1$  and  $L_2$  we construct a new instance I':

$$V'_1 = V_1^k = V_1 \times \cdots \times V_1$$

$$V'_2 = V_2^k = V_2 \times \cdots \times V_2$$

$$L'_1 = L_1^k = L_1 \times \cdots \times L_1$$

$$L'_2 = L_2^k = L_2 \times \cdots \times L_2$$

$$E' = E^k = E \times \cdots \times E$$

An edge  $((x_1, \ldots, x_k), (y_1, \ldots, y_k))$  whose end-points are labelled by  $(\ell_1^x, \ldots, \ell_k^x)$  and  $(\ell_1^y, \ldots, \ell_k^y)$  is happy if  $(\ell_i^x, \ell_i^y) \in R_{x_i, y_i}$  for all *i*.



## If I is regular than also I'.

## If I has the uniqueness property than also I'.

## Did the gap increase?

- Suppose we have labelling (5...5) that satisfies just an or fraction of edges in ().
- How many edges are happy?



## If I is regular than also I'.

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## Did the gap increase?

- Suppose we have labelling doubt that satisfies just an or fraction of edges in 4.
- We transfer this labeling to instance is vertex (a) and agets label (2) (a) and (1) (a) vertex (a) and (a) gets label (2) (a) (a) (a) (a) (b)
- How many edges are happy?



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If I has the uniqueness property than also I'.

Did the gap increase?

- Suppose we have labelling  $\ell_1, \ell_2$  that satisfies just an  $\alpha$ -fraction of edges in *I*.
- ▶ We transfer this labelling to instance I': vertex  $(x_1,...,x_k)$  gets label  $(\ell_1(x_1),...,\ell_1(x_k))$ , vertex  $(y_1,...,y_k)$  gets label  $(\ell_2(y_1),...,\ell_2(y_k))$ .
- How many edges are happy? only control out of confill (just an of fraction).



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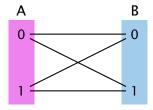


#### Non interactive agreement:

- Two provers A and B
- ▶ The verifier generates two random bits *b*<sub>*A*</sub>, and *b*<sub>*B*</sub>, and sends one to *A* and one to *B*.
- ► Each prover has to answer one of A<sub>0</sub>, A<sub>1</sub>, B<sub>0</sub>, B<sub>1</sub> with the meaning A<sub>0</sub> := prover A has been given a bit with value 0.
- The provers win if they give the same answer and if the answer is correct.



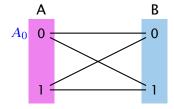
### The provers can win with probability at most 1/2.



Regardless what we do 50% of edges are unhappy!



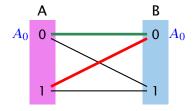
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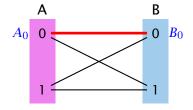
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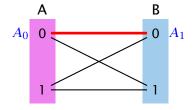
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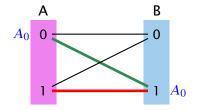
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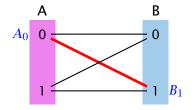
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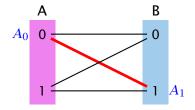
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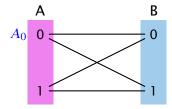


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### **Counter Example**

The provers can win with probability at most 1/2.



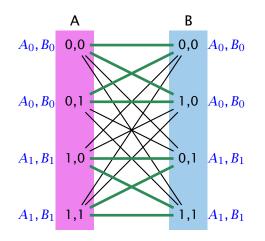
Regardless what we do 50% of edges are unhappy!



20 Hardness of Approximation

#### **Counter Example**

In the repeated game the provers can also win with probability 1/2:



### **Boosting**

#### Theorem 117

There is a constant c > 0 such if  $OPT(I) = |E|(1 - \delta)$  then  $OPT(I') \le |E'|(1 - \delta)^{\frac{ck}{\log L}}$ , where  $L = |L_1| + |L_2|$  denotes total number of labels in I.

proof is highly non-trivial



### **Boosting**

#### Theorem 117

There is a constant c > 0 such if  $OPT(I) = |E|(1 - \delta)$  then  $OPT(I') \le |E'|(1 - \delta)^{\frac{ck}{\log L}}$ , where  $L = |L_1| + |L_2|$  denotes total number of labels in I.

proof is highly non-trivial



# **Hardness of Label Cover**

#### Theorem 118

There are constants c > 0,  $\delta < 1$  s.t. for any k we cannot distinguish regular instances for Label Cover in which either

- OPT(I) = |E|, or
- OPT(*I*) =  $|E|(1 \delta)^{ck}$

# unless each problem in NP has an algorithm running in time $\mathcal{O}(n^{\mathcal{O}(k)})$ .

#### Corollary 119

There is no  $\alpha$ -approximation for Label Cover for any constant  $\alpha$ .



