



A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

# Introduction to the Theory of Computation

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# Introduction to the Theory of Computation

Lesson 1

## 2.5. State minimization

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- Lecturer: Dr. István Heckl, [Istvan.Heckl@gmail.com](mailto:Istvan.Heckl@gmail.com)
- <http://oktatas.mik.uni-pannon.hu>
  - registration
    - full name
    - neptun code
  - course: A számítástudomány alapjai (Heckl)
    - enlisting code: kincs
    - see: download materials
  - course: Theory of the elements of computation (Heckl)



- Subject code:
  - for full time student: VEMISAB512S
  - for part time student: VEMLSAB512S
- Signature: at least 50% result at ZH
- Subject name in Hungarian: A számítástudomány alapjai
- Literature: Harry R. Lewis, Christos H. Papadimitriou:  
Elements of the Theory of Computation, Prentice Hall, Inc.,  
1998. (second edition)
  - this presentation is based on this book
  - Bach Iván: Formális nyelvek
- Precondition: A digitális számítás elmélete

- Finite automaton:
  - state minimization
  - algorithmic aspects of finite automata
- Context-free languages: Parse trees
- Algorithms for context-free grammars, Determinism and parsing

- Extensions of Turing machines
- Random access Turing machines, Nondeterministic Turing machines
- Grammars
- Numerical functions

- Undecidability
  - unsolvable problems about Turing machines
  - unsolvable problems about grammars
  - an unsolvable tiling problem
  - properties of recursive languages
- Computational complexity: The class P, Problems from P

- Boolean satisfiability, The class NP
- Polynomial-time reductions
- Cook's Theorem
- More NP-complete problems
- Coping with NP-completeness

- Definition of:
  - deterministic finite automaton, M
  - state diagram of a DFA
  - configuration of a DFA  $M=(K, \Sigma, \delta, s, F)$
  - yield in one step of a DFA,  $|-\_M$
  - computation by DFA M
  - yield of a DFA,  $|-\_M^*$
  - word accepted by DFA
  - language accepted by DFA M,  $L(M)$
  - equivalence of two DFAs

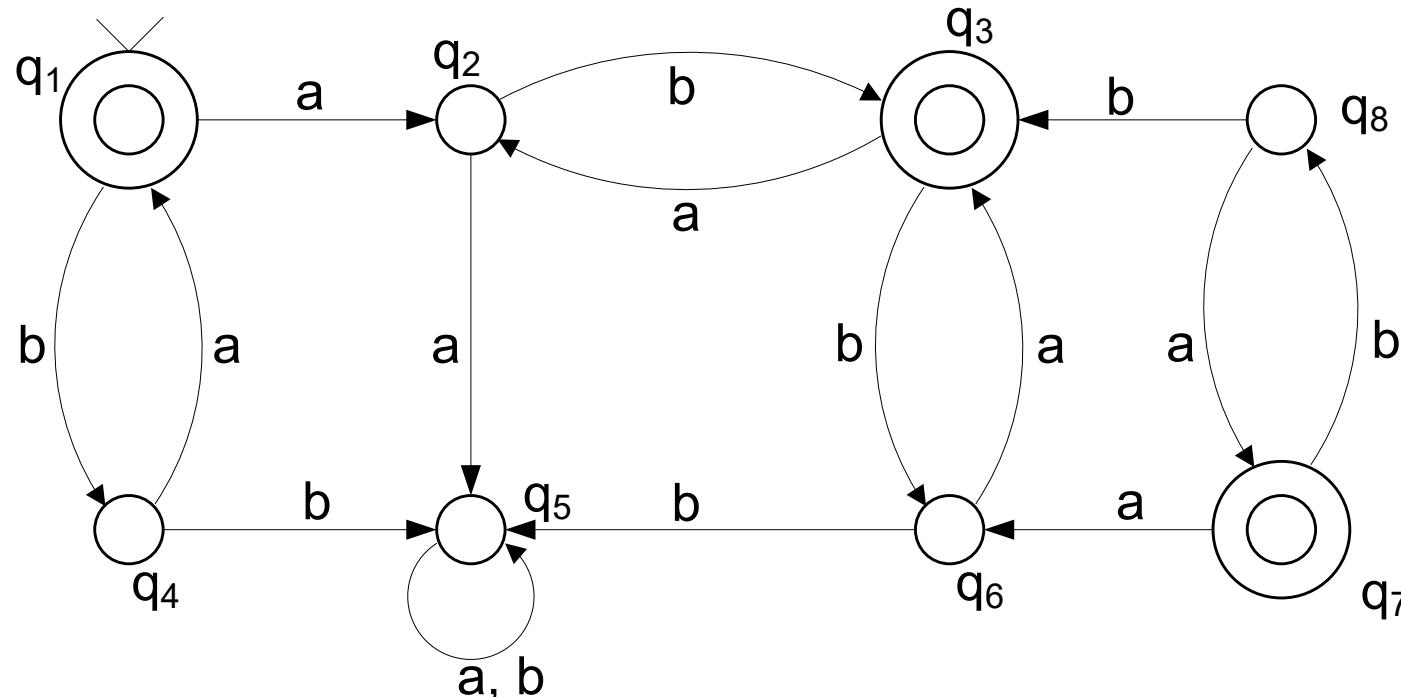
- Definition of:
  - regular expression
  - equivalence relation
  - closure
  - language

# State minimization

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- For practical reasons it is important to be able to minimize the number of states of a given DFA
- Motivating example: let DFA  $M$  such that  $L(M) = (ab \cup ba)^*$

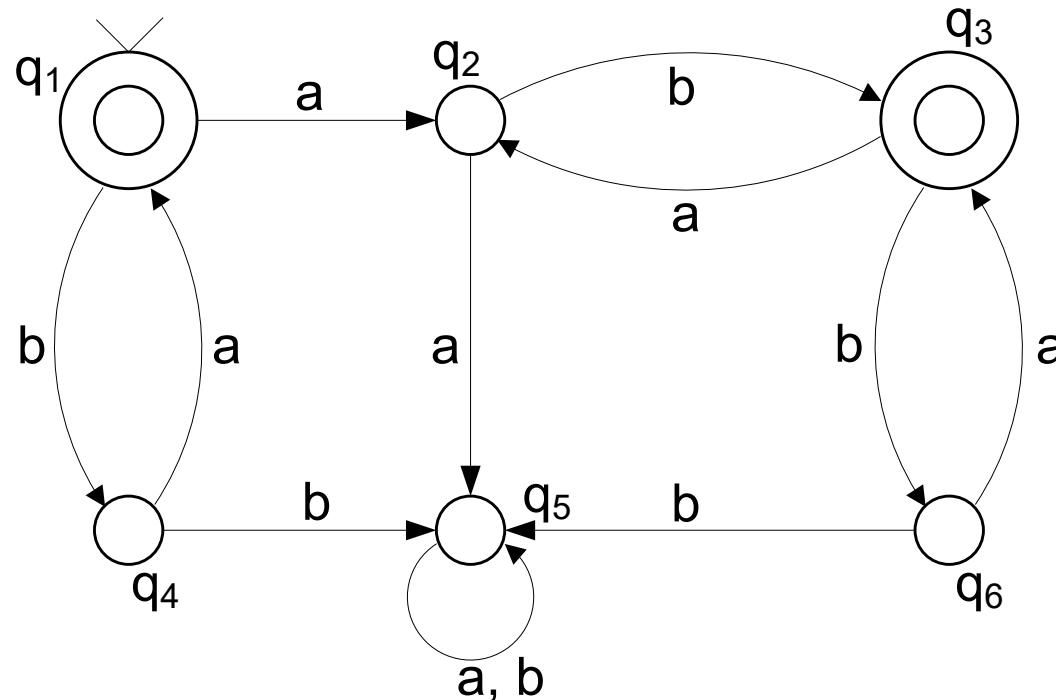


- This is the simplest kind of optimization
- State  $q_7$  and  $q_8$  are unreachable because there is no directed path from the start state to them
- Identifying the reachable states can be done in polynomial time
  - the set of reachable states can be defined as the closure of  $\{s\}$  under the relation  $\{(p, q) : \delta(p, a) = q \text{ for some } a \in \Sigma\}$
  - all closure can be implemented in polynomial time

– reachable states can be computed by this algorithm:

```
R := { s } ;
```

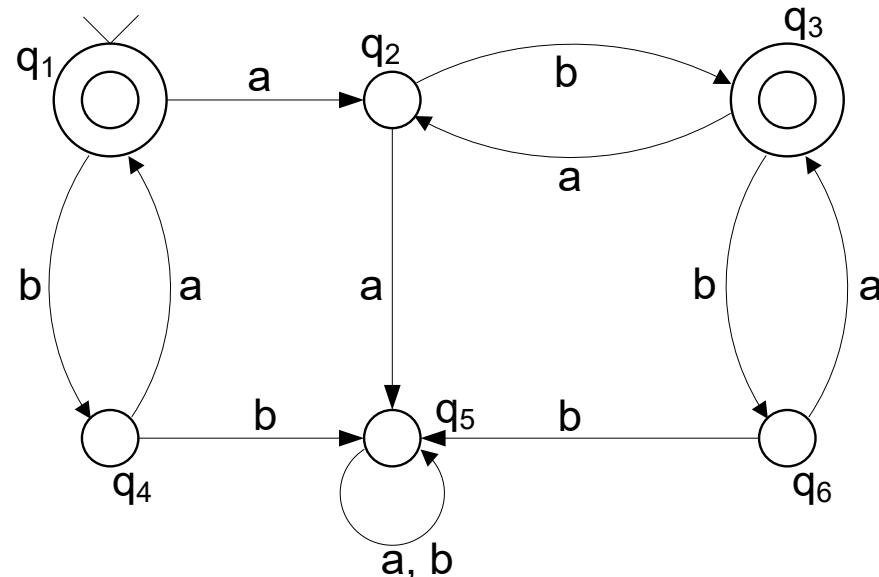
```
while there is a state p ∈ R and a ∈ Σ  
such that δ(p, σ) ∈ R do  
add δ(p, σ) to R
```



- States  $q_1$ , and  $q_3$  are equivalent
  - from either state, precisely the same strings lead the automaton to acceptance
- Equivalent states can be merged into one state

- Definition of  $x$  and  $y$  are equivalent with respect to  $L$ ,  
 $x \approx_L y$ :
  - let  $L \subseteq \Sigma^*$ ,  $x, y \in \Sigma^*$
  - if  $\forall z \in \Sigma^*, xz \in L \leftrightarrow yz \in L \rightarrow x \approx_L y$ 
    - either both concatenated strings are in  $L$  or both are not in  $L$
  - DFA is not mentioned here, you cannot check all  $z$
  - we don't know if  $xz$  and  $yz$  drive the DFA into the same state

- E.g.:
  - $ab \approx_L baba$ 
    - $z = aa \rightarrow abaa \notin L, babaaa \notin L$
    - $z = ba \rightarrow abba \in L, bababa \in L$
  - $a \approx_L baaba$  (we do not say that  $x, y \in L$ )
    - $z = ab \rightarrow aab \notin L, baabaab \notin L$
    - $z = b \rightarrow ab \in L, baabab \in L$

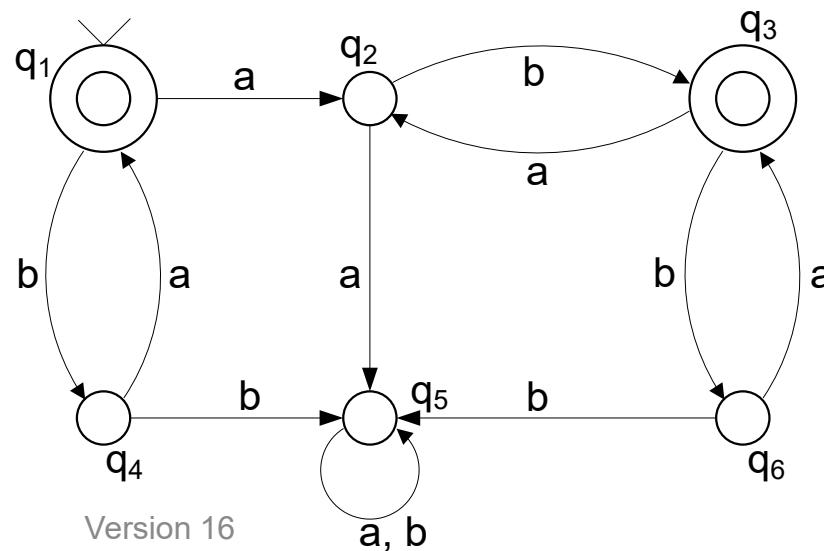


# Equivalent strings

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $\approx_L$  is an equivalence relation, so it defines equivalence classes
  - $\approx_L$  the relation holds between any two elements of a class
  - e.g.:  $e \approx_L ab \approx_L ba \approx_L abab \approx_L \dots$
  - an equivalence class is denoted with its simplest element (lexicographically)
    - e.g.:  $[e]$

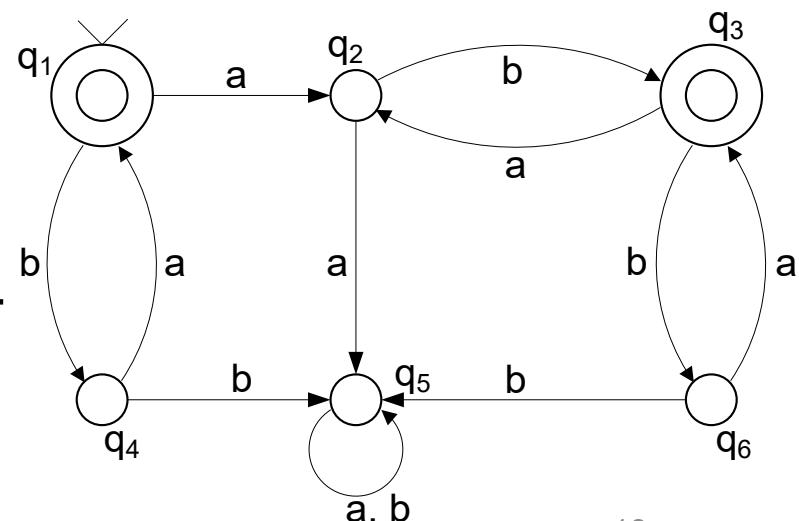


# Equivalent strings

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
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- The equivalence classes of  $\approx_L$ :
  - $[e] = L, \{q_1, q_3\}$ 
    - the elements of the language
  - $[a] = La, \{q_2\}$ 
    - a word of the language followed by 'a'
    - any  $x \in La$  needs  $z$  of the form  $bL$  in order for  $xz$  to be in  $L$
  - $[b] = Lb, \{q_4, q_6\}$
  - $[aa] = L(aa \cup bb)\Sigma^*, \{q_5\}$ 
    - no matter what  $z$  is  $xz \notin L$



- Definition of  $x$  and  $y$  are equivalent with respect to  $M$ ,  $\sim_M$ :
  - let  $M = (K, \Sigma, \delta, s, F)$  be a DFA,  $x, y \in \Sigma^*$
  - if  $\exists q \in K$  such that  $(s, x) \vdash_M^* (q, e)$ ,  $(s, y) \vdash_M^* (q, e) \rightarrow x \sim_M y$ 
    - both strings drive  $M$  from  $s$  to the same state

# Equivalent strings

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
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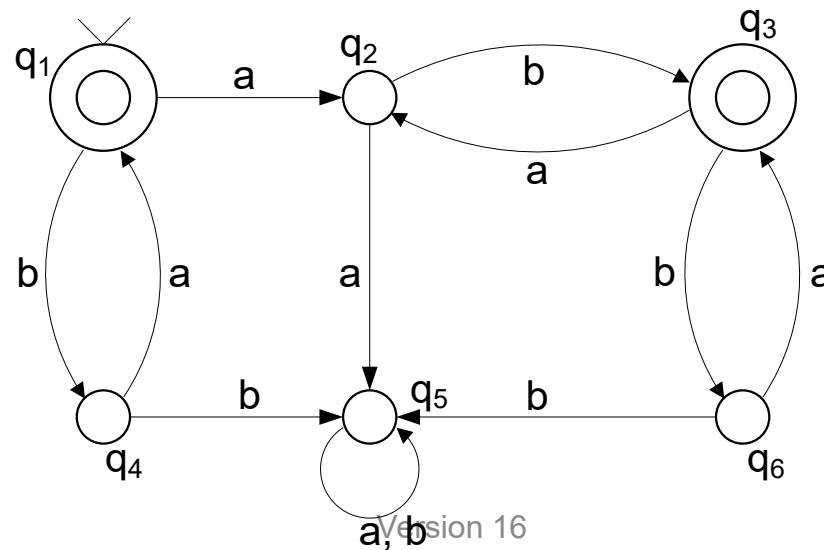
- E.g.:
  - $ba \sim_M baba$ 
    - $(s, ba) |-_M^* (q_1, e)$
    - $(s, baba) |-_M^* (q_1, e)$
  - $a \sim_M baababa$  (we do not say that  $x, y \in L$ )
    - $(s, a) |-_M^* (q_2, e)$
    - $(s, baababa) |-_M^* (q_2, e)$

# Equivalent strings

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $\sim_M$  is an equivalence relation, so it defines equivalence classes
  - $\sim_M$  the relation holds between any two elements of a class
  - e.g.:  $e \sim_M ba \sim_M baba \sim_M \dots$
  - there is  $|K|$  equivalence class denoted by:  $E_{q_1}, E_{q_2}, \dots$

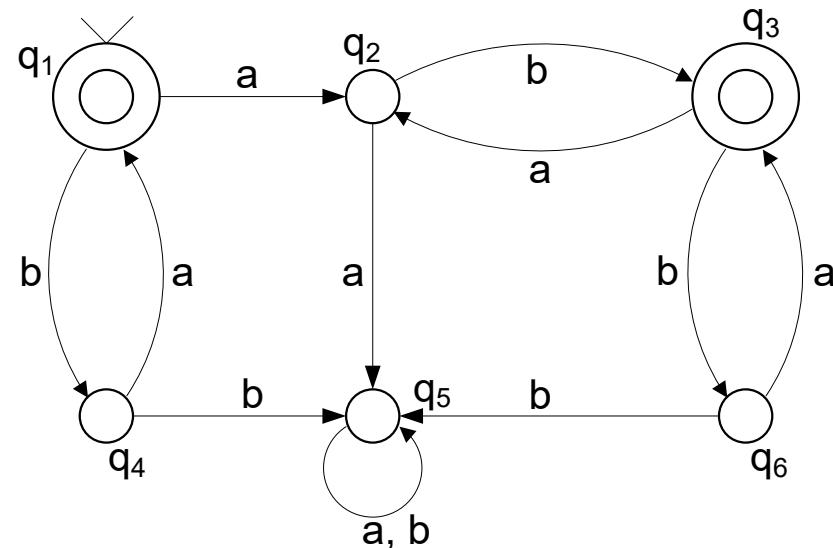


# Equivalent strings

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The equivalence classes of  $\sim_M$ :
  - $E_{q_1} = (ba)^*$
  - $E_{q_2} = La$
  - $E_{q_3} = (ab)^*abL$ 
    - the elements of the language preceded by ab
  - $E_{q_4} = b(ab)^*$
  - $E_{q_5} = L(bb \cup aa)\Sigma^*$
  - $E_{q_6} = (ab)^*abLb$



- Theorem: for any DFA  $M = (K, \Sigma, \delta, s, F)$ ,  $x, y \in \Sigma^*$ ,  
if  $x \sim_M y \rightarrow x \approx_{L(M)} y$  (the left side is more strict)
- Proof:
  - informal: if  $x$  and  $y$  drives  $M$  to the same  $q \rightarrow xz$  and  $yz$  drives  $M$  to the same  $q \rightarrow$  both  $xz$  and  $yz$  are accepted or rejected
  - let  $q(x) \in K$  such that  $(s, x) \vdash_M^* (q(x), e)$
  - $(s, x) \vdash_M^* (q(x), e) \leftrightarrow (s, xz) \vdash_M^* (q(x), z)$
  - let  $z \in \Sigma^*$  such  $xz \in L(M) \leftrightarrow (s, xz) \vdash_M^* (f, e)$ ,  $f \in F$
  - $xz \in L(M) \leftrightarrow (s, xz) \vdash_M^* (q(x), z) \vdash_M^* (f, e)$ ,  $f \in F$ ,  
from the previous two
  - $q(x) = q(y)$ , by the definition of  $x \sim_M y$
  - $(s, yz) \vdash_M^* (q(y), z) \vdash_M^* (f, e)$ ,  $f \in F$  (previous two)
  - $(s, yz) \vdash_M^* (f, e)$ ,  $f \in F \leftrightarrow yz \in L(M)$
  - both  $xz, yz \in L(M)$  which is the definition of  $x \approx_{L(M)} y$

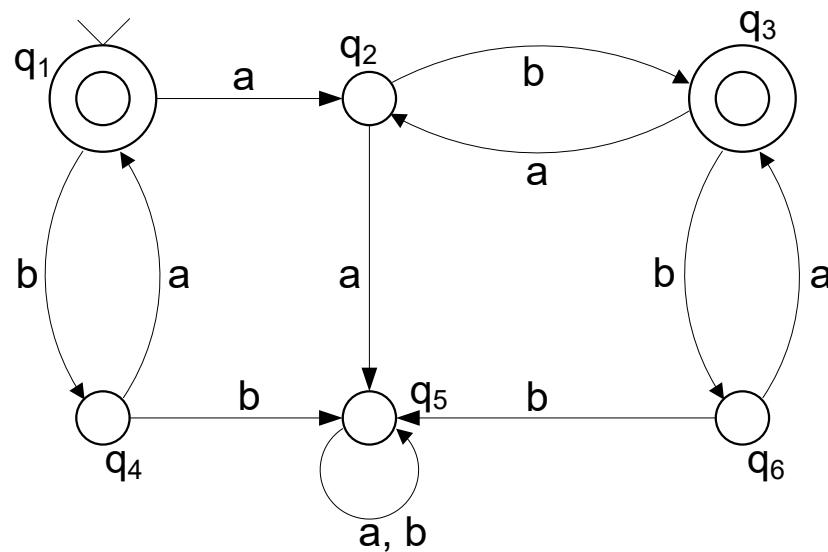
- The equivalence relation  $\sim$  is a refinement of  $\approx$ :  
 $x \sim y \rightarrow x \approx y$ 
  - each equivalence class of  $\approx$  is the union of one or more equivalence classes of  $\sim$

# Equivalent strings

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
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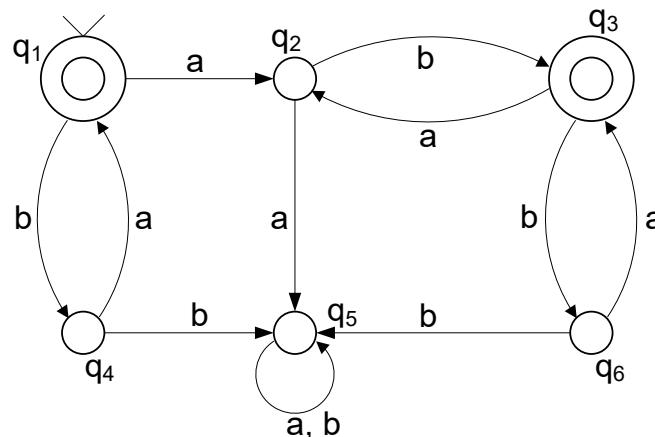
- E.g.:
  - $\sim_M$  is a refinement of  $\approx_{L(M)}$ 
    - $[e] = E_{q_1} \cup E_{q_3}$
    - $[a] = E_{q_2}$
    - $[b] = E_{q_4} \cup E_{q_6}$
    - $[aa] = E_{q_5}$



- Theorem (Myhill-Nerode theorem): if  $L \subseteq \Sigma^*$  is a regular language,  $L = L(M)$  where  $M$  is a DFA  $\rightarrow |K| =$  the number of equivalence classes in  $\approx_L$ 
  - $K$  is the set of state of  $M$
  - for regular languages always exists  $M$  with minimal states

- Definition of set  $A_M$ : let  $M = (K, \Sigma, \delta, s, F)$  be a DFA,  
 $(q, w) \in A_M \leftrightarrow (q, w) \vdash_M^* (f, e), f \in F$   
–  $w$  drives  $M$  from  $q$  to a final state

- Definition of equivalent states of DFA M,  $q \equiv p$ :  
if  $(\forall z \in \Sigma^*, (q, z) \in A_M \leftrightarrow (p, z) \in A_M) \rightarrow q \equiv p$ 
  - if each  $z$  which drives M from  $q$  to a final state also drives M from  $p$  to a final state  $\rightarrow q \equiv p$ 
    - also true for non-final states
  - if each string drives M from  $q$  and  $p$  to the same state (in terms of acceptance)  $\rightarrow q \equiv p$
  - the problem again is we cannot check all  $z$

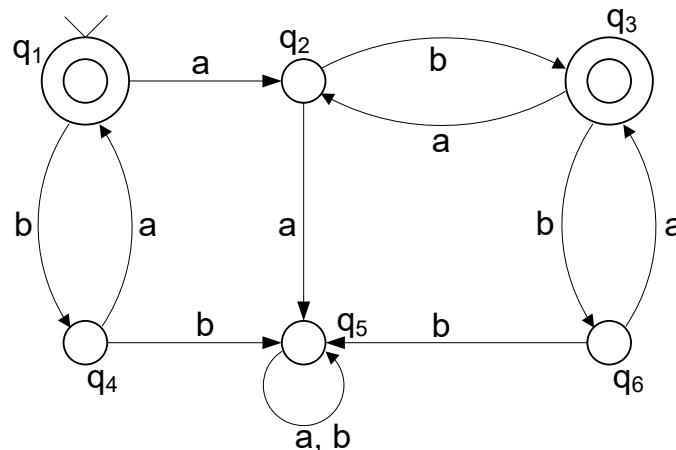


# Merging equivalent states

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- If  $q \equiv p \rightarrow E_q, E_p$  are subsets of the same equivalence class of  $\approx_L$ 
  - e.g.:  $q_1 \equiv q_3 \rightarrow E_{q1}, E_{q3} \subseteq [e]$
- Our strategy is not to determine relation  $\equiv$  directly but regard each state equivalent and check if some refinements can be done

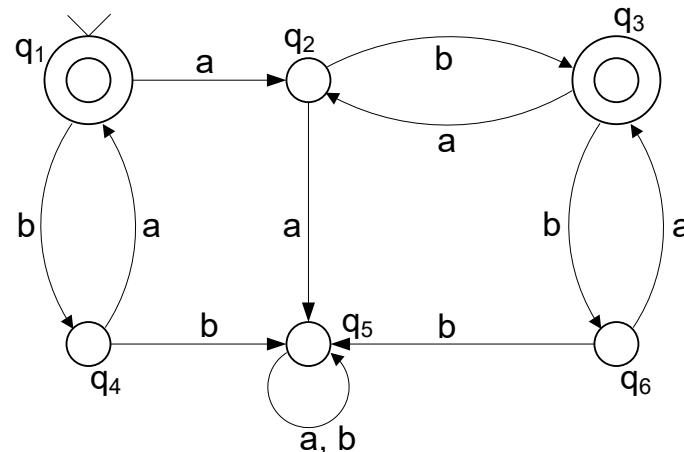


# Merging equivalent states

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of relation  $\Xi_n$ : if  $(\forall z \in \Sigma^*, |z| \leq n, (q, z) \in A_M \leftrightarrow (p, z) \in A_M) \rightarrow q \Xi_n p$ 
  - if each  $z$  of length up to  $n$  which drives  $M$  from  $q$  to a final state also drives  $M$  from  $p$  to a final state  $\rightarrow q \Xi_n p$
  - if  $n$  long strings drive  $M$  from  $q$  and  $p$  to the same state (in terms of acceptance)  $\rightarrow q \Xi_n p$



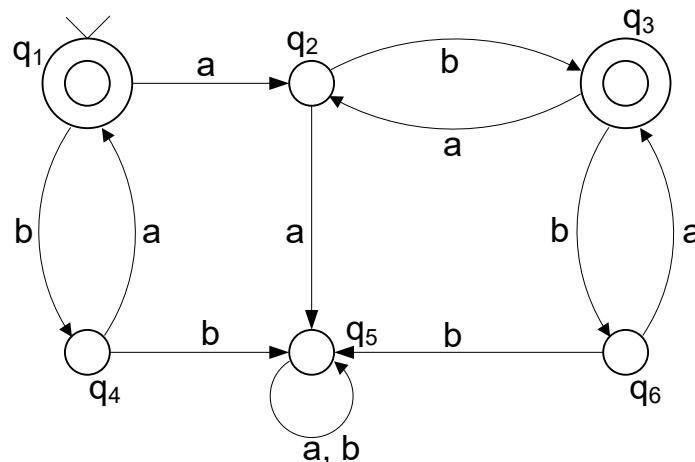
- $q \equiv_0 p$ 
  - $z = e$
  - if  $e$  drives  $M$  from  $q$  to a final state also drives  $M$  from  $p$  to a final state  $\rightarrow q \equiv_0 p$
  - if  $q, p \in F$  or  $q, p \notin F \rightarrow q \equiv_0 p$ 
    - $e$  means 0 yield step
    - $q \equiv_0 p$  holds if both states are either final or both are non-final
    - this can be checked easily
  - $\equiv_0$  has two equivalence classes:  $F, K-F$

# Merging equivalent states

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem:  $q, p \in K, n \geq 1, q \equiv_n p \leftrightarrow q \equiv_{n-1} p$  and  
 $\forall \sigma \in \Sigma, \delta(q, \sigma) \equiv_{n-1} \delta(p, \sigma)$

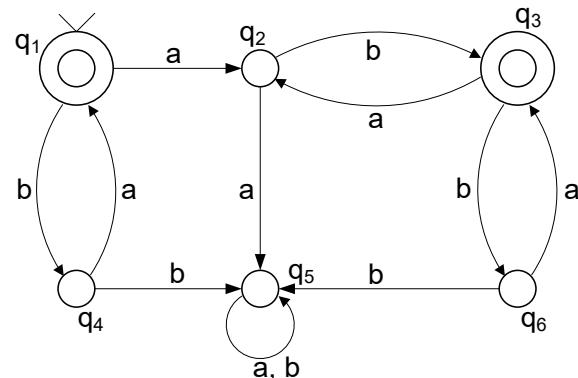


# Merging equivalent states

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Proof:
  - $q \equiv_n p \rightarrow q \equiv_{n-1} p$ 
    - if  $n$  long strings drive  $M$  to the same state (in terms of acceptance) then  $n-1$  long strings do the same
  - $w = \sigma v$ ,  $|w| = n$
  - left side:  $\forall$   $n$  long,  $w = \sigma v$ , strings drives  $q$  and  $p$  the same way (in terms of acceptance)
  - right side:  $\forall$   $n-1$  long  $v$  strings drives states  $\delta(q, \sigma)$  and  $\delta(p, \sigma)$  the same way for  $\forall \sigma$



- If  $\forall \sigma \in \Sigma$  drives M from q and p to same equivalence classes of  $\Xi_{n-1} \rightarrow q \equiv_n p$ 
  - if the equivalence classes of  $\Xi_{n-1}$  are known, we can check if  $\delta(q, \sigma) \equiv_{n-1} \delta(p, \sigma), \forall \sigma \in \Sigma$
- The equivalence classes of  $\Xi_0$  is known

- If two elements of an equivalence class of  $\Xi_{n-1}$  are not equivalent according to  $\Xi_n$  then the class has to be refined (split)
- Each equivalence relation in  $\Xi_0, \Xi_1, \Xi_2, \dots$  is a refinement of the previous one

initially the equivalence classes of  $\equiv_0$   
are F and K-F

repeat,  $n = 1, 2, \dots$

    compute equivalence classes of  $\equiv_n$   
    from  $\equiv_{n-1}$

until  $\equiv_n$  and  $\equiv_{n-1}$  are the same

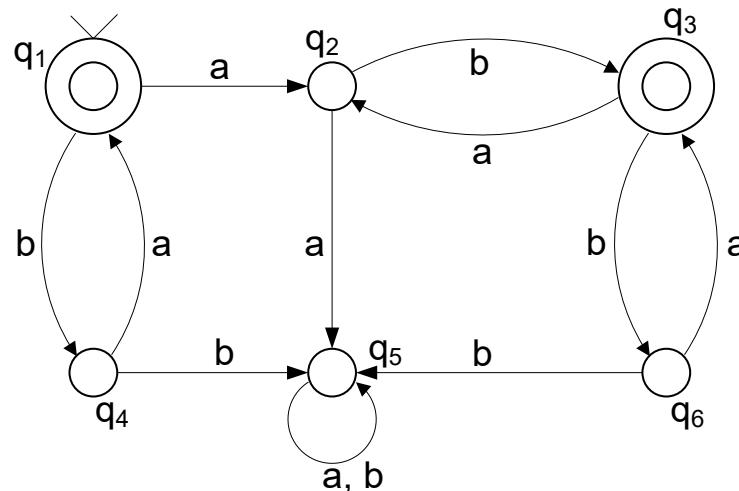
- The algorithm is finite as the maximum number of iterations is  $|K|-1$

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The equivalence classes of  $\Xi_0$  are:  $\{q_1, q_3\} = F$ ,  
 $\{q_2, q_4, q_5, q_6\} = K - F$ 
  - denote a class with its smallest element
    - e.g.:  $\{q_1, q_3\} = Q_1$ ,  $\{q_2, q_4, q_5, q_6\} = Q_2$
  - the sets are disjoint, so the identifiers are unique

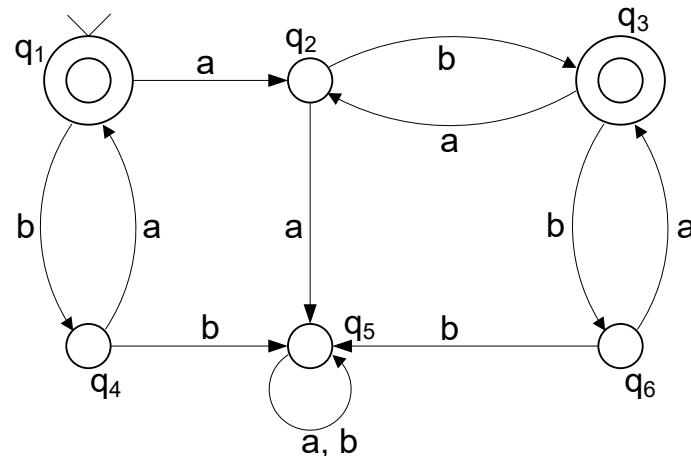


# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Determine the equivalence classes of  $\Xi_1$ 
  - do we split  $\{q_1, q_3\}$ ?
    - $\delta(q_1, a) \equiv_0 \delta(q_3, a)$ ,  $q_2 \equiv_0 q_2$ ,  $Q_2 = Q_2$
    - $\delta(q_1, b) \equiv_0 \delta(q_3, b)$ ,  $q_4 \equiv_0 q_6$ ,  $Q_2 = Q_2$
    - 'a' and 'b' drive M to the same equivalence class of  $\Xi_0$  so we do not split  $\{q_1, q_3\}$

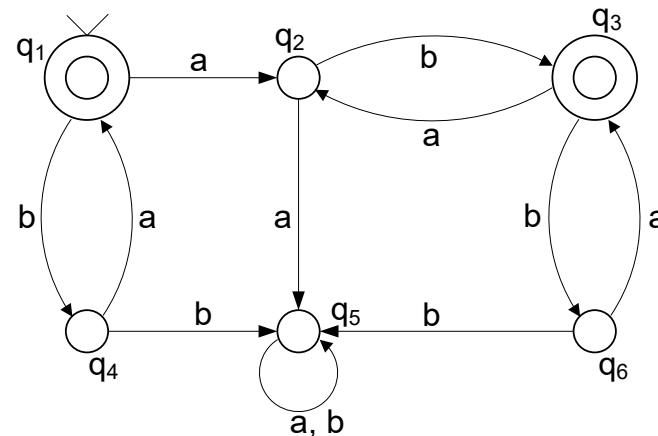


# Example

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
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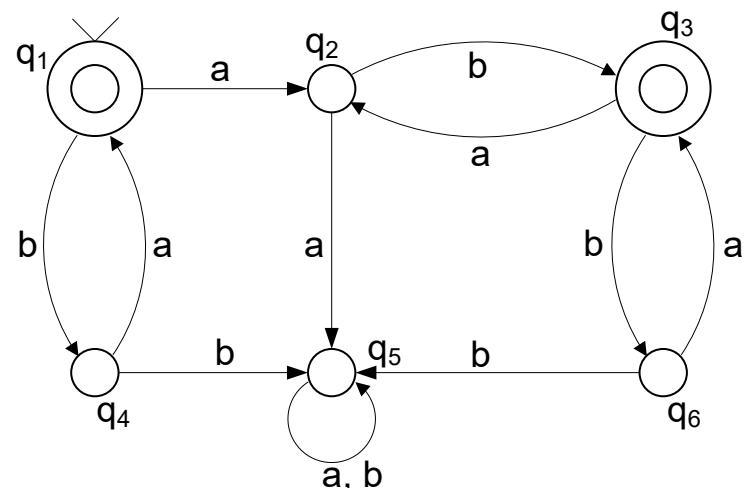
- if  $\delta(q_j, \sigma_1) = Q_{j1}, \delta(q_j, \sigma_2) = Q_{j2}, \dots \rightarrow$  let us denote  
 $\delta(q_j, \Sigma) = (Q_{j1}, Q_{j2}, \dots)$
- do we split  $\{q_1, q_3\}$ ?
  - $\delta(q_1, \Sigma) = (Q_2, Q_2)$
  - $\delta(q_3, \Sigma) = (Q_2, Q_2)$ 
    - the elements within a column are the same  $\rightarrow$  no split



– do we split  $\{q_2, q_4, q_5, q_6\}$  ?

- $\delta(q_2, \Sigma) = (Q_2, Q_1)$
  - $\delta(q_4, \Sigma) = (Q_1, Q_2)$
  - $\delta(q_5, \Sigma) = (Q_2, Q_2)$
  - $\delta(q_6, \Sigma) = (Q_1, Q_2)$
  - the elements within a column are not the same  
→ split
  - $\delta(q_2, \Sigma) = (Q_2, Q_1)$
- 
- $\delta(q_4, \Sigma) = (Q_1, Q_2)$
  - $\delta(q_6, \Sigma) = (Q_1, Q_2)$
- 
- $\delta(q_5, \Sigma) = (Q_2, Q_2)$

$$\begin{aligned}\{q_1, q_3\} &= Q_1 \\ \{q_2, q_4, q_5, q_6\} &= Q_2\end{aligned}$$



– the equivalence classes of  $\Xi_1$  are:  $\{q_1, q_3\}, \{q_2\}, \{q_4, q_6\}, \{q_5\}$

# Example

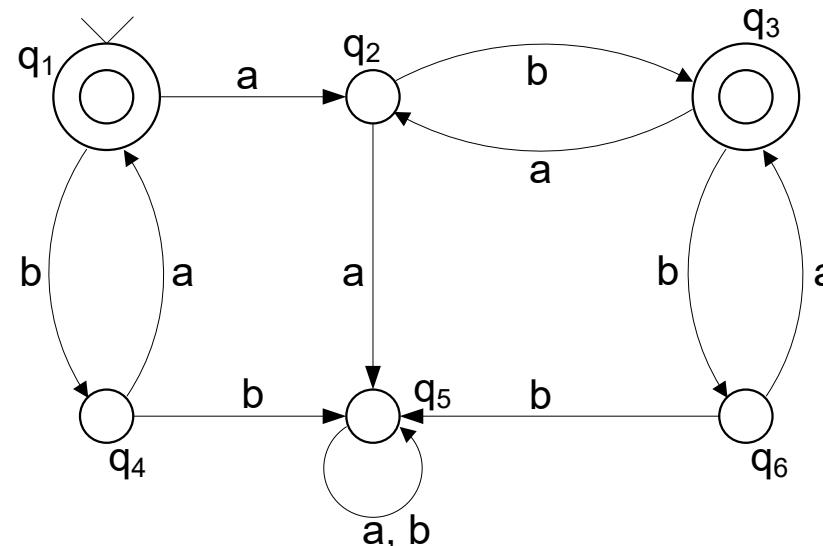
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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

– do we split  $\{q_4, q_6\}$  ?

- singleton set like  $\{q_2\}$  can not be split
- $\delta(q_4, \Sigma) = (Q_1, Q_5)$
- $\delta(q_6, \Sigma) = (Q_1, Q_5)$
- the elements within a column are the same → no split

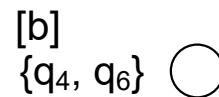
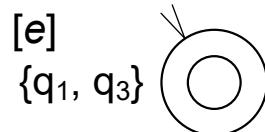
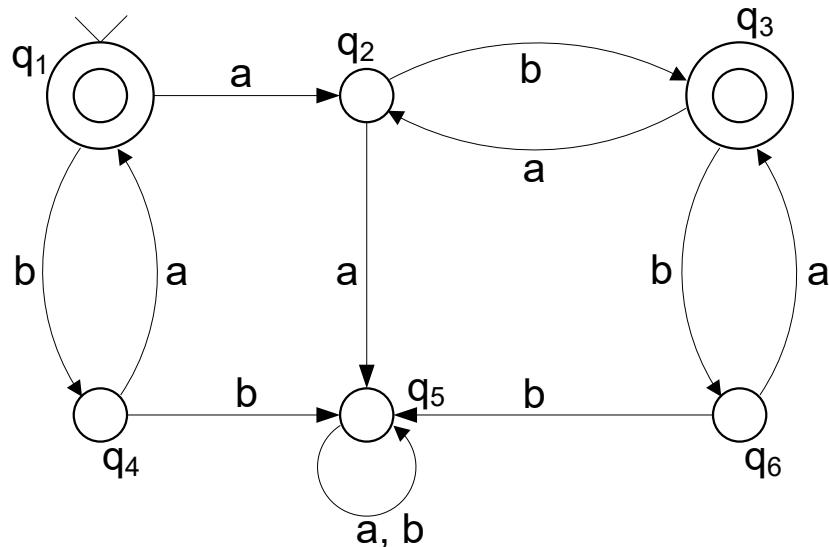
$$\begin{aligned}\{q_1, q_3\} &= Q_1 \\ \{q_2\} &= Q_2 \\ \{q_4, q_6\} &= Q_4 \\ \{q_5\} &= Q_5\end{aligned}$$



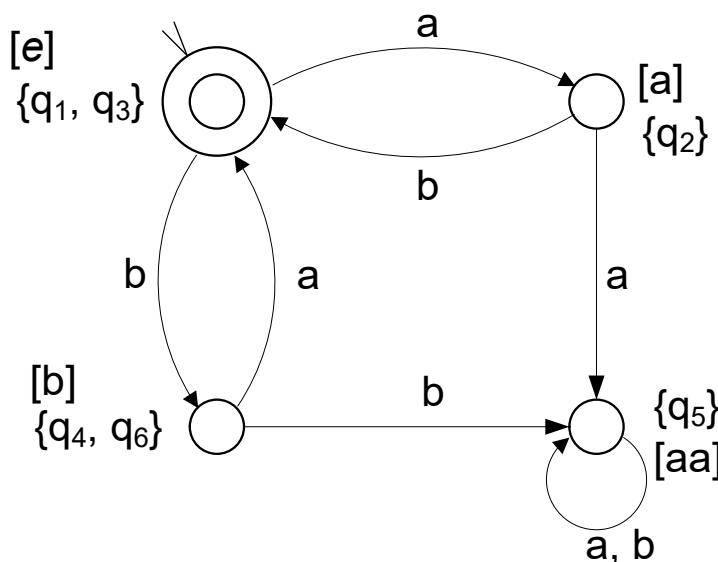
# Merging equivalent states

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen



- Draw the nodes
- Draw the arcs





# Introduction to the Theory of Computation

## Lesson 2

2.6. Algorithmic aspects of finite automat

3.2. Parse trees

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- Theorem: there are algorithms to accomplish the next tasks with the given complexities
  - $\text{NFA} \rightarrow \text{DFA}$  exponential
  - $\text{RE} \rightarrow \text{NFA}$  polynomial
  - $\text{NFA} \rightarrow \text{RE}$  exponential
  - $\text{DFA} \rightarrow \text{minimal DFA}$  polynomial

- Theorem: regard RE  $r$  and NFA  $M$  such that  $L(M) = L(r)$   
 $\rightarrow |K| \sim |r|$
- Proof:
  - the basic machines which accept a simple symbol have two states
  - union, Kleene star introduce 1 new state
  - $|K| < 2^*|r|$  according to the Thomson's construction theorems
- There is at most 2 transitions from each state,  $|\delta| \sim |r|$ 
  - it is not true for NFA in general but only to those which were created by the construction theorems

- It is possible to simulate an NFA by constructing a DFA on-the-fly
  - $S$  is the current DFA state, i.e., a set of NFA states

# Space-time tradeoff

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

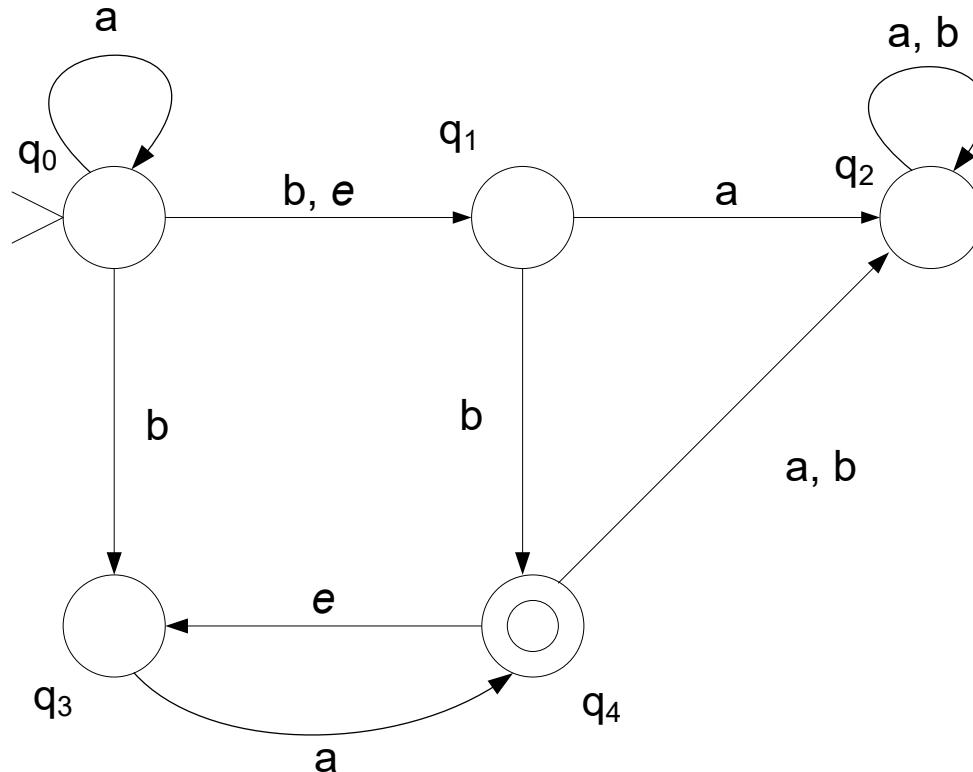
```
S0 = E(s)
a = nextchar()
while a ≠ eos
    n = n + 1
    Sn = ∪{E(q), ∀ p ∈ Sn-1, (p, a, q) ∈
    Δ}
    a = nextchar()
end
if Sn-1 ∩ F = ∅ return no
else return yes
```

# Example

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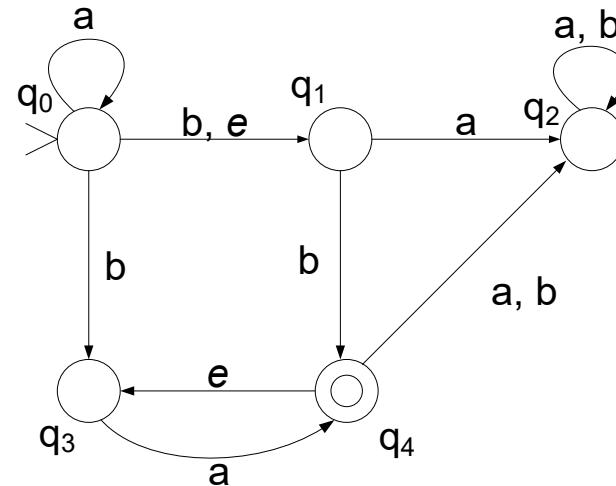
A felsőfokú oktatás minőségének és hozzáférhetőségének  
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- Check directly if aaaba is accepted by the NFA below!