#### Theorem 120

There exist regular Label Cover instances s.t. we cannot distinguish whether

- all edges are satisfiable, or
- at most a  $1/\log^2(|L_1||E|)$ -fraction is satisfiable

unless NP-problems have algorithms with running time  $O(n^{O(\log \log n)})$ .

choose  $k \ge \frac{2}{c} \log_{1/(1-\delta)} (\log(|L_1||E|)) = \mathcal{O}(\log\log n)$ .



#### Partition System (s, t, h)

- universe U of size s
- ► t pairs of sets  $(A_1, \bar{A}_1), \dots, (A_t, \bar{A}_t);$  $A_i \subseteq U, \bar{A}_i = U \setminus A_i$
- choosing from any h pairs only one of A<sub>i</sub>, A<sub>i</sub> we do not cover the whole set U

#### we will show later:

for any *h*, *t* with  $h \le t$  there exist systems with  $s = |U| \le 4t^2 2^h$ 



Given a Label Cover instance we construct a Set Cover instance;

The universe is  $E \times U$ , where U is the universe of some partition system; ( $t = |L_1|$ ,  $h = \log(|E||L_1|)$ )

for all  $u \in V_1, \ell_1 \in L_1$ 

 $S_{u,\ell_1} = \{((u,v),a) \mid (u,v) \in E, a \in A_{\ell_1}\}$ 

for all  $v \in V_2, \ell_2 \in L_2$ 

 $S_{v,\ell_2} = \{((u,v),a) \mid (u,v) \in E, a \in \bar{A}_{\ell_1}, \text{ where } (\ell_1,\ell_2) \in R_{(u,v)}\}$ 



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The universe is  $E \times U$ , where U is the universe of some partition system; ( $t = |L_1|$ ,  $h = \log(|E||L_1|)$ )

for all  $u \in V_1$ ,  $\ell_1 \in L_1$   $S_{u,\ell_1} = \{((u,v),a) \mid (u,v) \in E, a \in A_{\ell_1}\}$ for all  $v \in V_2$ ,  $\ell_2 \in L_2$   $S_{v,\ell_2} = \{((u,v),a) \mid (u,v) \in E, a \in \tilde{A}_{\ell_1}, \text{ where } (\ell_1,\ell_2) \in R_{(u,v)}\}$ note that  $S_{v,\ell_2}$  is well defined because of uniqueness property



Given a Label Cover instance we construct a Set Cover instance;

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Suppose that we can make all edges happy.

Choose sets  $S_{u,\ell_1}$ 's and  $S_{v,\ell_2}$ 's, where  $\ell_1$  is the label we assigned to u, and  $\ell_2$  the label for v. ( $|V_1|+|V_2|$  sets)

For an edge (u, v),  $S_{v,\ell_2}$  contains  $\{(u, v)\} \times A_{\ell_2}$ . For a happy edge  $S_{u,\ell_1}$  contains  $\{(u, v)\} \times \overline{A}_{\ell_2}$ .

Since all edges are happy we have covered the whole universe.



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Since all edges are happy we have covered the whole universe.



#### Lemma 121

Given a solution to the set cover instance using at most  $\frac{h}{8}(|V_1| + |V_2|)$  sets we can find a solution to the Label Cover instance satisfying at least  $\frac{2}{h^2}|E|$  edges.

If the Label Cover instance cannot satisfy a  $2/h^2$ -fraction we cannot cover with  $\frac{h}{8}(|V_1| + |V_2|)$  sets.

Since differentiating between both cases for the Label Cover instance is hard, we have an O(h)-hardness for Set Cover.



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Since differentiating between both cases for the Label Cover instance is hard, we have an  $\mathcal{O}(h)$ -hardness for Set Cover.