# Example

- The various values for the set  $S_n$  are shown below
  - $S_0 = \{q_0, q_1\}$
  - $S_1 = \{q_0, q_1, q_2\}$
  - $S_2 = \{q_0, q_1, q_2\}$
  - $S_3 = \{q_0, q_1, q_2\}$
  - $S_4 = \{q_1, q_2, q_3, q_4\}$
  - $S_5 = \{q_2, q_3, q_4\}$ 
    - it contains a final state so the word is accepted



- Definition of the size of a finite automaton,  $|\delta|$ : the size of the transition table
- Definition of the time (speed) of a finite automaton: the number of step required to process  $w \in \Sigma^*$

- Regard RE  $r$ ,  $w \in \Sigma^*$ , decide if  $w \in L(r)$ 
  - method 1:
    - construct NFA  $M_1$  such that  $L(M_1) = L(r)$
    - construct DFA  $M_2$  such that  $L(M_2) = L(M_1)$
    - space:  $|\delta_2| \sim 2^{|\Delta_1|} \sim 2^{|r|}$ 
      - if  $K$  increases  $\rightarrow \delta$  also increases
    - time:  $|w|$ 
      - one symbol is processed in every step

## – method 2:

- construct NFA  $M_1$  such that  $L(M_1) = L(r)$
- simulate  $M_1$  directly
- space:  $|\delta_1| \sim |r|$
- time:  $|w|^*|r|^2 \sim |w|^*|K_1|^2$ 
  - $|w|$  iteration of the while loop
  - at most  $|K_1|^2$  steps is required to calculate the new  $S$  set
    - »  $|K_1|^2$  triplet,  $(p, a, q) \in \Delta$ , is to be checked

# Space-time tradeoff

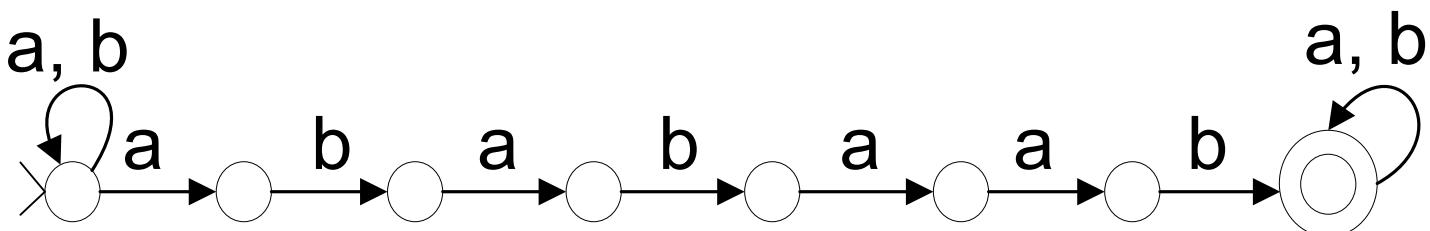
EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

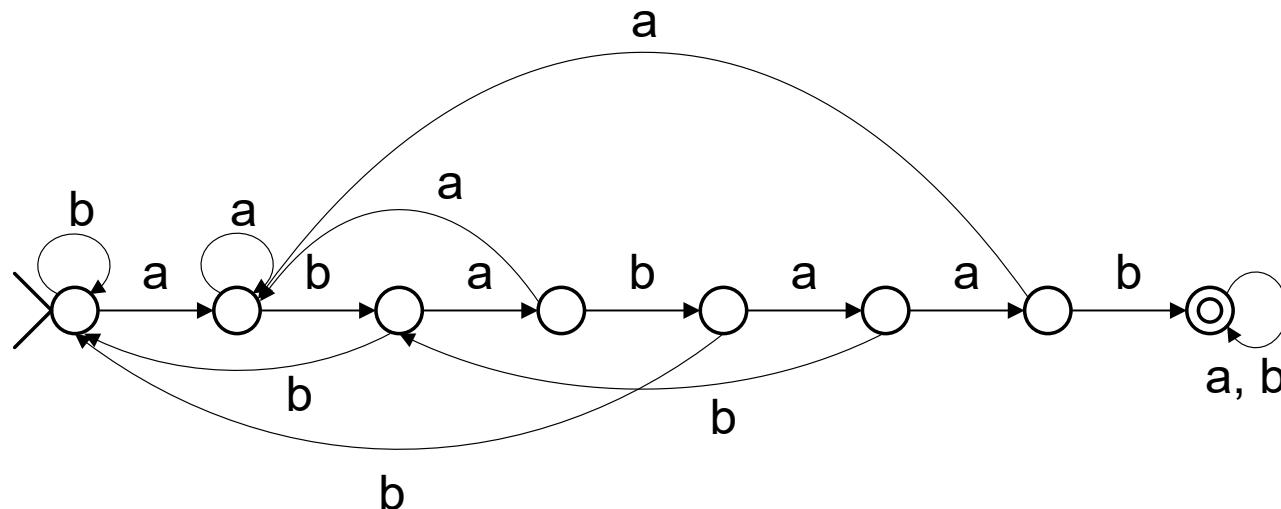
	Method 1	Method 2
space	$O(2^{ r })$	$O( r )$
time	$O( w )$	$O( w ^* r ^2)$

- There are combined techniques which is better for certain type of problems

- Given a fixed string  $x$ , determine if  $x$  is a substring of  $w$ !
  - remember that  $x$  is fixed and not an input
- The NFA below recognize pattern ababaab
  - $\Sigma = \{a, b\}$



- Though the NFA → DFA conversion may increases the number of states exponentially at this example this does not happen
  - the size of the DFA does not increase exponentially
  - in practice, another kind of algorithm is used because  $\Sigma$  is usually large



- Definition of:
  - regular grammar
  - context-free grammar
  - derivation of a string
  - partially defined string

- Give grammar  $G$  which can create a sequence of balanced parentheses!
  - every right parenthesis can be paired with a unique preceding left parenthesis
  - this is not a regular language
- Let  $G$  be the grammar  $(V, \Sigma, R, S)$ , where
  - $V = \{S, (, )\}$
  - $\Sigma = \{(, )\}$
  - $R = \{S \rightarrow e \mid SS \mid (S)\}$ 
    - it is in CFG but not in RG

- A string  $w \in L(G)$  may have many derivations in  $G$
- E.g.: if  $G$  is the CFG that generates the language of balanced parentheses, then the string  $()()$  can be derived from  $S$  by at least two distinct derivations
  - $S \Rightarrow SS \Rightarrow (S)S \Rightarrow ()S \Rightarrow ()(S) \Rightarrow ()()$
  - $S \Rightarrow SS \Rightarrow S(S) \Rightarrow (S)(S) \Rightarrow (S)() \Rightarrow ()()$

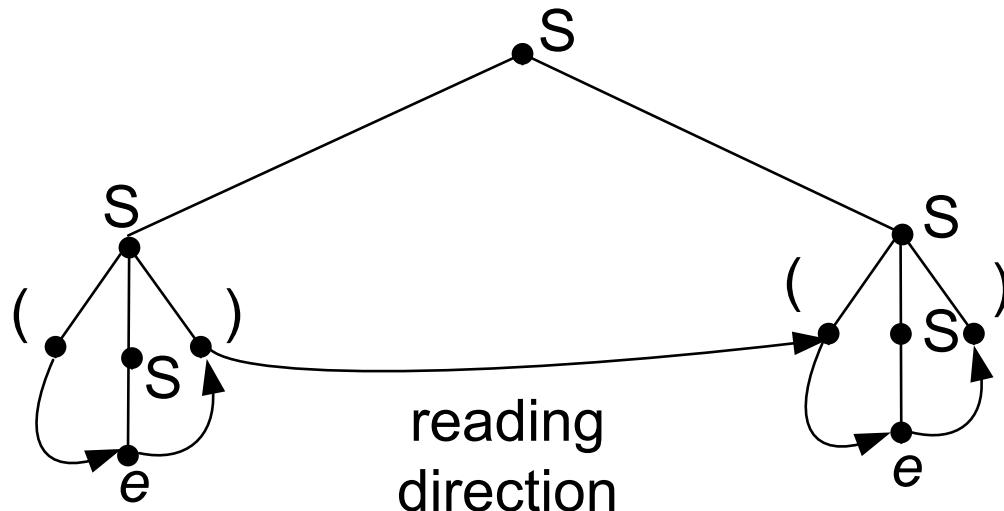
- These two derivations are "the same" in a sense
  - the rules used are the same
  - the rules are applied for the same non-terminals in the partially defined string
  - the only difference is in the order in which the rules are applied

# Parse trees

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Both derivations can be pictured like this



- Such a picture is called a parse tree
  - the points are called nodes
  - each node carries a label that is a symbol in V
  - the topmost node is called the root
  - the nodes along the bottom are called leaves
  - all leaves are labeled by terminals or possibly the empty string e
- Definition of the yield of the parse tree: concatenate the labels of the leaves from left to right
  - the derived string

- Definition of the parse tree for CFG  $G(V, \Sigma, R, S)$ :
  - a
    - this is the parse tree for each  $a \in \Sigma$ 
      - its root: 'a'
      - its leaf: 'a'
      - its yield: 'a'

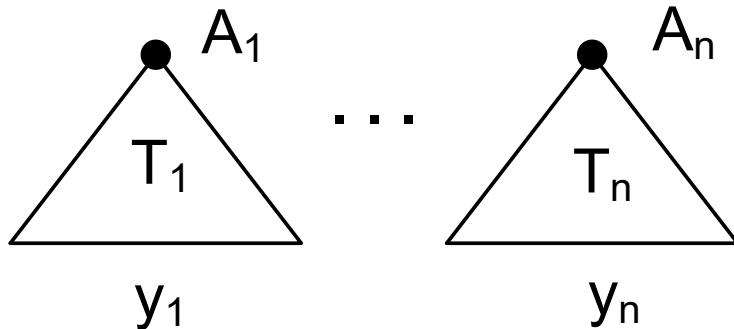
# Parse trees

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen



- this is a parse tree if  $A \rightarrow e$  is a rule in R
  - its root: A
  - its leaf: e
  - its yield: e

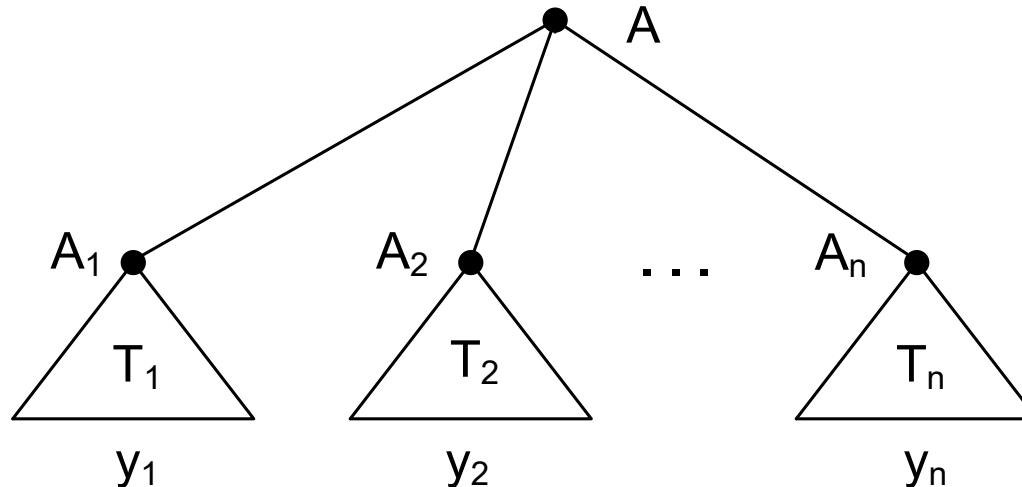


- if the above graphs are parse trees, where  $n \geq 1$ , with roots labeled  $A_1, A_2, \dots, A_n$  respectively, and with yields  $y_1 \dots y_n$ , and  $A \rightarrow A_1 A_2 \dots A_n$  is a rule in  $R$  then:

# Parse trees

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen



- this is a parse tree
  - its root: A
  - its leaves: the leaves of its constituent parse trees
  - its yield:  $y_1 y_2 \dots y_n$
- nothing else is a parse tree

- Parse trees are ways of representing derivations of strings in  $L(G)$ 
  - the superficial differences between derivations, owing to the order of application of rules, are suppressed

- Parse trees represent equivalence classes of derivations
  - several derivation may belong to the same parse tree
  - there is a precedence relation between the derivations belonging to the same parse tree

- We say a derivation precedes another if
  - the two derivations are identical except for two consecutive steps
  - in these steps the two nonterminals are replaced by the same two strings but in opposite order in the two derivations
  - derivation A precedes derivation B if in derivation A the leftmost of the two mentioned nonterminals is replaced first

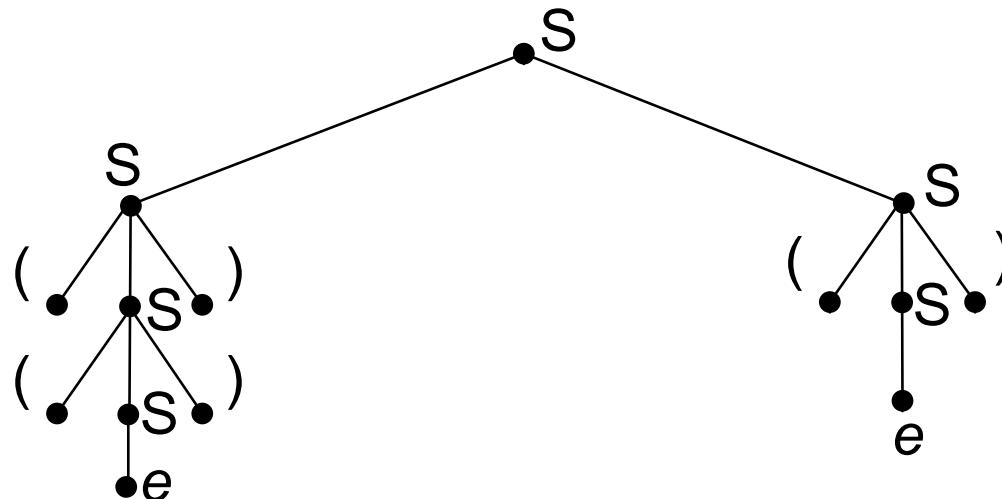
- Definition of precedence relation,  $\leq$ :
  - let  $G(V, \Sigma, R, S)$  is a CFG
  - $D=x_1 \Rightarrow x_2 \Rightarrow \dots \Rightarrow x_n$ ,  $D'=x'_1 \Rightarrow x'_2 \Rightarrow \dots \Rightarrow x'_n$  are derivations in  $G$ 
    - where  $x_i, x'_i \in V^*$  for  $i = 1, \dots, n$ ,  $x_1, x'_1 \in V - \Sigma$ ,  
 $x_n = x'_n \in \Sigma^*$
  - $D$  and  $D'$  are derivations of terminal strings from a single nonterminal

– we say that D precedes D', written  $D \leq D'$ , if  $n > 2$ ,  $\exists k$ ,  $1 < k < n$  such that

- $\forall i \neq k, x_i = x'_i$
- $x_{k-1} = x'_{k-1} = uAvBw$  (they have this form)
  - $u, v, w, y, z \in V^*$ ,  $A, B \in V - \Sigma$
- $x_k = uyvBw$ , where  $A \rightarrow y \in R$
- $x'_k = uAvzw$ , where  $B \rightarrow z \in R$
- $x_{k+1} = x'_{k+1} = uyvzw$

## Example

- Recall CFG G generating all strings of balanced parentheses
  - Consider these three derivations  $D_1, D_2, D_3$  of  $((())()$ 
    - $D_1 = S \Rightarrow SS \Rightarrow (S)S \Rightarrow ((S))S \Rightarrow ((())S \Rightarrow ((())(S) \Rightarrow ((())()$
    - $D_2 = S \Rightarrow SS \Rightarrow (S)S \Rightarrow ((S))S \Rightarrow ((S))(S) \Rightarrow ((())(S) \Rightarrow ((())()$
    - $D_3 = S \Rightarrow SS \Rightarrow (S)S \Rightarrow ((S))S \Rightarrow ((S))(S) \Rightarrow ((S))() \Rightarrow ((())()$



- $D_1 \leq D_2$ 
  - they differ for  $k=5$
  - $D_1$ :  $S \rightarrow ()$  for the 1st  $S$ ,  $S \rightarrow (S)$  for the 2nd  $S$
  - $D_2$ :  $S \rightarrow (S)$  for the 2nd  $S$ ,  $S \rightarrow ()$  for the 1st  $S$
- $D_2 \leq D_3$ 
  - they differ for  $k=6$
- neither  $D_1 \leq D_3$  nor  $D_3 \leq D_1$  is true because the two derivations differ in more than one intermediate string
  - $\leq$  is not total order

- Definition of similarity relation: the reflexive, symmetric, transitive closure of  $\leq$ 
  - it is an equivalence relation
- Two derivations are similar if
  - they can be transformed into each other via a sequence of "switching" in the order in which rules are applied
    - a "switching" transform a derivation either by one that precedes it, or by one that it precedes
  - they belong to the same parse tree

- The equivalence class of the derivations of  $((())()$  corresponding to the tree showed previously contains the derivations  $D_1, D_2, D_3$  and also
  - $D_4 = S \Rightarrow SS \Rightarrow (S)S \Rightarrow (S)(S) \Rightarrow ((S))(S) \Rightarrow ((())(S) \Rightarrow ((())()$
  - $D_5 = S \Rightarrow SS \Rightarrow (S)S \Rightarrow (S)(S) \Rightarrow ((S))(S) \Rightarrow ((S))() \Rightarrow ((())()$
  - $D_6 = S \Rightarrow SS \Rightarrow (S)S \Rightarrow (S)(S) \Rightarrow (S)() \Rightarrow ((S))() \Rightarrow ((())()$
  - $D_7 = S \Rightarrow SS \Rightarrow S(S) \Rightarrow (S)(S) \Rightarrow ((S))(S) \Rightarrow ((())(S) \Rightarrow ((())()$
  - $D_8 = S \Rightarrow SS \Rightarrow S(S) \Rightarrow (S)(S) \Rightarrow ((S))(S) \Rightarrow ((S))() \Rightarrow ((())()$
  - $D_9 = S \Rightarrow SS \Rightarrow S(S) \Rightarrow (S)(S) \Rightarrow (S)() \Rightarrow ((S))() \Rightarrow ((())()$
  - $D_{10} = S \Rightarrow SS \Rightarrow S(S) \Rightarrow S() \Rightarrow (S)() \Rightarrow ((S))() \Rightarrow ((())()$
- These derivations are all similar

# Example

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együttes javítása a Pannon Egyetemen

- These ten derivations are related by  $\leq$  as shown here

$$\begin{array}{ccccccccc} D_1 & \leq & D_2 & \leq & D_3 & \leq & D_5 & \leq & D_6 \\ & & & & \swarrow & & & & \searrow \\ & & & & D_4 & & & & D_9 \\ & & & & \leq & D_7 & \leq & D_8 & \leq \end{array} \quad D_9 \leq D_{10}$$

- Definition of the leftmost derivation: each equivalence class of derivations under similarity relation, i.e., each parse tree, contains a derivation that is not preceded by any other derivation
- Definition of one step leftmost derivation,  $\Rightarrow^L$ :
  - if  $x = wA\beta$ ,  $y = wa\beta$ ,  $w \in \Sigma^*$ ,  $a, \beta \in V^*$ ,  $A \in V - \Sigma$ ,  
 $A \rightarrow a \in R$
  - then  $x \Rightarrow^L y$
- A leftmost derivation exists in every parse tree
  - repeatedly replace the leftmost nonterminal in the current string
  - a leftmost derivation contains only  $\Rightarrow^L$  derivation

- Definition of the rightmost derivation: each equivalence class of derivations under similarity relation, i.e. each parse tree, contains a derivation that does not preceded any other derivation
- Definition of one step rightmost derivation  $=>^R$ :
  - if  $x = \beta Aw$ ,  $y = \beta\alpha w$ ,  $w \in \Sigma^*$ ,  $\alpha, \beta \in V^*$ ,  $A \in V - \Sigma$ ,  
 $A \rightarrow \alpha \in R$
  - then  $x =>^R y$
- A rightmost derivation exists in every parse tree
  - repeatedly replace the rightmost nonterminal in the current string
  - a rightmost derivation contains only  $=>^R$  derivation

- The leftmost derivation is unique since at each step the leftmost nonterminal is to be replaced
- The rightmost derivation is unique since at each step the rightmost nonterminal is to be replaced
- In the previous example  $D_1$  is the leftmost derivation and  $D_{10}$  is the rightmost one

- Theorem:
  - $G(V, \Sigma, R, S)$  is a CFG,  $A \in V - \Sigma$ ,  $w \in \Sigma^*$
  - the following statements are equivalent
    - $A \Rightarrow^* w$
    - there is a parse tree with root  $A$  and yield  $w$
    - there is a leftmost derivation  $A \Rightarrow^{L*} w$
    - there is a rightmost derivation  $A \Rightarrow^{R*} w$

- There may be a string in a language generated by a context-free grammar with two derivations that are not similar
  - the string has distinct parse trees (or equivalently distinct rightmost derivations, and two distinct leftmost derivations)

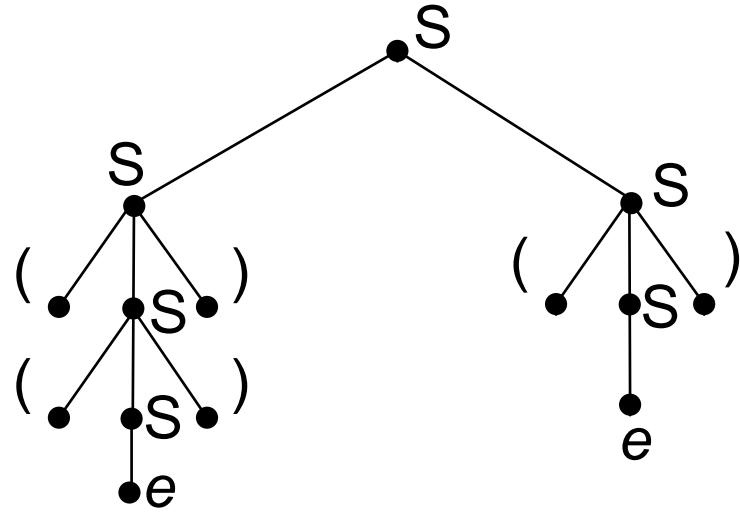
- Recall CFG G generating all strings of balanced parentheses
- There are other derivations of  $((())()$  that are not similar to the ones already introduced ( $D_1 - D_{10}$ )
- They do not belong to the parse tree shown previously
- E.g.:
  - $S \Rightarrow SS \Rightarrow SSS \Rightarrow S(S)S \Rightarrow S((S))S \Rightarrow S((())S \Rightarrow S((())(S) \Rightarrow S((())() \Rightarrow ((())()$

# Example

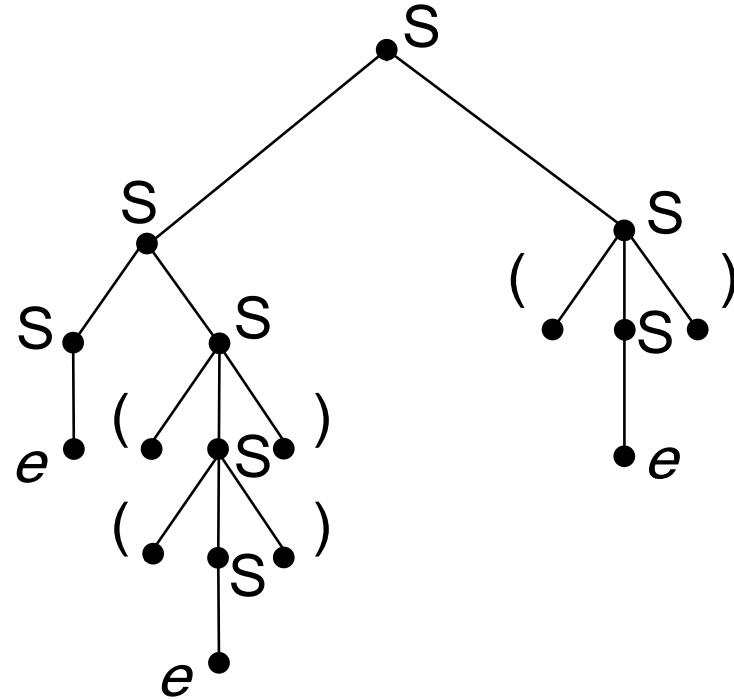
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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

old



new



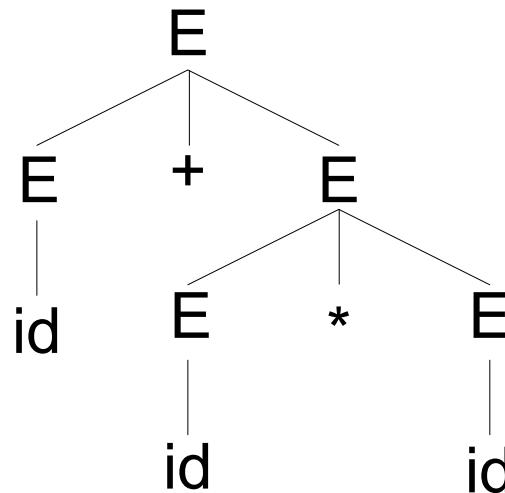
- Consider  $G'(V, \Sigma, R, E)$ , where
  - $V = \{+, *, (, ), id, E\}$
  - $\Sigma = \{+, *\}, (, ), id\}$
  - $R = \{E \rightarrow E+E \mid E^*E \mid (E) \mid id\}$
- $G'$  generates all arithmetic expressions over  $id$ 
  - recall that we already generated this language with grammar  $G$  which contained nonterminals for term and factor
  - $L(G') = L(G)$ , but  $G'$  does not express that addition and multiplication has different precedence

# Example

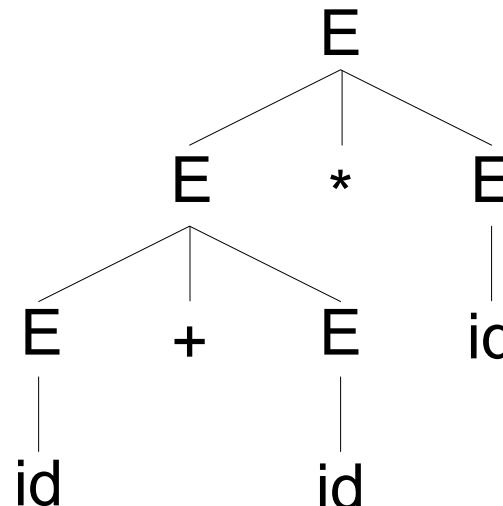
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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- There are two parse trees for  $\text{id} + \text{id}^* \text{id}$  in  $G'$
- The first parse tree corresponds to the natural meaning of this expression (with  $*$  taking precedence over  $+$ ), the other is "wrong"



(a)



(b)

- Definition of ambiguous grammar: such a CFG whose language contains such a string that have two or more distinct parse trees
  - e.g.:  $G'$
- Definition of parsing: assigning a parse tree to a given string
- Unambiguous grammars are important for programming languages

- We can "disambiguate" some ambiguous grammar
  - e.g. for  $G'$  we can introduce the nonterminals  $T$  and  $F$
  - e.g. there is an unambiguous grammar to create balanced strings of parentheses
- Definition of inherently ambiguous languages: there is no unambiguous CFG to generate them
  - programming languages are never inherently ambiguous



# Introduction to the Theory of Computation

## Lesson 3

Chomsky normal form, deterministic context-free languages

3.6. Algorithms for context-free grammars

3.7. Determinism

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Version 16



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Európai Unió  
Európai Strukturális  
és Beruházási Alapok



BEFEKTETÉS A JÖVŐBE

Case		English Name
Upper	Lower	
A	$\alpha$	alpha
B	$\beta$	beta
$\Gamma$	$\gamma$	gamma
$\Delta$	$\delta$	delta
E	$\epsilon$	epsilon
Z	$\zeta$	zeta
H	$\eta$	eta
$\Theta$	$\theta$	theta
I	$\iota$	iota
K	$\kappa$	kappa
$\Lambda$	$\lambda$	lambda
M	$\mu$	mu

Case		English Name
Upper	Lower	
N	$\nu$	nu
$\Xi$	$\xi$	xi
O	$\circ$	omicron
$\Pi$	$\pi$	pi
P	$\rho$	rho
$\Sigma$	$\varsigma$	sigma
T	$\tau$	tau
$\Upsilon$	$\upsilon$	upsilon
$\Phi$	$\phi$	phi
X	$\chi$	chi
$\Psi$	$\psi$	psi
$\Omega$	$\omega$	omega

Recursively enumerable languages  
Unrestricted grammars generating a language  
Turing machines accepting a language

Recursive languages  
Unrestricted grammars computing a function  
Turing machines deciding a language  
 $\mu$ -recursive functions

Context-free languages  
Context-free grammars  
**Pushdown automata**

Regular languages  
Regular grammar  
Regular expressions  
Deterministic finite automata  
Nondeterministic finite automata

- Definition of:
  - PDA
  - simple PDA
  - PDA → simple PDA conversion
  - simple PDA → CFG conversion

- Theorem: there is a polynomial algorithm which, given a
  - context-free grammar, construct an equivalent pushdown automaton
  - pushdown automaton, construct an equivalent context-free grammar
  - a context-free grammar  $G$  and a string  $w$ , decides whether  $w \in L(G)$

# Comparing theorems about regular and context-free languages

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A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Similarity: language acceptors can be transformed to generators and vice versa
- Differences:
  - it is trivial to test if  $w \in \Sigma^*$  belongs to a given regular language
    - construct the appropriate DFA, run it with  $w$
    - the introduced PDA is nondeterministic, so it is not obvious that  $w$  is accepted or not
  - there is no algorithm to
    - test if  $L(M_1) = L(M_2)$ ,  $M_1$  and  $M_2$  are PDA

- Definition of Chomsky normal form: the rules of a CFG has the next form  $R \subseteq (V - \Sigma) \times V^2$ 
  - the right-hand side of every rule must have length two
  - grammars in this form cannot produce strings of length less than 2
    - otherwise they may produce anything

- Theorem: for any context-free grammar  $G \exists$  a context-free grammar  $G'$  in Chomsky normal form such that  $L(G') = L(G) - \{\Sigma \cup \{e\}\}$ 
  - the construction of  $G'$  can be carried out in time polynomial in the size of  $G$

- Proof by construction:
  - the rules of G may violate the constraints of Chomsky normal form
    - long rules: right-hand side has length 3 or more
    - e-rules:  $A \rightarrow e$
    - short rules:  $A \rightarrow a$  or  $A \rightarrow B$

- Step 1: eliminate  $\forall$  long rule:  $A \rightarrow B_1B_2 \dots B_n \in R$ ,  
 $B_1, B_2, \dots, B_n \in V$ ,  $n \geq 3$ 
  - replace the previous long rule with  $n-1$  new rules:  
 $A \rightarrow B_1A_1$   
 $A_1 \rightarrow B_2A_2$   
...  
 $A_{n-2} \rightarrow B_{n-1}B_n$
  - $A_1, A_2, \dots, A_{n-2}$  are new nonterminals
    - they do not appear anywhere else

- the resulting grammar is
  - equivalent to the original one
  - has rules with right-hand sides of length 2 or less

- Step 2: eliminate  $\forall \epsilon\text{-rule: } A \rightarrow \epsilon, A \in V - \Sigma$ 
  - determine the set of erasable nonterminals:  $\epsilon$   
 $\epsilon = \{A \in V - \Sigma : A \Rightarrow^* \epsilon\}$   
 $\epsilon := \emptyset$   
while  $\exists \alpha \in \epsilon^*, \exists A \rightarrow \alpha$ , and  $A \notin \epsilon$   
add  $A$  to  $\epsilon$
  - delete the  $\epsilon$ -rules

- for  $\forall$  rule:  $A \rightarrow BC$  or  $A \rightarrow CB$ ,  $B \in \epsilon$  add new rule  
 $A \rightarrow C$ 
  - before deleting the  $B \rightarrow \epsilon$ :  $A \Rightarrow BC \Rightarrow^* C$
  - now:  $A \Rightarrow^* C$
- the new grammar can not generate  $\epsilon$

– Step 3a: eliminate  $\forall$  short rule:  $A \rightarrow a$  or  $A \rightarrow B$

- determine the replacement sets:

$$D(A) = \{B \in V : A \Rightarrow^* B\}, \forall A \in V$$

–  $D(A)$  - set of single symbols that can be derived from  $A$  in the grammar

$$D(A) := \{A\}$$

while  $\exists B \in D(A), \exists B \rightarrow C, C \notin D(A)$

add  $C$  to  $D(A)$

- delete the short rules
- for  $\forall$  rule:  $A \rightarrow BC$  add new rules  $A \rightarrow B'C'$ ,  
 $B' \in D(B)$ ,  $C' \in D(C)$ 
  - $|D(B)| * |D(C)|$  new rules

- before deleting  $B \rightarrow B'$ :  $A \Rightarrow BC \Rightarrow^* B'C$
- now:  $A \Rightarrow^* B'C$
- Step 3b: for  $\forall$  rule  $A \rightarrow BC$  where  $A \in D(S) - \{S\}$   
add rules  $S \rightarrow BC$ 
  - before deleting  $S \rightarrow A$ :  $S \Rightarrow^* A \Rightarrow BC$
  - now:  $S \Rightarrow^* BC$
- the new grammar can not generate strings of length 1 but it produces the same language as the previous one

- Transform the next grammar into Chomsky normal form!
  - the CFG generates the set of balanced parenthesis
    - $R = \{S \rightarrow SS, S \rightarrow (S), S \rightarrow e\}$
  - removing long rule:  $S \rightarrow (S)$ 
    - adding new rules:  $S \rightarrow (S_1, S_1 \rightarrow S)$
    - if the long rule would have been  $S \rightarrow ((())()$ , the new rules would be:  $S \rightarrow (S_1, S_1 \rightarrow (S_2, S_2 \rightarrow )S_3, S_3 \rightarrow (S_4, S_4 \rightarrow ))$

- removing e-rules
  - $\varepsilon = \{S\}$
  - omit e-rule:  $S \rightarrow e$
  - add new rules:
    - $S \rightarrow S$  because of  $S \rightarrow SS$
    - $S_1 \rightarrow )$  because of  $S_1 \rightarrow S)$

- if the rules would have been  $S \rightarrow A$ ,  $A \rightarrow B$ ,  $B \rightarrow e$ ,  
 $C \rightarrow B$ ,  $D \rightarrow AB$ ,  $E \rightarrow AC$ ,  $F \rightarrow Aa$  then
  - $\varepsilon = \{B, A, S, C, D, E\}$

- the current set of rules:  $S \rightarrow SS$ ,  $S \rightarrow (S_1, S_1 \rightarrow S)$ ,  
 $S \rightarrow S$ ,  $S_1 \rightarrow )$
- removing short rules:
  - determine the replacement sets:
    - $D(S_1) = \{S_1, \}\}$
    - $D(A) = \{A\}, \forall A \in V - \{S_1\}$
  - omit the short rules:  $S_1 \rightarrow )$ ,  $S \rightarrow S$
  - add new rule:  $S \rightarrow ()$  because of  $S \rightarrow (S_1$
- the grammar in Chomsky form is:  
 $S \rightarrow SS$ ,  $S \rightarrow (S_1, S_1 \rightarrow S)$ ,  $S \rightarrow ()$

- If the rules would have been  $S \rightarrow A$ ,  $A \rightarrow B$ ,  $B \rightarrow C$ ,  
 $B \rightarrow E$ ,  $F \rightarrow AB$ ,  $F \rightarrow bB$  then
  - $D(S) = \{S, A, B, C, E\}$
  - $D(A) = \{A, B, C, E\}$
  - $D(B) = \{B, C, E\}$
  - $D(F) = \{F\}$

- If the replacement sets would have been  $D(A) = \{A, B, D\}$ ,  $D(B) = \{B, D\}$ , and there is a rule  $S \rightarrow AB$ 
  - the new rules would be  $S \rightarrow AB$ ,  $S \rightarrow AD$ ,  $S \rightarrow BB$ ,  $S \rightarrow BD$ ,  $S \rightarrow DB$ ,  $S \rightarrow DD$

- Lemma: the next algorithm determines all  $N[i, i + s]$ 
  - the algorithm is based on dynamic programming
    - first it determines the smallest sub-problems:  
 $N[1, 1], N[2, 2], \dots$
    - then gradually the bigger ones using all the previous results (dynamic programming)

- Definition of  $N[i, i + s]$ :  $\{A \in V : A \Rightarrow^* x_i \dots x_{i+s}, x_j \in V\}$ 
  - given CFG G in Chomsky normal form
  - given string  $x = x_1 x_2 \dots x_n$ ,  $n \geq 2$
  - $N[i, i + s]$  is the set of symbols which can derive in G string  $x_i \dots x_{i+s}$

```
for i:= 1 to n do N[i, i] := {xi}  
for s:= 1 to n - 1 do  
    for i:= 1 to n - s do  
        for k:= i to i + s - 1 do  
            if ∃ A → BC ∈ R, B ∈ N[i, k],  
C ∈ N[k+1, i+s] then  
                add A to N[i, i+s]
```

Legend:

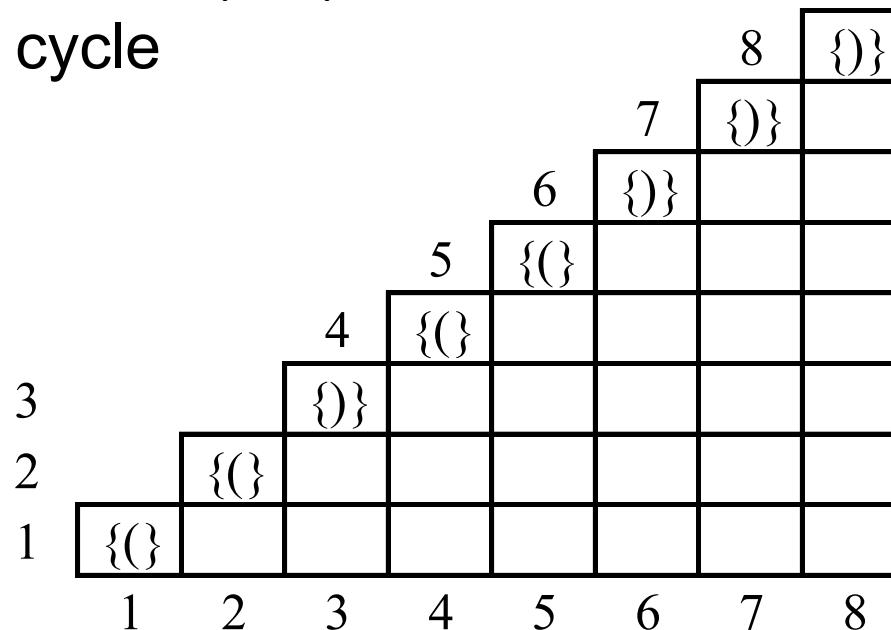
s: the size of the current problem - 1  
i: the start of the first subproblem  
k: the end of the first subproblem  
first subproblem: N[i, k]  
second subproblem: N[k+1, i+s]

- For a given G the complexity is  $O(n^3)$ 
  - calculation of the if depend on  $|R|$  and  $|\Sigma|$
- Main idea:
  - if:  $A \rightarrow BC$ ,  $B \Rightarrow^* x_i \dots x_k$ ,  $C \Rightarrow^* x_{k+1} \dots x_{i+s}$
  - then:  $A \Rightarrow^* x_i \dots x_{i+s}$ 
    - $A \Rightarrow BC \Rightarrow^* x_i \dots x_k C \Rightarrow^* x_i \dots x_k x_{k+1} \dots x_{i+s}$

- Theorem:
  - given:
    - CFG G in Chomsky normal form
    - string  $x = x_1 x_2 \dots x_n$ ,  $n \geq 2$
    - the  $N[i, i + s]$  sets
  - $x \in L(G) \leftrightarrow S \in N[1, n]$
- Proof: according to the definition of  $N[i, i + s]$   
if  $S \in N[1, n] \rightarrow S \Rightarrow^* x_1 x_2 \dots x_n = x$

# Example

- Decide if  $((())())$  can be generated by the next grammar!
  - the CFG generates the set of balanced parenthesis
    - $R = \{S \rightarrow SS, S \rightarrow (S_1, S_1 \rightarrow S), S \rightarrow ()\}$
  - perform the initial cycle



# Example

– calculate the  $s = 1$  iteration

- $N[1, 1] = "(" = \{\{\}\}, N[2, 2] = "(" = \{\{\}\}, \text{there is no rule such } A \rightarrow ((, A \in V \rightarrow N[1, 2] = "()" = \{\})$   
–  $N[1, 2]$  means row 1, column 2
- $N[2, 2] = \{\{\}\}, N[3, 3] = \{\{\}\}, \exists S \rightarrow () \rightarrow N[2, 3] = "()" = \{S\}$

								8	$\{\{\}\}$
								7	$\{\{\}\}$
								6	$\{\{\}\}$
								5	$\{\{\}\} \quad \{S\}$
								4	$\{\{\}\} \quad \{\}$
								3	$\{\{\}\} \quad \{\}$
								2	$\{\{\}\} \quad \{S\}$
								1	$\{\{\}\} \quad \{\}$
	1	2	3	4	5	6	7	8	

# Example

- calculate the table one diagonal after another

1	{}	{}	{}					
2	{}	{S}	{}					
3		{}	{}	{}				
4		{}	{}	{}				
5		{}	{S}	{S1}				
6		{}	{}	{}				
7		{}	{}	{}				
8		{}	{}	{}				

# Example

- a string can be divided in different ways, e.g. to calculate  $N[2, 7] = "(()())"$
- examine each partition  $N[2, k], N[k + 1, 7], 2 \leq k < 7$
- check if  $\exists A \rightarrow BC, B \in N[2, k], C \in N[k+1, 7]$

$k = 3, S \in N[2, 3] = "()", S \in N[4, 7],$

$\exists S \rightarrow SS,$

so  $S$  is added to  $N[2, 7]$


- $S \in N[1, 8]$  so the  $((())())$  can be generated in  $G$

- Definition of consistent strings,  $x$  and  $y$ :  $x$  is a prefix of  $y$  or  $y$  is a prefix of  $x$
- Definition of compatible transitions  $((p, a, \beta), (q, \gamma))$  and  $((p, a', \beta'), (q', \gamma'))$ : ' $a$ ' is consistent with  $a'$  and  $\beta$  is consistent with  $\beta'$ 
  - $\exists$  such a configuration when both transitions are applicable
  - in configuration  $(q, aabba, xyx)$  both  $((q, aa, x), (p, y))$  and  $((q, a, xy), (p, e))$
- Definition of deterministic PDA: it has no two distinct compatible transitions
  - it may have  $e$  transition

- Definition of deterministic context-free language  $L$ :  
 $L\$ = L(M)$  for some deterministic PDA  $M$ 
  - $\$$  is appended to each input string, it marks the end of string,  $\$ \notin \Sigma$
  - $L$  is recognized by such a deterministic PDA  $M$  which can sense the end of the input string
- Theorem: every deterministic context-free language is a context-free language

- If the end of string is not marked  $\rightarrow L = a^* \cup a^n b^n$  can not be recognized by a deterministic PDA
  - the number of 'a' must be counted if we want to check the number of b by pushing symbols into the stack
  - on the other hand  $a^*$  must be also accepted
    - if  $\$$  is met then the stack can be cleared