- $n_u$ : number of  $S_{u,i}$ 's in cover
- $n_v$ : number of  $S_{v,j}$ 's in cover
- at most 1/4 of the vertices can have  $n_u, n_v \ge h/2$ ; mark these vertices
- at least half of the edges have both end-points unmarked, as the graph is regular
- ▶ for such an edge (u, v) we must have chosen  $S_{u,i}$  and a corresponding  $S_{v,j}$ , s.t.  $(i, j) \in R_{u,v}$  (making (u, v) happy)
- we choose a random label for u from the (at most h/2) chosen S<sub>u,i</sub>-sets and a random label for v from the (at most h/2) S<sub>v,j</sub>-sets
- (u, v) gets happy with probability at least  $4/h^2$
- hence we make a 2/h<sup>2</sup>-fraction of edges happy



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- hence we make a  $2/h^2$ -fraction of edges happy



#### Set Cover

#### Theorem 122

There is no  $\frac{1}{32}\log n$ -approximation for the unweighted Set Cover problem unless problems in NP can be solved in time  $\mathcal{O}(n^{\mathcal{O}(\log \log n)})$ .



Set  $h = \log(|E||L_1|)$  and  $t = |L_1|$ ; Size of partition system is  $s = |U| = 4t^22^h = 4|L_1|^2(|E||L_1|)^2 = 4|E|^2|L_1|^4$ 

The size of the ground set is then

 $n = |E||U| = 4|E|^3|L_2|^4 \le (|E||L_2|)^4$ 

for sufficiently large |E|. Then  $h \ge \frac{1}{4} \log n$ .

If we get an instance where all edges are satisfiable there exists a cover of size only  $|V_1| + |V_2|$ .

Set  $h = \log(|E||L_1|)$  and  $t = |L_1|$ ; Size of partition system is  $s = |U| = 4t^2 2^h = 4|L_1|^2 (|E||L_1|)^2 = 4|E|^2 |L_1|^4$ 

The size of the ground set is then

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If we get an instance where all edges are satisfiable there exists a cover of size only  $|V_1| + |V_2|$ .

## **Partition Systems**

#### Lemma 123

# Given h and t with $h \le t$ , there is a partition system of size $s = \ln(4t)h2^h \le 4t^22^h$ .

We pick t sets at random from the possible  $2^{|U|}$  subsets of U.

Fix a choice of h of these sets, and a choice of h bits (whether we choose  $A_i$  or  $\bar{A}_i$ ). There are  $2^h \cdot {t \choose h}$  such choices.



20 Hardness of Approximation

## **Partition Systems**

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Fix a choice of h of these sets, and a choice of h bits (whether we choose  $A_i$  or  $\bar{A}_i$ ). There are  $2^h \cdot {t \choose h}$  such choices.



The probability that an element  $u \in A_i$  is 1/2 (same for  $\overline{A}_i$ ).

The probability that u is covered is  $1 - \frac{1}{2h}$ .

The probability that all u are covered is  $(1 - \frac{1}{2h})^s$ 

The probability that there exists a choice such that all u are covered is at most

$$\binom{t}{h} 2^h \left( 1 - \frac{1}{2^h} \right)^s \le (2t)^h e^{-s/2^h} = (2t)^h \cdot e^{-h \ln(4t)} < \frac{1}{2} \ .$$



The probability that an element  $u \in A_i$  is 1/2 (same for  $\overline{A}_i$ ).

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The probability that an element  $u \in A_i$  is 1/2 (same for  $\overline{A}_i$ ).

The probability that *u* is covered is  $1 - \frac{1}{2h}$ .

The probability that all u are covered is  $(1 - \frac{1}{2h})^s$ 

The probability that there exists a choice such that all u are covered is at most

$$\binom{t}{h} 2^h \left( 1 - \frac{1}{2^h} \right)^s \le (2t)^h e^{-s/2^h} = (2t)^h \cdot e^{-h \ln(4t)} < \frac{1}{2} \ .$$



# **Advanced PCP Theorem**

#### Theorem 124

For any positive constant  $\epsilon > 0$ , it is the case that  $NP \subseteq PCP_{1-\epsilon,1/2+\epsilon}(\log n, 3)$ . Moreover, the verifier just reads three bits from the proof, and bases its decision only on the parity of these bits.

It is NP-hard to approximate a MAXE3LIN problem by a factor better than  $1/2 + \delta$ , for any constant  $\delta$ .

It is NP-hard to approximate MAX3SAT better than  $7/8 + \delta$ , for any constant  $\delta$ .