- We suppose that not every context-free language is a deterministic context-free language
  - e.g.:  $L = \{a^n b^m c^p : m, n, p \geq 0, n \neq m \text{ or } m \neq p\}$
  - $L$  is context-free
  - $L$  can be deterministic context-free language if we know in advance that the 'a' or the 'b' is to be counted

- Definition of dead end configuration of PDA  $M$ ,  $(q, \omega, \alpha)$ : for each configuration  $(q', \omega', \alpha')$  which follows  $(q, \omega, \alpha)$ ,  $(q, \omega, \alpha) \vdash_M^+ (q', \omega', \alpha')$ , the followings are true:  
 $\omega' = \omega$  and  $|\alpha'| \geq |\alpha|$ 
  - no progress can be made starting from dead end in terms of either reading the input or reducing the stack
  - $M$  executes a never-ending sequence of e-moves

- If  $M$  does not have dead end then
  - for any configuration we can find such a following configuration which has less input or less stack
  - the procedure is finite as the input and the number of elements in the stack are finite

- Theorem: the class of deterministic context-free languages is closed under complement
- Proof by construction:
  - let  $M = (K, \Sigma, \Gamma, \Delta, s, F)$  is any deterministic PDA
  - transform  $M$  to  $M'$ , such that
    - $L(M) = L(M')$
    - $M'$  is a simple PDA (and still deterministic)

- $M'$  may reject the input if
  - the input is read and  $M'$  is not in final state
  - the input is read and the stack is not empty
  - $M'$  is stuck as there is no applicable transition
  - $M'$  reaches a dead end (does not read input)

- the dead ends of a simple, deterministic PDA can be determined
  - not detailed here, how exactly
- transform  $M'$  to  $M''$ , such that
  - $L(M') = L(M'')$
  - $M''$  have some new, non-final states
  - $M''$  moves into the new states and empty the stack whenever  $M'$  would reject the input based on condition 2, 3, or 4
  - $M''$  has no dead end
- create  $M'''$  by inverting  $F''$ 
  - $L(M''') = L(M'')^C = L(M')^C = L(M)^C$

- Theorem: the class of deterministic context-free languages is properly contained in the class of context-free languages

- Proof by contradiction:
  - suppose  $L = \{a^n b^m c^p : m, n, p \geq 0, n \neq m \text{ or } m \neq p\}$  is deterministic context-free
  - $L^C = \{a^n b^m c^p : m, n, p \geq 0, n = m = p\} \cup \{w \notin a^* b^* c^*\}$  is deterministic context-free according to the previous theorem
  - $L^C \cap \{w \in a^* b^* c^*\}$  is context-free according to a former theorem (context-free  $\cap$  regular = context free)
  - but  $L^C \cap \{w \in a^* b^* c^*\} = \{w \in a^n b^n c^n\}$  is not context-free

- Interesting:
  - DFA is equivalent with NFA
  - PDA is not equivalent with deterministic PDA
  - Turing-machine is equivalent to nondeterministic Turing-machine



# Introduction to the Theory of Computation

Lesson 4

3.7. Parsing

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BEFEKTETÉS A JÖVŐBE

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- More information in the course: Compilers
- Definition of parsing: given some grammar  $G$ , use a deterministic automaton (parser)
  - to decide if  $w \in L(G)$ 
    - see previous lesson
  - if  $w \in L(G)$ , give the derivation of  $w$

- Top-down parsing constructs the leftmost derivation of  $w$ 
  - start the construction from the root
- CFGs and deterministic PDAs are considered here
- Our aim is to construct a deterministic PDA for a CFG

# Example 1

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Given CFG  $G = (\{a, b, S\}, \{a, b\}, \{S \rightarrow aSb, S \rightarrow e\}, S)$   
construct a parser to it!
  - $L(G) = \{a^n b^n\}$
  - construct a PDA equivalent with G
    - PDA  $M_1 = (\{p, q\}, \{a, b\}, \{a, b, S\}, \Delta_1, p, \{q\})$ 
      - $\Delta_1 = \{((p, e, e), (q, S)), ((q, e, S), (q, aSb)), ((q, e, S), (q, e)), ((q, a, a), (q, e)), ((q, b, b), (q, e))\}$
    - $M_1$  is not deterministic because in configuration  $(q, abb, Sb)$  both transition 2 and 3 are applicable

- Idea of lookahead:
  - non-determinism can be resolved by reading 1 character from the input before decision
  - at this example if
    - next symbol: 'a' then use  $S \rightarrow aSb$
    - next symbol: b then use  $S \rightarrow e$

- in general
  - read 1 symbol without consulting the stack
  - go to a new state corresponding to the read symbol
    - $q_a$  is a state coming from state  $q$  after reading 'a'
  - with the new states a deterministic PDA  $M_2$  can be constructed such that  $L(M_1) = L(M_2)$

- $M_2 = (\{p, q, q_a, q_b, q_{\$}\}, \{a, b\}, \{a, b, S\}, \Delta_2, p, \{q_{\$}\})$ 
  - $\Delta_2$  contains the following transitions:
    - (1) ((p, e, e), (q, S))
    - (2) ((q, a, e), (q<sub>a</sub>, e)) // new
    - (3) ((q<sub>a</sub>, e, S), (q<sub>a</sub>, aSb)) // S is top of the stack
    - (4) ((q<sub>a</sub>, e, a), (q, e))
    - (5) ((q, b, e), (q<sub>b</sub>, e)) // new
    - (6) ((q<sub>b</sub>, e, S), (q<sub>b</sub>, e)) // S is top of the stack
    - (7) ((q<sub>b</sub>, e, b), (q, e))
    - (8) ((q, \\$, e), (q<sub>\\$</sub>, e)) // new
      - we can decide between 3 and 6 based on the state (which is based on the lookahead)

Step	State	Unread input	Stack	Transition used	Rule
0	p	ab\$	e	-	
1	q	ab\$	S	1	
2	q <sub>a</sub>	b\$	S	2	
3	q <sub>a</sub>	b\$	aSb	7	S → aSb
4	q	b\$	Sb	3	
5	q <sub>b</sub>	\$	Sb	4	
6	q <sub>b</sub>	\$	b	8	S → e
7	q	\$	e	5	
8	q <sub>\$</sub>	e	e	6	

- Definition of LL(1) parser: such a deterministic PDA which is constructed to be a top-down parser and to read at most 1 input symbol at a time
- Not all CFG have a deterministic PDA equivalent
- Some have a deterministic PDA equivalent but not an LL(1) parser
  - reading ahead just 1 character may not be enough to resolve uncertainties
- Some grammar need transformations before an LL(1) parser can be constructed

# Example 2

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Given CFG G that generates arithmetic expressions with operations + and \*, construct the LL(1) parser!
- $G = (V, \Sigma, R, E)$ , where
  - $V = \{+, *, (, ), \text{id}, T, F, E\}$ 
    - E - expression, T - term, F - factor
  - $\Sigma = \{+, *, (, ), \text{id}\}$
  - $R = \{E \rightarrow E + T, E \rightarrow T, T \rightarrow T * F, T \rightarrow F, F \rightarrow (E), F \rightarrow \text{id}, F \rightarrow \text{id}(E)\}$ 
    - the last rule corresponds to functions such as  $\sqrt{x * x + 1}$ ,  $f(z)$

# Example 2

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The PDA equivalent with  $G$  contains the following transitions:
  - (0)  $((p, e, e), (q, \mathsf{E}))$
  - (1)  $((q, e, \mathsf{E}), (q, \mathsf{E} + \mathsf{T}))$
  - (2)  $((q, e, \mathsf{E}), (q, \mathsf{T}))$
  - (3)  $((q, e, \mathsf{T}), (q, \mathsf{T}^* \mathsf{F}))$
  - (4)  $((q, e, \mathsf{T}), (q, \mathsf{F}))$
  - (5)  $((q, e, \mathsf{F}), (q, (\mathsf{E})))$
  - (6)  $((q, e, \mathsf{F}), (q, \mathsf{id}))$
  - (7)  $((q, e, \mathsf{F}), (q, \mathsf{id}(\mathsf{E})))$
  - $((q, a, a), (q, e)) \in \Delta$  for  $\forall a \in \Sigma$

- Looking ahead one symbol does not help to choose between transitions 6 and 7
  - they both can be applied in configuration  $(q, \text{id}(5)\$, F)$
- The source of the problem are the rules:  
 $F \rightarrow \text{id}$ ,  $F \rightarrow \text{id}(E)$ 
  - identical left side,  $F$ , and identical first nonterminal,  $\text{id}$

- Construct again the PDA
  - (6') ((q, e, F), (q, idA))
  - (7') ((q, e, A), (q, e))
  - (8') ((q, e, A), (q, (E)))
  - if the lookahead = '(' then transition 8' is to be used

- Given CFG G
  - if it has rules
    - $A \rightarrow \alpha\beta_1 \mid \dots \mid \alpha\beta_n$ 
      - $\alpha \neq e, n \geq 2$
    - then replace them by
      - $A \rightarrow \alpha A', A' \rightarrow \beta_1 \mid \dots \mid \beta_n$ 
        - $A'$  is a new nonterminal

# Eliminating left recursion

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Looking ahead one symbol does not help to choose between transitions 1 and 2 ( belonging to  $E \rightarrow E + T | T$ )
  - they both can be applied in configuration (q, id+id\$, E)
    - if transition 2 is selected then the PDA will stuck
    - if transition 1 is repeated 2, 3, ... times then the PDA will stuck

# Eliminating left recursion

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of left recursion in CFG G: G contains such a rule whose left side is the same as the leftmost element of the right side
  - at the problem in question:  $E \rightarrow E + T$
- Replace the previous rules by:  $E \rightarrow TE'$ ,  $E' \rightarrow +TE' | e$ 
  - if the lookahead = '+' then we know which rule must be used

# Eliminating left recursion

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Construct again the PDA from  $G' = (V', \Sigma, R, E)$

- (1)       $E \rightarrow TE'$
- (2)       $E' \rightarrow +TE'$
- (3)       $E' \rightarrow e$
- (4)       $T \rightarrow FT'$
- (5)       $T' \rightarrow *FT'$
- (6)       $T' \rightarrow e$
- (7)       $F \rightarrow (E)$
- (8)       $F \rightarrow idA$
- (9)       $A \rightarrow e$
- (10)       $A \rightarrow (E)$

# Eliminating left recursion

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Given CFG G
  - if it has rules
    - $A \rightarrow A\alpha_1 \mid \dots \mid A\alpha_n$  and  $A \rightarrow \beta_1 \mid \dots \mid \beta_m$ 
      - for  $\forall i \beta_i$  does not start with A,  $n > 0$
    - e.g. before:  $A \Rightarrow A\alpha_1 \Rightarrow^* A\alpha_1\dots\alpha_1 \Rightarrow \beta_1\alpha_1\dots\alpha_1$
  - then replace them by
    - $A \rightarrow \beta_1A' \mid \dots \mid \beta_mA'$ ,  $A' \rightarrow \alpha_1A' \mid \dots \mid \alpha_nA' \mid e$ 
      - $A'$  is a new nonterminal
    - e.g. after:  $A \Rightarrow \beta_1A' \Rightarrow \beta_1\alpha_1A' \Rightarrow^* \beta_1\alpha_1\dots\alpha_1A' \Rightarrow \beta_1\alpha_1\dots\alpha_1$

- From given CFG create a deterministic PDA
  - perform left factoring
  - eliminate left recursion
  - perform CFG → PDA transformation
  - introduce lookahead states
  - apply the PDA to the word
- Inspecting the stack we can see the rules creating a word
- Not every CFG can be transformed in this way

- Bottom-up parsing constructs the rightmost derivation of  $w$ 
  - start the construction from the leaves
- There is a new method to construct a PDA equivalent with a CFG
  - given  $G = (V, \Sigma, R, S)$

- PDA  $M = (K, \Sigma, \Gamma, \Delta, p, F)$ ,  $K = \{p, q\}$ ,  $\Gamma = V$ ,  $F = \{q\}$ 
  - $\Delta$  contains the following transitions:
    - $((p, a, e), (p, a))$ , for  $\forall a \in \Sigma$ 
      - » shift the input into the stack symbol by symbol
    - $((p, e, \alpha^R), (p, A))$ , for  $\forall$  rule  $A \rightarrow \alpha$  in  $R$ 
      - » if a reverse of the right side of a rule is at the top of the stack then replace it with the corresponding left side
    - $((p, e, S), (q, e))$ 
      - » end operation by removing the start symbol
  - $L(G) = L(M)$

# Bottom-up parsing

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $G = (V, \Sigma, R, E)$ 
  - $V = \{+, *, (, ), \text{id}, T, F, E\}$
  - $\Sigma = \{+, *, (, ), \text{id}\}$
  - $R =$ 
    - $E \rightarrow E + T \quad (R1)$
    - $E \rightarrow T \quad (R2)$
    - $T \rightarrow T * F \quad (R3)$
    - $T \rightarrow F \quad (R4)$
    - $F \rightarrow (E) \quad (R5)$
    - $F \rightarrow \text{id} \quad (R6)$

- $\Delta$  of M:
  - $(p, a, e), (p, a)) \text{ for } \forall a \in \Sigma$   $(\Delta 0)$
  - $(p, e, T + E), (p, E))$   $(\Delta 1)$
  - $(p, e, T), (p, E))$   $(\Delta 2)$
  - $(p, e, F * T), (p, T))$   $(\Delta 3)$
  - $(p, e, F), (p, T))$   $(\Delta 4)$
  - $(p, e, )E(), (p, F))$   $(\Delta 5)$
  - $(p, e, id), (p, F))$   $(\Delta 6)$
  - $(p, e, E), (q, e))$   $(\Delta 7)$

step	state	Unread input	Stack	Transition used	Rule
0	p	$\text{id}^*(\text{id})$	e		
1	p	$^*(\text{id})$	id	$\Delta 0$	
2	p	$^*(\text{id})$	F	$\Delta 6$	R6
3	p	$^*(\text{id})$	T	$\Delta 4$	R4
4	p	( $\text{id}$ )	$^*T$	$\Delta 0$	
5	p	$\text{id})$	$(^*T$	$\Delta 0$	
6	p	)	$\text{id}(^*T$	$\Delta 0$	
7	p	)	$F(^*T$	$\Delta 6$	R6
8	p	)	$T(^*T$	$\Delta 4$	R4
9	p	)	$E(^*T$	$\Delta 2$	R2
10	p	e	$)E(^*T$	$\Delta 0$	
11	p	e	$F^*T$	$\Delta 5$	R5
12	p	e	T	$\Delta 3$	R3
13	p	e	E	$\Delta 2$	R2
14	q	e	e	$\Delta 7$	

- The implied rightmost derivation:
  - apply the rules from end to start

$$\begin{aligned} E &\Rightarrow T \\ &\Rightarrow T^* F \\ &\Rightarrow T^* (E) \\ &\Rightarrow T^* (T) \\ &\Rightarrow T^* (F) \\ &\Rightarrow T^* (id) \\ &\Rightarrow F^* (id) \\ &\Rightarrow id^* (id) \end{aligned}$$

- M is nondeterministic:  $\Delta_0$  is compatible with all the transitions of type  $\Delta_1$  through  $\Delta_8$
- To make M deterministic the next conflicts must be resolved
  - shift or reduce
  - reduce or another reduce

- Resolving the shift-reduce conflict
- Construct precedence relation  $P$ 
  - based on
    - $b$  = next input symbol
    - $a$  = top stack symbol
  - $P \subseteq V \times (\Sigma \cup \{\$\})$
  - the construction of  $P$  is not given here
- Decision:
  - reduce if  $(a, b) \in P$
  - shift otherwise

# Bottom-up parsing

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Relation P for the example, e.g.:  $(T, +) \in P$ , if there is a mark in the table

	(	)	id	+	*	\$	b
(							
)		✓		✓	✓	✓	
id		✓		✓	✓	✓	
+							
*							
E							
T		✓		✓		✓	
F		✓		✓	✓	✓	
a							

- Resolving the reduce-reduce conflict: reduce the longest possible string from the top of the stack
- Definition of weak precedence grammars: the bottom-up parser corresponding to this grammar is deterministic if the shift-reduce and reduce-reduce conflicts are resolved



# Introduction to the Theory of Computation

## Lesson 5

### 4.3. Extensions of Turing machines

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- Definition of:
  - Turing machine
  - basic machines
  - Machine schema (variable like notation)
  - simple machines (copy, delete, shift, search)

- Turing machines can be extended
- The extensions:
  - do not increase computational power
    - we will see an extended TM that can be simulated with a conventional TM
  - do increase speed and ease of use

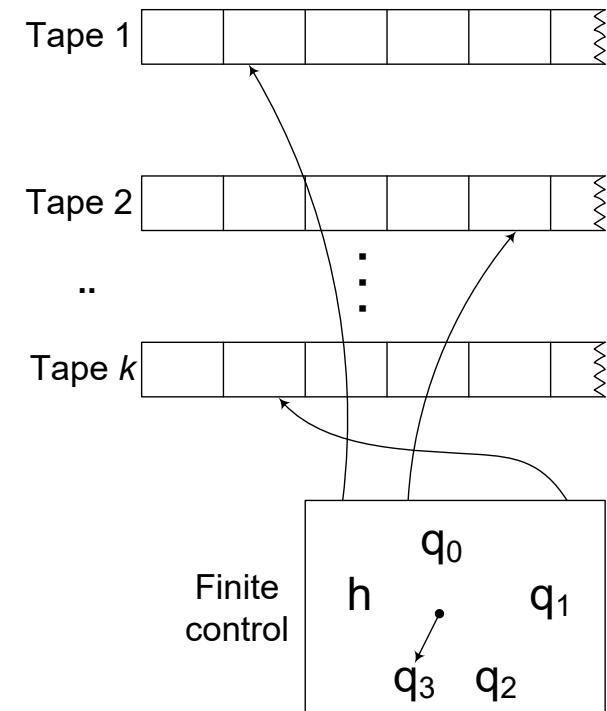
- Possible extensions of the Turing machine:
  - multiple tapes
  - two-way infinite tapes
  - two-dimensional tapes
  - multiple heads
  - random access
  - nondeterministic TM

# Multiple tape TM

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- k-tape TM: a TM equipped with  $k$  tapes and corresponding heads ( $k \geq 1$ )
  - each tape is connected to the finite control by means of a read/write head



- Definition of k-tape TM: a quintuple  $(K, \Sigma, \delta, s, H)$ , where
  - $k \in \mathbb{N}, k \geq 1$
  - $K$  set of states (finite)
  - $\Sigma$  alphabet (finite)
    - containing  $\sqcup, \triangleright$ , not containing  $\leftarrow, \rightarrow$
  - $s \in K$  the initial state
  - $H \in K$  the set of halting states (finite)

–  $\delta$  transition function,  $(K - H) \times \Sigma^k \rightarrow K \times (\Sigma \cup \{\leftarrow, \rightarrow\})^k$

- $\forall q \in K - H, \forall (a_1, \dots, a_k) \in \Sigma^k, \delta(q, (a_1, \dots, a_k)) = (p, (b_1, \dots, b_k))$ 
  - $p$ : the new state
  - $b_j$ : the action taken by M at tape j
- if  $\exists j$  such that  $a_j = \triangleright \rightarrow b_j = \rightarrow$
- there is no  $j$  such that  $b_j = \triangleright$
- in a single step for each tape the actual symbol can be modified or the head can be moved

- Definition of the configuration of a k-tape TM: an element of  $K \times (\triangleright \Sigma^* \times (\Sigma^*(\Sigma - \{\sqcup\}) \cup \{e\}))^k$ 
  - identify the state, the tape contents, and the head positions for each tape

- The initial configuration
  - contains the input on the 1<sup>st</sup> tape the same way as on simple TM
  - the other tapes are empty with heads on the leftmost blank position
- The output of a k-tape TM is stored on the first tape

- Construct the copying machine with a 2-tape TM:
  - transform  $\triangleright \underline{L} w \underline{L}$  into  $\triangleright \underline{L} w \underline{L} w \underline{L}$ ,  $w \in \{a, b\}^*$
- Operation:
  - start:
    - first tape:  $\triangleright \underline{L} w$ , second tape:  $\triangleright \underline{L}$

- move the heads on both tapes to the right, copying each symbol from the first tape onto the second tape
  - the first square of the second tape should be left blank
  - first tape:  $> \underline{L} w \underline{L}$ , second tape:  $> \underline{L} w \underline{L}$
  - part of  $\delta$ :  $\delta(q_1, a, \underline{L}) = (q_2, a, a)$ ,  
 $\delta(q_2, a, a) = (q_1, \rightarrow, \rightarrow)$

# Example

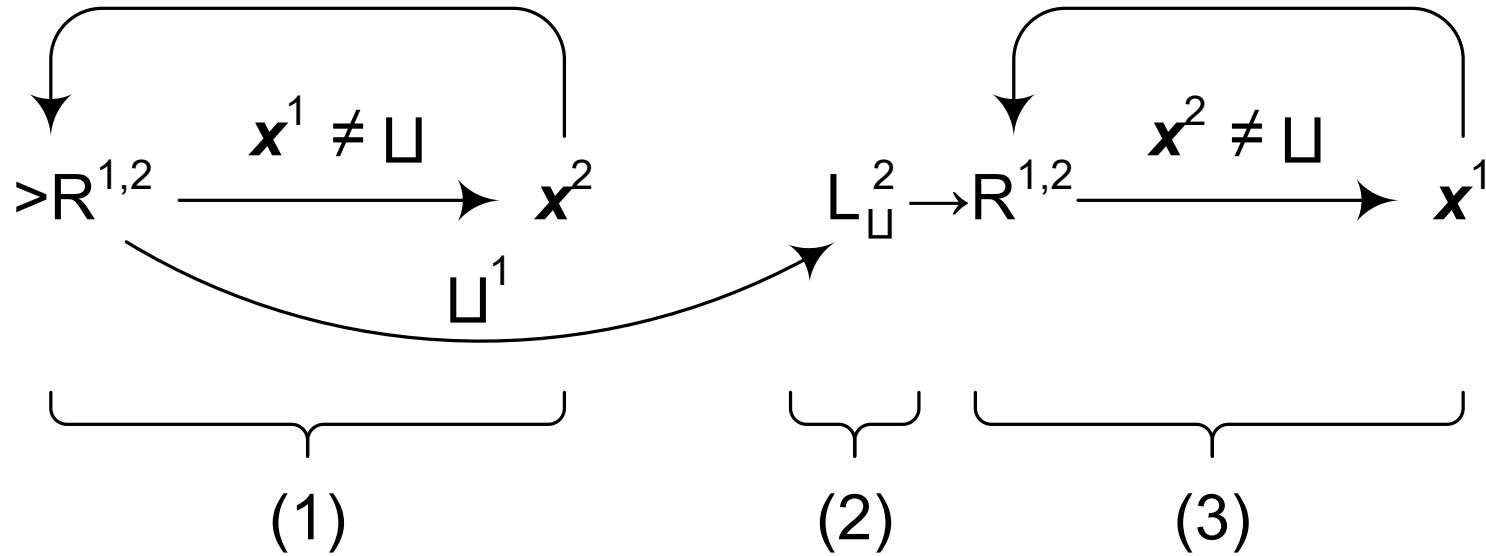
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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- move the head on the second tape to the left until blank is found
  - first tape:  $\triangleright \underline{L} w \underline{L}$ , second tape:  $\triangleright \underline{L} w$
- move the heads on both tapes to the right copying symbols from the second tape onto the first tape
  - first tape:  $\triangleright \underline{L} w \underline{L} w \underline{L}$ , second tape:  $\triangleright \underline{L} w \underline{L}$
- It is easier to implement C (the copy machine) on a 2-tape machine and the operation is quicker

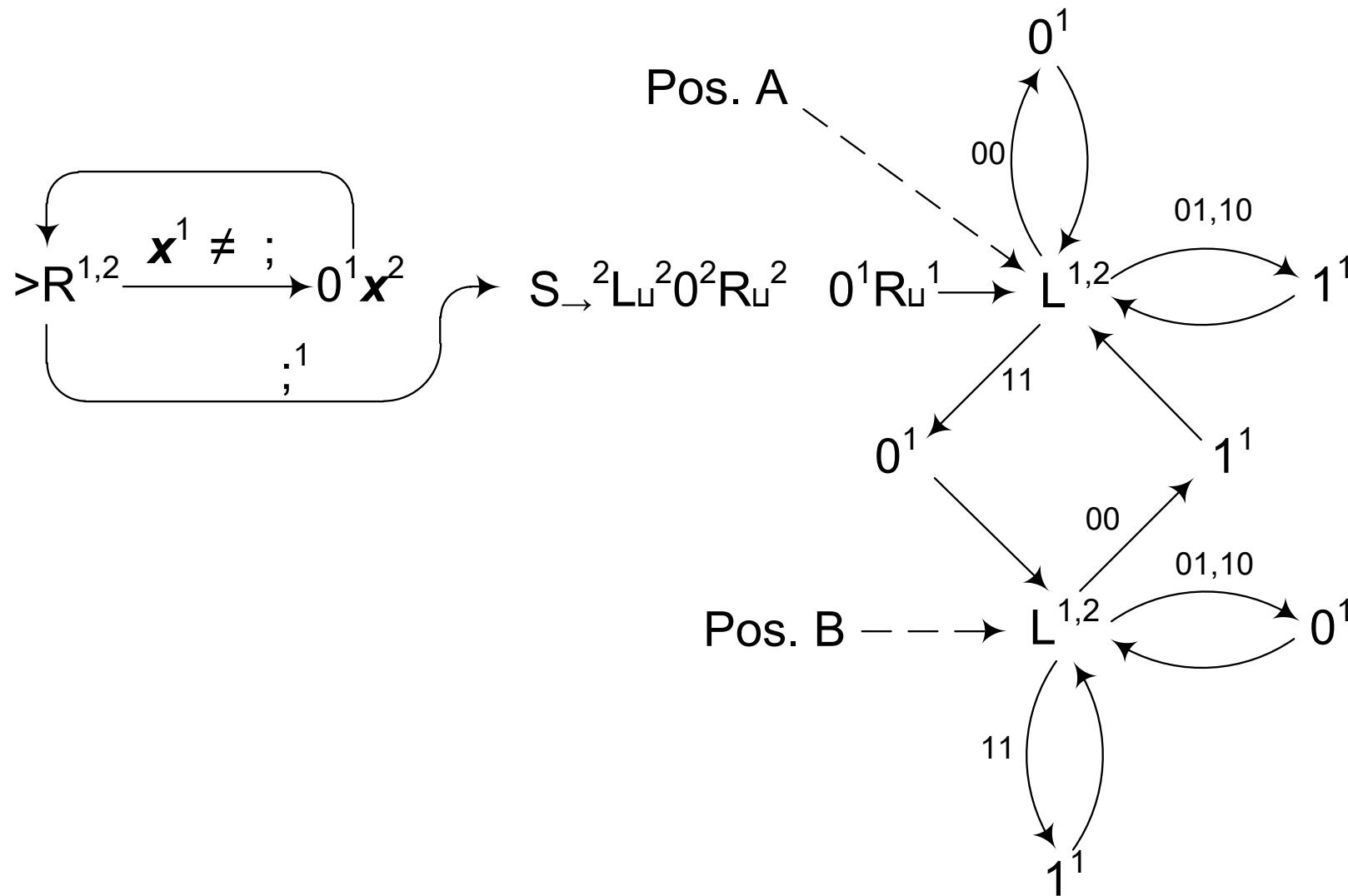
- We use the following conventions to depict a k-tape TM:
  - a superscript is attached to each sub-TM denoting the tape on which it operates
    - e.g.:
      - $U^2$  writes a blank on the second tape
      - $L_U^1$  searches to the left for a blank on the first tape
      - $R^{1,2}$  is a shorthand for  $R^1R^2$
    - beware:  $R^2$  meant  $RR$  before

- a superscript is attached to the symbols on the arrows
  - e.g.:  $a^1$  means "follow this arrow if the previous sub-TM finished and 'a' is under the head on the first tape"
  - a pair of bits, e.g. 01, on an arrow label is a shorthand for  $a^1 = 0, a^2 = 1$



- Construct a machine to add two binary numbers with a 2-tape TM:
  - $n_i \in \mathbb{N}$ ,  $w_i = \text{numBin}(n_i)$ ,  $i = 1, 2, 3$ ,  $n_3 = n_1 + n_2$
  - transform  $\triangleright \underline{w_1; w_2} \sqcup$  into  $\triangleright \sqcup 0^* w_3 \sqcup$ ,  $w_i \in \{0, 1\}^*$

# Example



# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Operation:
  - start:
    - first tape:  $\triangleright \underline{L} w_1; w_2 \underline{L}$ , second tape:  $\triangleright \underline{L}$
  - copies the first binary integer in its second tape writing zeros in its place
    - and in the place of the ';' separating the two integers

- first tape:  $\triangleright \underline{0}^* w_2 \underline{\triangleright}$ , second tape:  $\triangleright \underline{0} w_1 \underline{\triangleright}$ 
  - $S \rightarrow^2 L \downarrow^2 0^2 R \downarrow^2 : \triangleright \underline{w}_1 \underline{\triangleright} \rightarrow \triangleright \underline{0} w_1 \underline{\triangleright}$ 
    - » the adding stops when one of the numbers is finished, now the carry bit will be written

- performs binary addition by the "school method", starting from the least significant bit of both integers
  - writes the result in the first tape
  - sub-machines indicate if there is a carry bit
    - » there is carry bit at the lower  $L^{1,2}$  machine
  - first tape:  $\triangleright \underline{U} 0^* w_3 \underline{U}$ , second tape:  $\triangleright \underline{U} w_1 \underline{U}$

- T1: >010;011, T2: >
- T1: >0000011, T2: >010 // Pos. A
- T1: >0000011, T2: >010 // Pos. A
- T1: >0000011, T2: >010
- T1: >0000011, T2: >010 // Pos. A
- T1: >0000001, T2: >010
- T1: >0000001, T2: >010 // Pos. B
- T1: >0000101, T2: >010
- T1: >0000101, T2: >010 // Pos. A

- Theorem:
  - let  $M = (K, \Sigma, \delta, s, H)$  is a  $k$ -tape TM,  $k \geq 1$
  - $\exists$  a standard TM  $M' = (K', \Sigma', \delta', s', H')$ ,  $\Sigma \subseteq \Sigma'$  such that:
    - $M$  halts on input  $x$  with output  $y$  on its first tape  
 $\leftrightarrow M'$  halts on input  $x$  at the same halting state and with the same output
    - if on input  $x$   $M$  halts after  $t$  steps  $\rightarrow M'$  halts after  $O(t^*(|x|+t))$  steps
  - $k$ -tape TM and standard TM are equivalent

# Corollary from the k-tape Turing machine

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A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Proof: in the book
  - idea: each cell of tape of  $M'$  contains ( $2k$ -tuples)
  - it means  $\Sigma' = (\Sigma \times \{0, 1\})^k$
  - $\{0, 1\}$  is needed to sign the head position of a head

# Corollary from the k-tape Turing machine

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Corollary: Any function that is computed or language that is decided or semidecided by a k-tape Turing machine is also computed, decided, or semidecided, respectively, by a standard Turing machine

- The tape is infinite in both directions
  - there is no  $\triangleright$  symbol
- All squares are initially blank except those containing the input, tape:  $\underline{L}w$

- It can be simulated by a 2-tape machine:
  - tape 1 contains the part of the tape to the right of the square containing the first input symbol
  - tape 2 contains the part of the tape to the left of the square containing the first input symbol in reverse order
- The simulation needs only linear (and not quadratic) time because in each step only one tape is active

- TM has a single tape and multiple heads
- In one step, the heads all sense the scanned symbols and move or write independently
  - some convention must be used what happens if two heads try to modify the same symbol

- Idea for simulation:
  - the tape is divided into tracks
    - one corresponds to the real tape
    - the others contain only the head positions
- The simulation needs quadratic time

- Construct a 3-heads TM to decide  $\{a^n b^n c^n : n \in N\}$ !
  - position head 1 to the first 'a', head 2 to the first b, head 3 to the first c
  - enter a loop which moves all head to the right until head 1 reads 'a', head 2 reads b, head 3 read c
  - if at the same time head 1 reads b, head 2 reads c, and head 3 reads  $\sqcup$  → the input is accepted

- Construct a 2-heads TM to implement the copying machine!
  - position head 1 to the first input symbol, head 2 to the second  $\sqcup$  after the input
  - move both heads to the right while head 2 writes that head 1 reads

- TM store data on a two-dimensional grid instead of a tape
  - infinite to the right and to the bottom
  - there is end of tape markers on the left and the top sides
  - remember both  $N$  and  $N^2$  are countably infinite
- The input is placed into the second row
- The head can move on in four different directions
- The simulation needs polynomial time

- The different extension can be combined
- Example:
  - k-tape TM with multiple heads each
  - multiple 2-way infinite tape TM
  - TM with 4 dimensional tape

- Theorem: Any language decided or semidecided, and any function computed by Turing machines with several tapes, heads, two-way infinite tapes, or multi-dimensional tapes, can be decided, or computed, respectively, by a standard Turing machine



# Introduction to the Theory of Computation

## Lesson 6

### 4.3. Random access Turing machines

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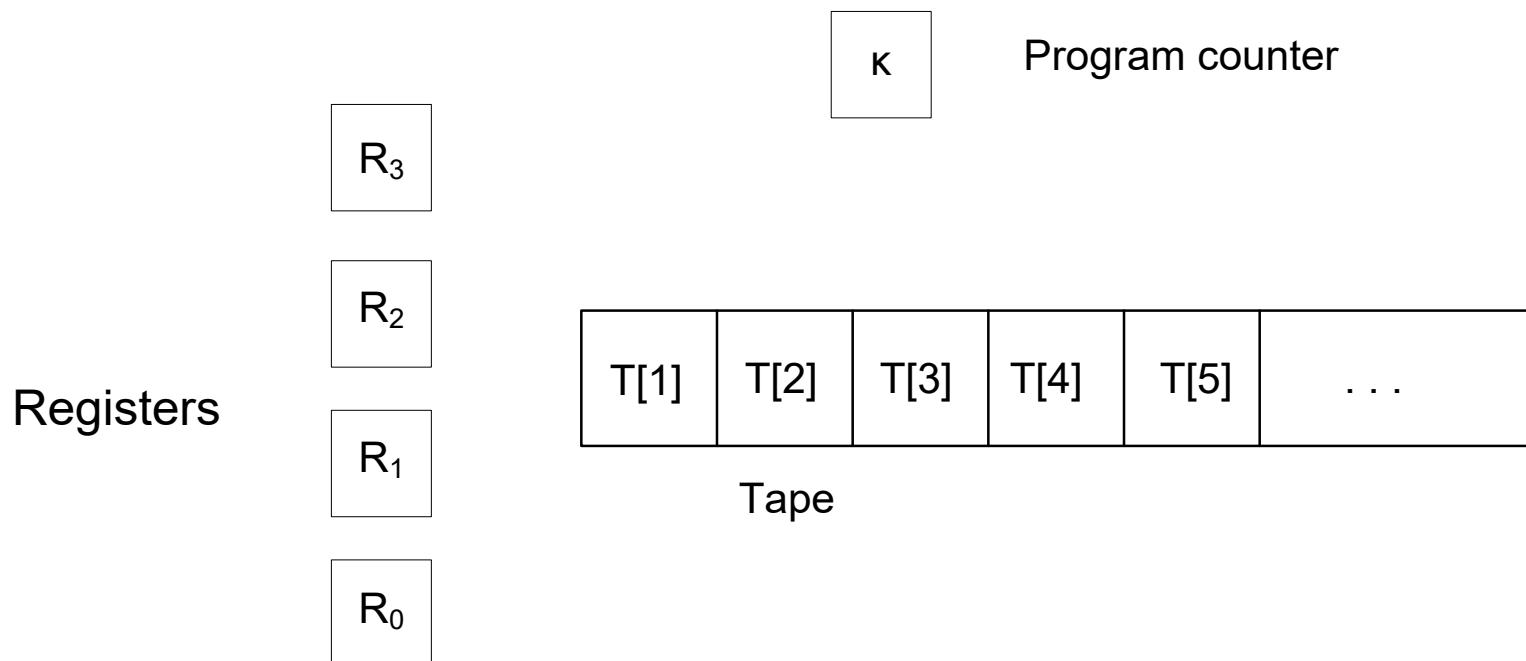
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# Random access TM

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen



- The TM and its previous enhancements had a common drawback their memories (tapes) are sequential
  - the head must be moved square by square to the desired location before a symbol can be read or written
  - real life computer has random access memory

- A random access TM has:
  - a fixed number of registers
    - initially they contain 0s
  - a one-way infinite tape
    - the registers and the tape squares can be accessed directly
    - the registers and the tape squares can contain natural numbers
    - the tape contains the input

- a program containing a finite sequence of instructions
- a program counter showing which instruction is to be executed next

Instruction	Operand	Semantics
read	j	$R_0 := T[R_j]$
write	j	$T[R_j] := R_0$
load	j	$R_0 := R_j$
load	=c	$R_0 := c$
store	j	$R_j := R_0$
add	j	$R_0 := R_0 + R_j$
add	=c	$R_0 := R_0 + c$
sub	j	$R_0 := \max\{R_0 - R_j, 0\}$

Instruction	Operand	Semantics
sub	=c	$R_0 := \max\{R_0 - c, 0\}$
half		$R_0 := \lfloor R_0 / 2 \rfloor$
jump	s	$\kappa := s$
jpos	s	if $R_0 > 0$ then $\kappa := s$
jzero	s	if $R_0 = 0$ then $\kappa := s$
halt		$\kappa := 0$

- $T[i]$  denotes the current contents of tape square  $i$
- $R_j$  denotes the current contents of register  $j$ ,  $0 \leq j < k$ 
  - $R_0$ , called accumulator, has a special role, all arithmetic and logical operations use this register

- $s$  denotes any instruction number in the program
- $c \in \mathbb{N}$
- All instructions change  $\kappa$  to  $\kappa+1$ , unless explicitly stated otherwise

- Definition of a random access TM: a pair  $M = (k, \Pi)$ , where:
  - $k > 0$ , the number of registers
  - $\Pi = (\pi_1, \pi_2, \dots, \pi_p)$ , the program, which is a finite sequence of instructions
    - we assume that  $\pi_p = \text{halt}$ , though the program may contain other halt instructions

- Definition of the configuration of a random access TM:  
a  $(k+2)$ -tuple of the current value of the program counter, the registers, and the tape
  - denoted by  $C = (k, R_0, R_1, \dots, R_{k-1}, T)$
  - $C \in \{0, 1, \dots p\} \times N^k \times T$

- if  $\kappa = 0 \rightarrow$  halted configuration
- $T \subseteq \{N-\{0\}\} \times \{N-\{0\}\}$ , there are no pairs  $(i, m), (j, n)$  such that  $i = j$ 
  - $(i, m) \in T \leftrightarrow T[i] = m$
  - if  $T[i]$  is not defined  $\rightarrow T[i] = 0$

- Definition of the yields in one step relation of a random access TM,  $\vdash_M$ : a relation between two "neighboring" configurations
  - $C = (\kappa, R_0, R_1, \dots, R_{k-1}, T)$ ,  $C' = (\kappa', R'_0, R'_1, \dots, R'_{k-1}, T')$
  - if the difference in the values of  $\kappa$  and  $\kappa'$ ,  $R_i$  and  $R'_i$ ,  $T$  and  $T'$  is exactly the difference caused by  $\pi_\kappa \rightarrow C \vdash C'$

– e.g.:

- if  $\pi_k = \text{read } j \rightarrow R_0' = T[R_j]$ ,  $k' = k+1$ , the other components are unchanged
- if  $\pi_k = \text{add } =c \rightarrow R_0' = R_0 + c$ ,  $k' = k+1$ , the other components are unchanged
- and so on

- Reading and writing the program of a random access TM is extremely cumbersome
- The conventional statements of an algorithm can be translated into a sequence of instructions
  - e.g.:  $R_1 := R_1 + R_2 - 1$  is an abbreviation of:
    1. load 1
    2. add 2
    3. sub =1
    4. store 1
  - the registers can be renamed to enhance readability
    - let  $R_1 := x ; R_2 := y \rightarrow x := x + y - 1$

```
while x > 0 do
    x := x - 3
```

– is an abbreviation of: ( $R_1 := x$ )

1. load 1
2. jzero 6
3. sub =3
4. store 1
5. jump 1
6. halt

– if labels are introduced before the important  
instructions → instruction numbers can be omitted

- Construct a random access TM to implement instruction  $\text{mply } j$ , which performs  $R_0 := R_0 * R_j$ !
  - suppose  $R_0$  initially contains  $x$  and  $R_1$  initially contains  $y$ , after TM will halt  $R_0$  will contain  $x*y$ 
    - the instructions to move  $x$  and  $y$  from the tape to the registers is neglected now
- This new instruction can be used later as a function
  - if  $\text{mply } j$  is given then it is easy to devise  $\text{mply } =c$

- Multiplication is done by successive additions
- The simple algorithm needs  $O(y)$  iterations
  - w contains the result

```
w := 0
while y > 0 do
begin
    w := w + x
    y := y - 1
end
halt
```

- Idea: perform the multiplication using "smaller" multiplications
  - e.g.:  $345 \cdot 247 = 345 \cdot (2 \cdot 10^2 + 4 \cdot 10^1 + 7 \cdot 10^0) = 2 \cdot 345 \cdot 10^2 + 4 \cdot 345 \cdot 10^1 + 7 \cdot 345 \cdot 10^0$ 
    - $7 = 247 \bmod 10$
    - $4 = \lfloor 247/10 \rfloor \bmod 10$
    - $2 = \lfloor 247/100 \rfloor \bmod 10$

$$\begin{array}{r} 345 \cdot 247 \\ 690 \\ 1380 \\ \hline 2415 \\ \hline 85215 \end{array}$$

- School method for multiplication
  - w contains the result
  - it needs  $O(\log y)$  step
  - contains invalid instructions

```
w := 0
for each digit k of yorig
begin
    b := ⌊yorig / 10k-1⌋ mod 10
    x := xorig10k-1
    w := w + bx
end
halt
```

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Idea: imagine  $y$  as a binary number

- e.g.:  $x_{\text{orig}} = 3, y_{\text{orig}} = 5 = 101_2$

$$3 \cdot 5 = 3 \cdot (1 \cdot 2^2 + 0 \cdot 2^1 + 1 \cdot 2^0) = \underbrace{1 \cdot 3 \cdot 2^2}_{b3} + \underbrace{0 \cdot 3 \cdot 2^1}_{x3} + \underbrace{1 \cdot 3 \cdot 2^0}_{w1}$$
$$\underbrace{\underbrace{b3 + x3}_{w2} + \underbrace{w1}_{w3}}_{w2 + w3}$$

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The advanced algorithm using binary representation
  - it contains only valid instructions

```
w := 0
while y > 0 do
begin
    z := half(y)
    if y - z - z ≠ 0 then w := w + x
    x := x + x
    y := z
end
halt
```

```
w := 0
for each digit k of yorig
begin
    b := ⌊yorig / 10k-1⌋ mod 10
    x := xorig10k-1
    w := w + bx
end
halt
```

– calculating  $b_k$  (the k-th least significant bit of y):

- $z := \text{half}(y)$ , if  $y - z - z \neq 0$
- e.g.:  $\lfloor y_{\text{orig}} / 2^{k-1} \rfloor \bmod 2$ 
  - $b_1 = \lfloor 5 / 2^{1-1} \rfloor \bmod 2 = 5 - 2 - 2 = 1$  (last bit)
  - $b_2 = \lfloor 5 / 2^{2-1} \rfloor \bmod 2 = 2 - 1 - 1 = 0$
  - $b_3 = \lfloor 5 / 2^{3-1} \rfloor \bmod 2 = 1 - 0 - 0 = 1$ 
    - »  $\lfloor y_{\text{orig}} / 2^{k-1} \rfloor$  shift to the right
    - » mod 2 check the last bit

– shifting to the right

- $y := z$

– calculating  $x_k = x_{\text{orig}} 2^{k-1}$

- $x := x + x$

- e.g.:

$$- x_1 = 3 * 2^{1-1} = x_{\text{orig}} = 3$$

$$- x_2 = 3 * 2^{2-1} = 3+3 = 6$$

$$- x_3 = 3 * 2^{3-1} = 6+6 = 12$$

– calculating  $w_k = \sum_k b_k x_k$

- if  $b_k$  is 1  $\rightarrow w := w + x$
- e.g.:

$$- w_1 = \sum_1 b_k * x_k = 1 * 3 * 2^{1-1} = 1 * 3 = 3$$

$$- w_2 = \sum_2 b_k * x_k = 0 * 3 * 2^{2-1} + 1 * 3 * 2^{1-1} = 0 * 6 + 1 * 3 = 3$$

$$\begin{aligned} - w_3 &= \sum_3 b_k * x_k = 1 * 3 * 2^{3-1} + 0 * 3 * 2^{2-1} + 1 * 3 * 2^{1-1} = \\ &1 * 12 + 0 * 6 + 1 * 3 = 15 \end{aligned}$$

- Variables at the start of an iteration:
  - $R_1 = y$ , contains  $\lfloor y_{\text{orig}} / 2^{k-1} \rfloor$
  - $R_2 = x_k$ , contains  $x_{\text{orig}} * 2^{k-1}$
  - $R_3 = z$ , contains  $\lfloor y_{\text{orig}} / 2^k \rfloor$
  - $R_4 = w_k$ , contains  $\sum_k b_k * x_k$  (partial result)
- Initial values:  $x = x_{\text{orig}}$ ,  $y = y_{\text{orig}}$ ,  $w = 0$ ,  $z = 0$
- Result:  $w$

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The values of the variables in the algorithm for  $x = 3$ ,  $y = 5$

Iteration:	0	1	2	3
x:	3	6	12	24
y:	5	2	1	0
z:	0	2	1	0
w:	0	3	3	15

- The advanced algorithm translated to instructions  
(R0 = R2 = x, R1 = y)

- |             |             |
|-------------|-------------|
| 1. store 2  | 11. add 2   |
| 2. load 1   | 12. store 4 |
| 3. jzero 19 | 13. load 2  |
| 4. half     | 14. add 2   |
| 5. store 3  | 15. store 2 |
| 6. load 1   | 16. load 3  |
| 7. sub 3    | 17. store 1 |
| 8. sub 3    | 18. jump 2  |
| 9. jzero 13 | 19. halt    |
| 10. load 4  |             |

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The computation of  $5*3$  (and not  $3*5$  as previously):
  - $x = 5, y = 3$
  - $5*3 = 5*(1*2^1 + 1*2^0) = 1*5*2^1 + 1*5*2^0$

Iteration:	0	1	2
x:	5	10	20
y:	3	1	0
z:	0	1	0
w:	0	5	15

- The computation of  $5^*3$ :

(1; 5, 3, 0, 0, 0;  $\emptyset$ ) | $-$  (2; 5, 3, 5, 0, 0;  $\emptyset$ ) | $-$   
(3; 3, 3, 5, 0, 0;  $\emptyset$ ) | $-$  (4; 3, 3, 5, 0, 0;  $\emptyset$ ) | $-$   
(5; 1, 3, 5, 0, 0;  $\emptyset$ ) | $-$  (6; 1, 3, 5, 1, 0;  $\emptyset$ ) | $-$   
(7; 3, 3, 5, 1, 0;  $\emptyset$ ) | $-$  (8; 2, 3, 5, 1, 0;  $\emptyset$ ) | $-$   
(9; 1, 3, 5, 1, 0;  $\emptyset$ ) | $-$  (10; 1, 3, 5, 1, 0;  $\emptyset$ ) | $-$   
(11; 0, 3, 5, 1, 0;  $\emptyset$ ) | $-$  (12; 5, 3, 5, 1, 0;  $\emptyset$ ) | $-$   
(13; 5, 3, 5, 1, 5;  $\emptyset$ ) | $-$  (14; 5, 3, 5, 1, 5;  $\emptyset$ ) | $-$   
(15; 10, 3, 5, 1, 5;  $\emptyset$ ) | $-$  (16; 10, 3, 10, 1, 5;  $\emptyset$ ) | $-$   
(17; 1, 3, 10, 1, 5;  $\emptyset$ ) | $-$  (18; 1, 1, 10, 1, 5;  $\emptyset$ ) | $-$   
(2; 1, 1, 10, 1, 5;  $\emptyset$ ) | $-^*$  (18; 0, 0, 20, 0, 15;  $\emptyset$ ) | $-$   
(2; 0, 0, 20, 0, 15;  $\emptyset$ ) | $-$  (3; 0, 0, 20, 0, 15;  $\emptyset$ ) | $-$   
(19; 0, 0, 20, 0, 15;  $\emptyset$ )

# Encoding of the tape

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Let
  - $\Sigma$  is an alphabet such that  $\sqcup \in \Sigma, \triangleright \notin \Sigma$
  - $E: \Sigma \rightarrow \{0, 1, \dots |\Sigma|-1\}$  is any bijection
    - $E(\sqcup) = 0$
  - $w = a_1a_2\dots a_n \in \Sigma^*$  is encoded on the tape as  
 $T = \{(1, E(a_1)), (2, E(a_2)), \dots (n, E(a_n))\}$

- A random access TM accepts a string if the initial configuration yields a halted configuration with  $R_0 = 1$
- A random access TM rejects a string if the initial configuration yields a halted configuration with  $R_0 = 0$

- Construct a random access TM deciding the language  $\{a^n b^n c^n : n \geq 0\}$  !

```
aCount := bCount := cCount := 0, n := 1
```

```
while T[n] = 1 do :
```

```
    n := n +1, aCount := aCount +1
```

```
while T[n] = 2 do :
```

```
    n := n +1, bCount := bCount +1
```

```
while T[n] = 3 do :
```

```
    n := n +1, cCount := cCount +1
```

```
if aCount = bCount = cCount and T[n] = 0 then
    load =1, halt
```

```
else load =0, halt
```

- $E(a) = 1, E(b) = 2, E(c) = 3$
- We are using the variables account, bcount, and ccount to stand for the number of 'a', b, and c

- Theorem: any recursive or recursively enumerable language, and any recursive function, can be decided, semidecided, and computed, respectively, by a random access TM
- Proof: book

- Theorem: any language decided or semidecided by and any function computable by a random access TM can be decided, semidecided, and computed, respectively, by a standard TM
  - the simulation needs polynomial time
- Corollary: standard and random access TM are equivalent



# Introduction to the Theory of Computation

Lesson 7

4.4. Nondeterministic Turing machines

4.6. Grammars

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- Definition of nondeterministic TM: a quintuple  $(K, \Sigma, \Delta, s, H)$ , where
  - $\Delta \subseteq ((K - H) \times \Sigma) \times (K \times (\Sigma \cup \{\leftarrow, \rightarrow\}))$ 
    - it has transition relation (not function)

- Definition of word accepted by nondeterministic TM:  
 $w \in (\Sigma - \{\triangleright, \sqcup\})^*$  is accepted by  $M$  if  $\exists h \in H, u, v \in \Sigma^*, a \in \Sigma$  such that  $(s, \triangleright \sqcup w) \xrightarrow{M}^* (h, \triangleright \sqcup uav)$ 
  - at least one final configuration is reachable from the initial configuration through yield operations
  - the input may lead to a lot of non-halting configurations
- $M$  accepts (semidecides) a language  $L$  if  $M$  accepts all words in  $L$

- Definition of deciding a language  $L$  by nondeterministic TM  $M = (K, \Sigma, \Delta, s, \{y, n\})$ :
  - if  $\forall w \in L, L \subseteq (\Sigma - \{\sqcup, \triangleright\})^*$
  - then
    - $\exists n \in \mathbb{N}$  depending on  $M$  and  $w$  such that there is no configuration  $C$  such that  $(s, \triangleright \underline{w}) \vdash_M^n C$ 
      - the length of any computation is bounded by  $n$

- $w \in L \leftrightarrow \exists u, v \in \Sigma^*, a \in \Sigma \text{ such that } (s, \triangleright \underline{w}) \vdash_{\stackrel{*}{M}} (y, \triangleright \underline{uav})$ 
  - at least one accepting configuration is reachable from the initial configuration
  - the input may lead to a lot of rejecting configuration

- Definition of computing a function  $f$  by nondeterministic TM  $M$ :
  - if  $w \in L$ ,  $L \subseteq (\Sigma - \{\sqcup, \triangleright\})^*$
  - then
    - $\exists n \in \mathbb{N}$ , depending on  $M$ , such that there is no configuration  $C$  such that  $(s, \triangleright \sqcup w) \vdash_M^n C$ 
      - the length of any computation is bounded by  $n$
    - $(s, \triangleright \sqcup w) \vdash_M^* (h, \triangleright \sqcup v)$ ,  $h \in H$ ,  $f(w) = v$  for  $\forall$  computation
      - all computations agrees on the output

- Construct nondeterministic TM M to decide  
 $L = \{w \in \{0, 1\}^*: w \text{ is the binary encoding of a composite number}\}$  !
- Operation of M:
  - choose nondeterministically (write 1, 0, or stop) two binary numbers  $1 < p, q < w$ , and write them on the tape, after w, separated by ;
    - e.g.:  $\triangleright \underline{\_}110011;111;1111\underline{\_}\underline{\_}$

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- multiply p and q and put the answer, A, on the tape, in place of p and q
  - e.g.:  $\triangleright \underline{1}10011;1011111\underline{\hspace{2mm}}\underline{\hspace{2mm}}$
- compare A and w, if equal, go to state y, else go to state n

- Theorem: if a nondeterministic TM N semidecides or decides a language or computes a function  $\rightarrow \exists$  deterministic TM D semideciding, deciding the same language or computing the same function

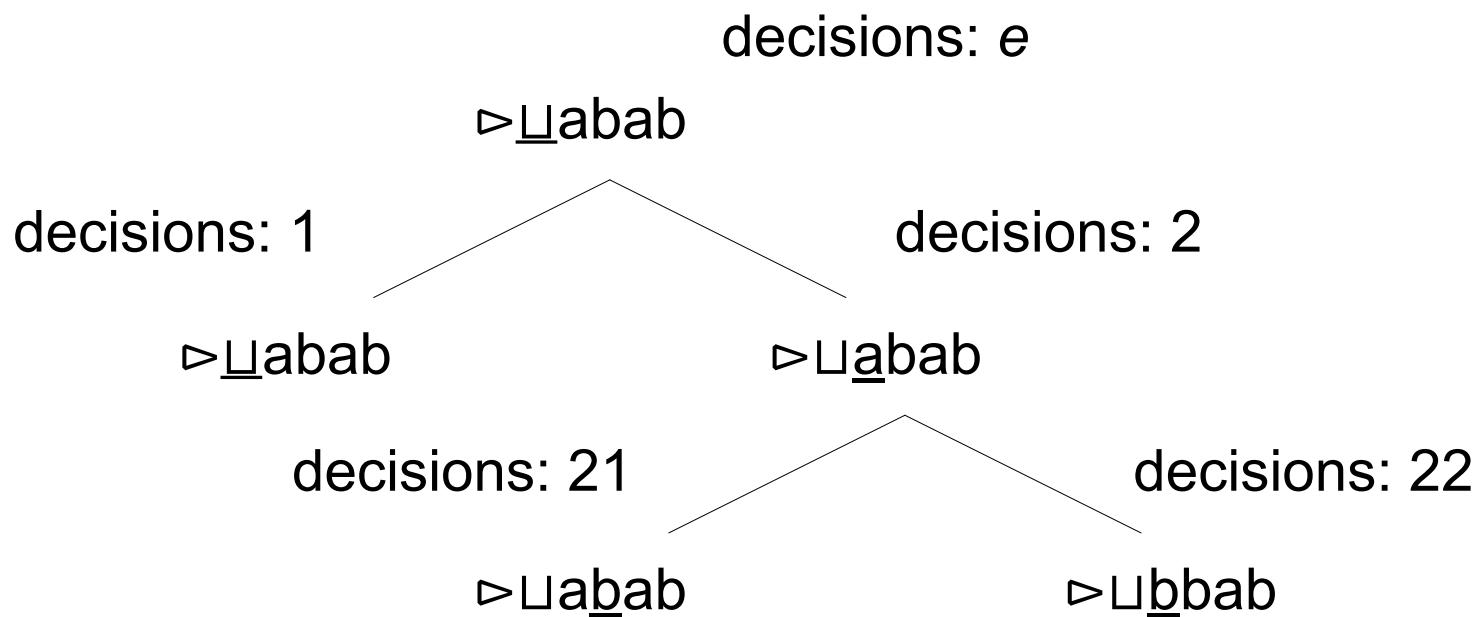
- Proof for semidecide:
  - construct D which simulates N by systematically trying out every possible computation of N
  - if such a computation is found which accepts the input then D also accepts the input

- Each configuration may lead to at most  $|K|^*(|\Sigma|+2)$  configurations
- The configurations can be organized into a computation tree
  - the root is the initial configuration

# Computation tree

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen



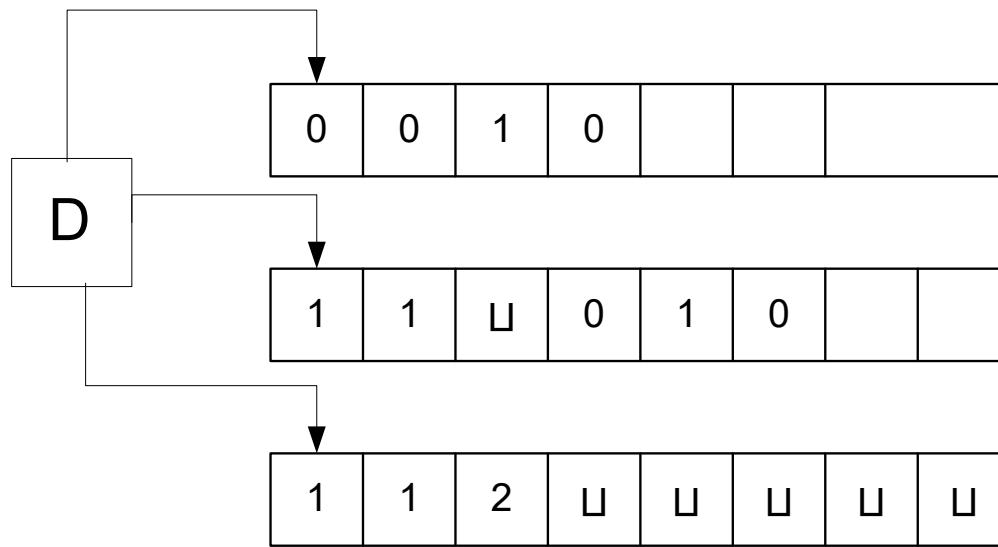
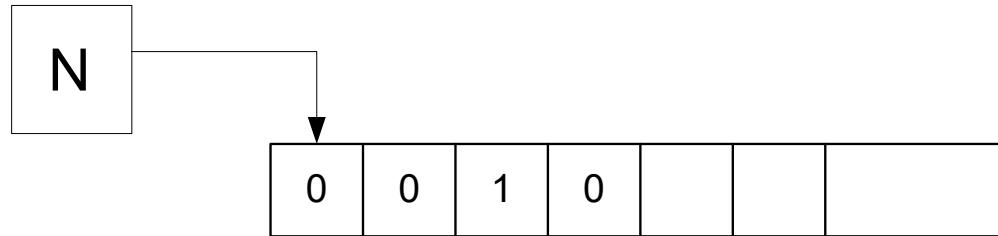
- A path of the computation tree of N represents a possible computation of D
  - a computation tree may be infinite for recursively enumerable languages but some branches are finite
  - a sequence of numbers correspond to each node representing the nondeterministic decisions made to reach that node

- one may create such a computation tree where every inner node has exactly  $r$  children
  - some children may be the same
- visiting the nodes corresponding to decisions  $e, 1, 2, \dots, r, 11, 12, \dots, 1r, 21, 22, \dots, rr, 111, \dots$  will encounter every node

# Nondeterministic TM $\leftrightarrow$ deterministic TM

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen



- D has three tapes:
  - tape 1: contains the input string and it is never altered
  - tape 2: a copy of N's tape, it changes according to a particular computation
  - tape 3: string of the actual non-deterministic decisions

- Operation:
  1. tape 1 contains the input w, tape 2, 3 are empty
  2. copy tape 1 to tape 2
  3. use tape 2 to simulate N with the help of tape 3
    - the current state, the input symbol, and tape 3 together determines the next configuration
      - string 142 in tape 3 means the 4<sup>th</sup> possible configuration has to be chosen (then go right)

- if an accepting configuration is encountered,  
accept the input
  - if no more symbols remain on tape 3 go to step  
4
4. prepare new computation: replace the string on  
tape 3 with the lexicographically next string, go to  
step 2 (e, 1, 2, ..., r, 11, 12, ... 1r, 21, 22, ..., rr, ...)

- Corollary: every nondeterministic TM has an equivalent deterministic TM

Recursively enumerable languages  
**Unrestricted grammars generating a language**  
Turing machines accepting a language

Recursive languages  
**Unrestricted grammars computing a function**  
Turing machines deciding a language  
 **$\mu$ -recursive functions**

Context-free languages  
Context-free grammars  
Pushdown automata

Regular languages  
Regular grammar  
Regular expressions  
Deterministic finite automata  
Nondeterministic finite automata

- TMs act as languages acceptors like weaker relatives:
  - the finite automata
  - pushdown automata
- TMs can recognize two important families of languages:
  - the recursive languages
  - recursively enumerable languages

- Language generators:
  - regular expressions
  - regular grammars
  - context-free grammars
  - new kind of language generator:  
(unrestricted) grammar

- Definition of (unrestricted) grammar: a quadtuple  $G = (V, \Sigma, R, S)$ , where
  - $V$  is an alphabet
  - $\Sigma \subseteq V$ , the set of terminal symbols
    - $V - \Sigma$ , the set of nonterminal symbols
  - $S \in V - \Sigma$ , is the start symbol
  - $R \subseteq (V^*(V - \Sigma)V^*) \times V^*$ , the rules
    - denote:  $u \rightarrow_G v$  if  $(u, v) \in R$

- Definition of the one step derivation of grammar G,  
 $\Rightarrow_G$ :
  - $u, v, x, y \in V^*$ ,  $u' \rightarrow_G v' \in R$
  - if  $u = xu'y$ ,  $v = xv'y \rightarrow u \Rightarrow_G v$
- When the grammar to which we refer to is obvious, we write  $A \rightarrow w$  and  $u \Rightarrow v$  instead of  $u' \rightarrow_G v'$  and  $u \Rightarrow_G v$
- Definition of derivation,  $\Rightarrow_G^*$ : the reflexive, transitive closure of  $\Rightarrow_G$

- We call the sequence  $w_0 \Rightarrow_G w_1 \Rightarrow_G \dots \Rightarrow_G w_n$  a derivation in  $G$  of  $w_n$  from  $w_0$ 
  - $w_0, w_1, \dots, w_n \in V^*$
  - $n \in \mathbb{N}$ , the length of the derivation
    - we say that the derivation has  $n$  steps
      - emphasizing the number of steps:  $w_0 \Rightarrow^n w_n$

- Definition of language generated by grammar G,  $L(G)$ : the set of strings generated by G
  - $L(G) = \{w \in \Sigma^* : S \Rightarrow_G^* w\}$
- Example: any context-free grammar is also a grammar
  - in CFG the right side of a rule is a member of  $V-\Sigma$ , rather than  $V^*(V-\Sigma)V^*$

- Give grammar G such that  $L(G) = \{a^n b^n c^n : n \geq 1\}$ !
  - $V = \{S, a, b, c, A, B, C, T_a, T_b, T_c\}$
  - $\Sigma = \{a, b, c\}$ , and
  - rules:
    - $S \rightarrow ABCS$
    - $S \rightarrow T_c$
    - $CA \rightarrow AC$
    - $BA \rightarrow AB$
    - $CB \rightarrow BC$

# Example

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $CT_c \rightarrow T_c C$
- $CT_c \rightarrow T_b C$
- $BT_b \rightarrow T_b b$
- $BT_b \rightarrow T_a b$
- $AT_a \rightarrow T_a a$
- $T_a \rightarrow e$

- Operation:
  - the first two rules generate string  $(ABC)^n T_c$
  - the next three rules allow to sort nonterminals A, B, C, so the string becomes  $A^n B^n C^n T_c$

- the remaining rules allow to migrate the token to the left while converting a nonterminal to a terminal one at a time
  - token: one of  $T_c$ ,  $T_b$ , or  $T_a$
  - while  $T_b$  migrates left, all  $B$  is transformed to  $b$ 
    - in the last step  $T_b$  is converted to  $T_a$
- Beware: some derivation does not yield a string because some nonterminals can not be eliminated

- $S \Rightarrow ABCS \Rightarrow ABCABCS \Rightarrow ABCABCT_c \Rightarrow ABACBCT_c \Rightarrow ABABCCT_c \Rightarrow AABBCCCT_c \Rightarrow AABBCT_c c \Rightarrow AABBT_b cc \Rightarrow AABT_b bcc \Rightarrow AABT_b bcc \Rightarrow AAT_a bbcc \Rightarrow AT_a abbcc \Rightarrow T_a aabbcc \Rightarrow aabbcc$
- Stuck:
  - $S \Rightarrow ABCS \Rightarrow ABCABCS \Rightarrow ABCABCT_c \Rightarrow ABCABT_b c \Rightarrow ABCAT_a bc \Rightarrow ABCT_a abc \Rightarrow ABCabc$

- Theorem: a language is generated by a grammar  $\leftrightarrow$  it is recursively enumerable
- Proof: only if →
  - $G = (V, \Sigma, R, S)$  is given, construct TM M such that  $L(M) = L(G)$

- M is nondeterministic and has 2 tapes:
  - tape 1: contains the input w, never changes
  - tape 2: M reconstructs a derivation of w from S
    - M starts writing S on tape 2
    - M selects nondeterministically a rule of G or a comparison of tape 1 and 2

- for each rule  $u' \rightarrow_G v'$ , M has a number of transitions:
  - » nondeterministically find  $u'$
  - » delete  $u'$
  - » shift the remaining part of the string to make space
  - » write  $v'$
- go back to select a new rule or to compare
- M stops only if the comparison of tape 1 and 2 results in a positive answer

- Proof: if  $\leftarrow$ 
  - TM M is given, construct  $G = (V, \Sigma, R, S)$  such that  $L(M) = L(G)$ 
    - suppose M deletes the tape before accepting
  - $V = \Sigma \cup K$
  - G simulates a backward computation of M
  - configuration  $(q, \triangleright u \underline{a} w)$  is represented by the string  $\triangleright u a q w \triangleleft$ 
    - $\triangleleft$  is the new end marker symbol

- rules have to be introduced to each transition of M:
  - if  $\delta(q, a) = (p, b)$ ,  $p \in K$ ,  $b \in \Sigma$ 
    - $(q, \triangleright u \underline{a} w) \vdash (p, \triangleright u \underline{b} w)$
    - $\triangleright u b p w \triangleleft \Rightarrow \triangleright u a q w \triangleleft$
    - G rule:  $bp \rightarrow aq$ , (p, q are nonterminals)

- if  $\delta(q, a) = (p, \rightarrow)$ ,  $p \in K$ 
  - $(q, \triangleright u \underline{a} b w) \vdash (p, \triangleright u a \underline{b} w)$
  - $\triangleright u a b p w \triangleleft \Rightarrow \triangleright u a q b w \triangleleft$
  - G rule:  $abp \rightarrow aqb$ , for all b

- if  $\delta(q, a) = (p, \rightarrow)$ ,  $p \in K$ 
  - $(q, \triangleright u \underline{a}) \vdash (p, \triangleright u a \underline{\sqcup})$
  - $\triangleright u a \sqcup p \triangleleft \Rightarrow \triangleright u a q \triangleleft$
  - G rule:  $a \sqcup p \triangleleft \rightarrow a q \triangleleft$

- if  $\delta(q, a) = (p, \leftarrow)$ ,  $p \in K$ ,  $a \neq \sqcup$  or  $w \neq e$ 
  - $(q, \triangleright ub\cancel{a}w) \vdash (p, \triangleright ub\cancel{a}w)$
  - $\triangleright ubpaw \triangleleft \Rightarrow \triangleright ubaqw \triangleleft$
  - G rule:  $pa \rightarrow aq$ , for all b
- if  $\delta(q, \sqcup) = (p, \leftarrow)$ ,  $p \in K$ 
  - $(q, \triangleright ub\sqcup) \vdash (p, \triangleright ub\cancel{b})$
  - $\triangleright ubp \triangleleft \Rightarrow \triangleright ub\sqcup q \triangleleft$
  - G rule:  $p \triangleleft \rightarrow \sqcup q \triangleleft$
- other rules:  $S \rightarrow \triangleright \sqcup h \triangleleft$ ,  $\triangleright \sqcup S \rightarrow e$ ,  $\triangleleft \rightarrow e$
- a derivation:  $S \Rightarrow \triangleright \sqcup h \triangleleft \Rightarrow^* \triangleright \sqcup sw \triangleleft \Rightarrow w \triangleleft \Rightarrow w$

- Definition of computation of function  $f$  by grammar  $G$ :
  - $\forall w, v \in \Sigma^*, SwS \Rightarrow^* v \leftrightarrow v = f(w)$
  - we say  $G$  computes  $f: \Sigma^* \rightarrow \Sigma^*$
  - $G$  yield exactly one string, the correct value of  $f(w)$  from  $SwS$ 
    - all derivations agree on the result
    - note that a CFG cannot change terminals
- The computation of  $f$  by  $G$  is always finite so it can be considered as an algorithm

- Definition of grammatically computable function  $f$ :  $\exists$  such grammar  $G$  that computes it
- Theorem: a function is recursive (Turing computable)  $\leftrightarrow$  it is grammatically computable
- Proof: similar to previous theorem



# Introduction to the Theory of Computation

Lesson 8

## 4.7. Numerical functions

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MAGYARORSZÁG  
KORMÁNYA

**Európai Unió**  
Európai Strukturális  
és Beruházási Alapok



**BEFEKTETÉS A JÖVŐBE**

- Functions from numbers to numbers
- They are a new representation of computation different from both TMs and unrestricted grammars

- E.g.:  $f(m, n) = mn^2 + 3m^{2m+17}$ 
  - we know that  $f$  can be calculated as it is the composition of addition, multiplication, and exponentiation
  - exponentiation can be defined recursively with multiplication
  - multiplication can be defined recursively with addition

- Definition of the basic functions:
  - k-ary zero function:  $0 \leq k$ ,  $\text{zero}_k(n_1, \dots, n_k) = 0$ ,  
 $n_1, \dots, n_k \in \mathbb{N}$
  - j-th k-ary identity function:  $0 < j \leq k$   $\text{id}_{k,j}(n_1, \dots, n_k) = n_j$ ,  $n_1, \dots, n_k \in \mathbb{N}$ 
    - we may omit subscript k from id and zero if it does not cause confusion
  - successor function:  $\text{succ}(n) = n + 1$ ,  $n \in \mathbb{N}$
- E.g.:
  - $\text{zero}_2(3, 5) = 0$ ,  $\text{id}_{3,1}(6, 2, 0) = 6$ ,  $\text{succ}(34) = 35$

- Definition of composition:
  - let
    - $0 \leq k, l$
    - $g: N^k \rightarrow N$
    - $f, h_1, \dots, h_k: N^l \rightarrow N$
  - the composition of  $g$  with  $h_1, \dots, h_k$ :  
 $f(n_1, \dots, n_l) = g(h_1(n_1, \dots, n_l), \dots, h_k(n_1, \dots, n_l))$ 
    - the 1<sup>st</sup> parameter of  $g$  is given by  $h_1$ , the 2<sup>nd</sup> is given by  $h_2, \dots$
  - if  $k = 1 \rightarrow f(n_1, \dots, n_l) = g(h(n_1, \dots, n_l))$

- Definition of recursion:

- let

- $0 \leq k$
    - $g: N^k \rightarrow N$
    - $f: N^{k+1} \rightarrow N$
    - $h: N^{k+2} \rightarrow N$

- function  $f$  defined recursively by  $g$  and  $h$ :

- $f(n_1, \dots, n_k, 0) = g(n_1, \dots, n_k)$ , base case
    - $f(n_1, \dots, n_k, m+1) = h(n_1, \dots, n_k, m, f(n_1, \dots, n_k, m))$   
 $n_1, \dots, n_k, m \in N$

- Definition of primitive recursive functions: the basic functions and the functions that can be obtained from them with composition and recursion

- Show that  $\text{plus2}(n) = n + 2$  is primitive recursive!
  - $\text{plus2}(n) = \text{succ}(\text{succ}(n))$
  - obtained by composition  $k = l = 1$  using  $f = \text{plus2}$ ,  
 $g = h_1 = \text{succ}$

- Show that  $\text{plus}(n, m) = n + m$  is primitive recursive!
  - $\text{plus}(n, 0) = \text{id}_1(n) = n$
  - $\text{plus}(n, m+1) = \text{succ}(\text{plus}(n, m)) =$   
 $\text{succ}(\text{id}_3(n, m, \text{plus}(n, m)))$ 
    - e.g.:  $\text{plus}(4, 2) = \text{plus}(4, 1) + 1$
  - $f$  is obtained by primitive recursion:  $k = 1$ ,  $f = \text{plus}$ ,  
 $g = \text{id}_1$ ,  $h(n, m, p) = \text{succ}(\text{id}_3(n, m, p))$
  - $h$  is primitive recursion obtained by composition

- Multiplication  $\text{mult}(n, m) = n*m$  is defined recursively:
  - $\text{mult}(n, 0) = \text{zero}(n) = 0$
  - $\text{mult}(n, m + 1) = \text{plus}(n, \text{mult}(n, m))$   
 $= \text{plus}(\text{id}_1(n, m, \text{mult}(n, m)), \text{id}_3(n, m, \text{mult}(n, m)))$

- Exponent function  $\exp(n, m) = n^m$  is defined as:
  - $\exp(n, 0) = \text{succ}(\text{zero}(n)) = 1$
  - $\exp(n, m + 1) = \text{mult}(n, \exp(n, m))$

- Constant functions such as  $f(n_1, \dots, n_k) = 3$  are primitive recursive
  - $f(n_1, \dots, n_k) = \text{succ}(\text{succ}(\text{succ}(\text{zero}(n_1, \dots, n_k))))$
- From now on we can use the usual arithmetic notation:  
 $x^*y+3$  instead of  $\text{plus}(\text{mult}(x, y), 3)$

- The sign function,  $\text{sgn}(m)$  is zero if  $m = 0$  otherwise, it is 1 is also primitive recursive, defined with recursion
  - $\text{sgn}(0) = 0 = \text{zero}()$ ,  $k = 0$
  - $\text{sgn}(m + 1) = 1 = \text{succ}(\text{zero}(m, \text{sgn}(m)))$

- The pred function,  $\text{pred}(m) = m-1$  if  $m > 0$ , 0 otherwise, is also primitive recursive, defined with recursion
  - $\text{pred}(0) = 0$
  - $\text{pred}(m+1) = m$

- The nonnegative subtraction function,  $n \sim m = \max \{n-m, 0\}$ , is also primitive recursive, defined with recursion:
  - $n \sim 0 = n$
  - $n \sim (m + 1) = \text{pred}(n \sim m)$

- Definition of primitive recursive predicate: such a primitive recursive function which takes only values 1, 0
- The iszero predicate ( $1$  if  $n = 0$ ,  $0$  otherwise) is also primitive recursive, defined with recursion
  - $\text{iszzero}(0) = 1$
  - $\text{iszzero}(n + 1) = 0$

- The isone predicate, (1 if  $n = 1$ , 0 otherwise) is also primitive recursive, defined with recursion
  - $\text{isone}(0) = 0$
  - $\text{isone}(n + 1) = \text{iszzero}(n)$
- The predicate  $\text{positive}(n)$  is the same as the  $\text{sgn}(n)$

- The predicate greater-than-or-equal( $m, n$ ):  $m \geq n$  is defined as  $\text{iszzero}(n \sim m)$
- The predicate less-than( $m, n$ ) is defined as:  
 $1 \sim \text{greater-than-or-equal}(m, n)$ 
  - the negation of any primitive recursive predicate is also primitive recursive

- The disjunction of two primitive recursive predicates,  $p(m, n)$ ,  $q(m, n)$  is also primitive recursive
  - $p(m, n)$  or  $q(m, n)$  defined as:  
 $1 \sim \text{iszzero}(p(m, n) + q(m, n))$
  - e.g.:  $p(m, n) = 1$ ,  $q(m, n) = 0 \rightarrow \text{iszzero}(p(m, n) + q(m, n)) = 0$ ,  $1 \sim 0 = 1 = 1$  or 0

- The conjunction of two primitive recursive predicates,  $p(m, n)$ ,  $q(m, n)$  is also primitive recursive
  - $p(m, n)$  and  $q(m, n)$  defined as:  
 $p(m, n) * q(m, n)$

- If  $g$  and  $h$  are primitive recursive and  $p$  is a primitive recursive predicate, all three with the same arity  $k \rightarrow$  the function defined by cases is also primitive recursive

$$f(n_1, \dots, n_k) = \begin{cases} g(n_1, \dots, n_k), & \text{if } p(n_1, \dots, n_k) \\ h(n_1, \dots, n_k), & \text{otherwise} \end{cases}$$

$$- f(n_1, \dots, n_k) = p(n_1, \dots, n_k) * g(n_1, \dots, n_k) + \\ (1 \sim p(n_1, \dots, n_k)) * h(n_1, \dots, n_k)$$

- The remainder function is also primitive recursive defined with recursion
  - $62 \% 6 = 61 \% 6 + 1 = (60 \% 6 + 1) + 1$
  - $10^*6 \% 6 = 0 \leftrightarrow (10^*6-1) \% 6 = 6-1$ 
    - $x^*y \% y = 0 \leftrightarrow (x^*y-1) \% y = y-1$
  - $\text{rem}(0, n) = 0$
  - $\text{rem}(m+1, n) = \begin{cases} 0, & \text{if equal}(\text{rem}(m, n), \text{pred}(n)) \\ \text{rem}(m, n)+1, & \text{otherwise} \end{cases}$

- The divide function is also primitive recursive defined with recursion
  - $65 / 6 = 64 / 6 = 63 / 6 = 62 / 6 = 61 / 6 = 60 / 6$
  - $(65 + 1) / 6 = 65 / 6 + 1$ ,  $\text{rem}(65, 6) = \text{pred}(6)$
  - $\text{div}(0, n) = 0$
  - $\text{div}(m+1, n) = \begin{cases} \text{div}(m, n)+1, & \text{if } \text{equal}(\text{rem}(m,n), \text{pred}(n)) \\ \text{div}(m,n), & \text{otherwise} \end{cases}$

- Theorem: there are functions which are not primitive recursive

- Proof:
  - primitive recursive functions are countable
    - they can be represented as strings over a finite alphabet
    - the alphabet contains symbols for
      - basic functions
      - recursion and composition
      - 0 and 1 are used to index in binary the basic functions, e.g.:  $d_{11,1}$
      - parenthesis, coma

- consider the strings which define primitive recursive functions and list them all in lexicographic order:  
 $f_0, f_1, \dots$
- let  $g(n) = f_n(n)+1$  (Diagonalization)
- suppose  $g(n)$  is primitive recursive
  - $\exists m \in \mathbb{N}$  such that  $f_m = g$
  - $f_m(m) = g(m) = f_m(m)+1$  which is contradiction
- $g$  is not primitive recursive

- Definition of minimalization of  $g$ :

- $0 \leq k$ ,  $g: N^{k+1} \rightarrow N$

- $f(n_1, \dots, n_k) = \begin{cases} \text{the least } m, \text{ such that } g(n_1, \dots, n_k, m) = 1, \\ \text{if such an } m \text{ exists} \\ 0, \text{ otherwise} \end{cases}$

- $f$  can be denoted:  $\mu m [g(n_1, \dots, n_k, m) = 1]$

- $\mu m$  - minimalization in  $m$

- though  $f$  is always well-defined there is no obvious method for computing it

- The obvious method to calculate  $m$  is not algorithm, because it may fail to terminate

```
m := 0;
```

```
while g(n1, ..., nk, m) ≠ 1 do m := m + 1;  
output m
```

- Definition of a minimalizable function  $g$ :
  - $0 \leq k$ ,  $g: \mathbb{N}^{k+1} \rightarrow \mathbb{N}$ ,  $n_1, \dots, n_k \in \mathbb{N}$ ,  $\exists m \in \mathbb{N}$  such that  $g(n_1, \dots, n_k, m) = 1$
  - if the above method always terminates

- Definition of  $\mu$ -recursive function: the basic functions and the functions that can be obtained from them with composition, recursion, and minimalization of minimalizable functions

- Define the logarithm function:  $\log(m, n)!$ 
  - to avoid the mathematical pitfalls the function  $\log$  is defined as:  $\log(m, n) = \lfloor \log_{m+2}(n+1) \rfloor$
  - $\log(m, n) =$   
 $\mu p[\text{greater-than-or-equal}((m + 2)^p, n + 1)]$ 
    - $\exists p \geq 0$  such that  $(m + 2)^p \geq n$
    - $\log(10, 1000) = 3$
  - $\log$  can be also defined without minimalization

- Theorem: a function  $f: N^k \rightarrow N$  is  $\mu$ -recursive  $\leftrightarrow$  it is recursive (Turing computable)

- Proof by induction: only if →
  - suppose  $f$  is  $\mu$ -recursive
  - construct random access TM  $M$  to compute  $f$
  - the basic functions are recursive:
    - zero function: delete input, write 0
    - identity function: delete part of the input, shift the remaining part to the left
    - successor function: if the number is represented unary → write 1, step right

- composition can be computed by M
  - $f: N^k \rightarrow N$  is the composition of the function  $g: N^l \rightarrow N$  with  $h_1, \dots, h_l: N^k \rightarrow N$ 
    - by induction we know that  $g$  and  $h$  functions are computable

- $f$  is computed as

$$m_1 := h_1(n_1, \dots, n_k);$$
$$m_2 := h_2(n_1, \dots, n_k);$$
$$m_3 := h_3(n_1, \dots, n_k);$$

...

$$m_l := h_l(n_1, \dots, n_k);$$

output  $g(m_1, \dots, m_l);$

- recursion can be carried out by M
  - f is defined recursively from g and h
    - by induction we know that g and h functions are computable

- $f(n_1, \dots, n_k, m)$  can be computed:

```
v := g(n1, ..., nk);  
if m = 0 then output v;  
else  
    for i := 1, 2, ..., m do  
        v := h(n1, ..., nk, i - 1, v);  
    output v
```

- minimization can be carried out by M
  - f is by the minimization of g
    - g is minimalizable
    - by induction we know that g is computable
  - $f(n_1, \dots, n_k)$  can be computed:  
 $m := 0;$   
while  $g(n_1, \dots, n_k, m) \neq 1$  do  $m := m + 1;$   
output  $m$
  - the three basic functions and the three operations  
can be computed with M
- Proof: if  $\leftarrow$  (in the book)

- Lambda calculus
- Combinatory logic
- Markov algorithm
- Register machine
- P"
- Source:  
[http://en.wikipedia.org/wiki/Theory\\_of\\_computation](http://en.wikipedia.org/wiki/Theory_of_computation)



# Introduction to the Theory of Computation

## Lesson 9

- 5.4. Unsolvable problems about Turing machines
- 5.5. Unsolvable problems about grammars
- 5.6. An unsolvable tiling problem
- 5.7. Properties of recursive languages

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MAGYARORSZÁG  
KORMÁNYA

**Európai Unió**  
Európai Strukturális  
és Beruházási Alapok



**BEFEKTETÉS A JÖVŐBE**

- Definition of:
  - recursive languages
  - recursively enumerable languages
  - recursive functions
  - universal Turing machine
  - halting problem
  - problems and the corresponding language

- Definition of reduction  $\tau$ , from language  $L_1$  to  $L_2$ : a recursive function (Turing computable)  $\tau: \Sigma^* \rightarrow \Sigma^*$  such that  $x \in L_1 \leftrightarrow \tau(x) \in L_2$ 
  - $L_2$  can contain more strings than  $L_1$ , strings which are not the image of any string in  $L_1$
  - the term reduction from  $L_1$  to  $L_2$  is misleading

- Theorem: if  $L_1$  is not recursive,  $\exists$  a reduction from  $L_1$  to  $L_2$  then  $L_2$  is not recursive either
  - a non-recursive language corresponds to an undecidable problem,  $L_1$  to  $P_1$ ,  $L_2$  to  $P_2$
  - $P_2$  is at least as complex as  $P_1$ ,  $P_1 \leq P_2$

- Proof by contradiction:
  - suppose  $L_2$  is recursive, so it is decided by Turing machine  $M_2$
  - let  $T$  be the Turing machine, that computes the reduction  $\tau$
  - the machine schema  $TM_2$  would decide  $L_1$
  - $L_1$  is not recursive, we have reached a contradiction

# Undecidable problems about TM

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: the following problem is undecidable: given a TM M, does M halt on the empty tape?

- Proof:
  - the corresponding language:  $L = \{\rho(M) : M \text{ halts on } e\}$
  - let  $M_w$  write  $w$  on its tape then starts to simulate  $M$ 
    - if  $w = a_1 \dots a_n \rightarrow M_w = Ra_1Ra_2R\dots Ra_nL \sqcup M$
  - $M$  halts on  $w \leftrightarrow M_w$  halts on  $e$ 
    - $\rho(M)\rho(w) \in H \leftrightarrow \rho(M_w) \in L$
  - there is  $\tau: \rho(M)\rho(w) \rightarrow \rho(M_w)$  a reduction from  $H$  to  $L$ 
    - it is obvious that  $\tau$  exists though we do not define it
      - it deletes  $\rho(w)$  and modify  $\rho(M)$  to  $\rho(M_w)$
  - $\tau$  exists,  $H$  is not recursive  $\rightarrow L$  is not recursive

- Theorem: the following problem is undecidable: given a TM M
  - is there any string on which M halts?
  - does M halts on every input?

- Proof:
  - first language:  $L' = \{\rho(M) : M \text{ halts on some input}\}$
  - second language:  $E = \{\rho(M) : M \text{ halts on every input}\}$
  - let  $M'$  which erases the input and simulate  $M$ ,  $M' = DM$

- $M \text{ halts on } e \leftrightarrow M' = DM \text{ halts on any / each } w$ 
  - $\rho(M) \in L \leftrightarrow \rho(M') \in L' = E$
  - $L = \{\rho(M) : M \text{ halts on } e\}$
- there is  $\tau$ :  $\rho(M) \rightarrow \rho(M')$  a reduction from  $L$  to  $L'$  and also from  $L$  to  $E$
- $\tau$  exists,  $L$  is not recursive  $\rightarrow L'$  and  $E$  are not recursive

# Undecidable problems about TM

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: the following problem is undecidable: given TM  $M_1$  and  $M_2$ , do they halt on the same input strings?

- Proof:
  - the corresponding language:  $S = \{\rho(M_1)\rho(M_2) : M_1$  and  $M_2$  halt on the same input}
  - let  $M_2 = y$  to be the machine that halts on every input
  - $M$  halts on every input  $\leftrightarrow$  both  $M$  and  $y$  halt on every input, thus, they halt on the same input
    - $\rho(M) \in E \leftrightarrow \rho(M)\rho(y) \in S$
  - there is  $\tau: \rho(M) \rightarrow \rho(M)\rho(y)$  a reduction from  $E$  to  $S$ 
    - $\tau$  exists though we do not define it
  - $\tau$  exists,  $E$  is not recursive  $\rightarrow S$  is not recursive

- Theorem:  $\exists$  a fixed TM M for which the following problem is undecidable: does M halt on a given input string?
- Proof: M is the universal TM

- Theorem: the following problems are undecidable
  - given grammar  $G$  and string  $w$ , determine if  $w \in L(G)$
  - given grammar  $G$ , determine if  $\epsilon \in L(G)$
  - given two grammars  $G_1, G_2$ , determine if  $L(G_1) = L(G_2)$
  - given grammar  $G$ , determine if  $L(G) = \emptyset$
  - $\exists$  fixed grammar  $G_0$ , determine if string  $w \in L(G_0)$

- Proof:
  - a reduction from the halting problem exists to (a) based on TM unrestricted grammar equivalence
  - simple reductions exist from (a) to the other cases
    - see hints at undecidable problems about TMs

- Theorem: the following problems about CFG are undecidable
  - a) given CFG  $G$ , determine if  $L(G) = \Sigma^*$
  - b) given CFG  $G_1, G_2$ , determine if  $L(G_1) = L(G_2)$
  - c) given PDA  $M_1, M_2$ , determine if  $L(M_1) = L(M_2)$
  - d) given PDA  $M$ , find an equivalent PDA with as few states as possible

- Proof a: in the book
- Proof b:
  - assume problem b) is decidable
  - consider CFG  $G_{\text{triv}}$  that generates  $\Sigma^*$
  - by checking if  $L(G) = L(G_{\text{triv}})$  we can decide if  $L(G) = \Sigma^*$  (problem a) which is contradiction
- Proof c:
  - assume problem c) is decidable
  - $\exists$  CFG  $G_1, G_2$  such that  $L(G_1) = L(M_1), L(G_2) = L(M_2)$
  - by checking if  $L(M_1) = L(M_2)$  we can decide if  $L(G_1) = L(G_2)$  (problem b) which is a contradiction

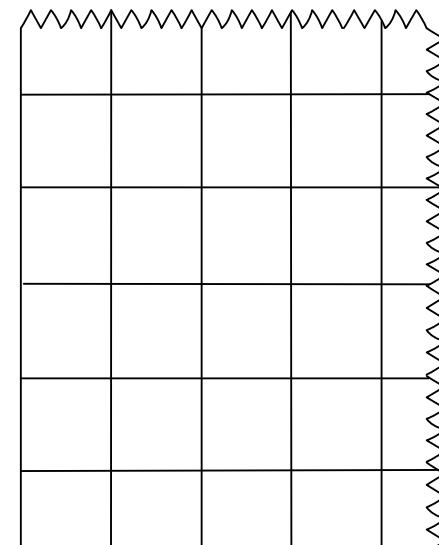
- Lemma: it is decidable if a one state PDA accepts  $\Sigma^*$
- Proof d:
  - assume we can optimize PDA M
  - let M' the optimized PDA
  - $\exists$  CFG G such that  $L(G) = L(M)$
  - by checking if M' has one state and  $L(M') = \Sigma^*$  we can decide if  $L(G) = \Sigma^*$  (problem a) which is a contradiction

# An unsolvable tiling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of the tiling problem: given different types of tiles, each one is a unit square, there is an infinite supply of copies of each type; tile the first quadrant of the plane
  - the origin tile must be placed in the corner
  - given which types may adjoin horizontally
  - given which types may adjoin vertically
  - tiles may not be rotated or turned over



# An unsolvable tiling problem

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of a tiling system: a quadruple  $\Delta = (D, d_0, H, V)$ , where
  - $D$  set of the types of the tiles
  - $d_0 \in D$ , the origin tile
  - $H \subseteq D^2$  set of those pairs whose 1st and 2nd component may adjoin horizontally
  - $V \subseteq D^2$  set of those pairs whose 1st and 2nd component may adjoin vertically

# An unsolvable tiling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of tiling by  $\Delta$ :  $f: \mathbb{N}^2 \rightarrow D$  such that
  - $f(0, 0) = d_0$
  - $(f(m, n), f(m+1, n)) \in H, \forall m, n \in \mathbb{N}$
  - $(f(m, n), f(m, n+1)) \in V, \forall m, n \in \mathbb{N}$

# An unsolvable tiling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: the following problem is undecidable: given a tiling system, determine if  $\exists$  a tiling by that system

# An unsolvable tiling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Proof idea:
  - reduce the complement of the halting problem to tiling problem
    - given TM M, determine if M fails to halt on input e
      - it is also undecidable
  - construct from any TM M a tiling system
    - the horizontal dimension represents the tape of M
    - the vertical dimension stands for time

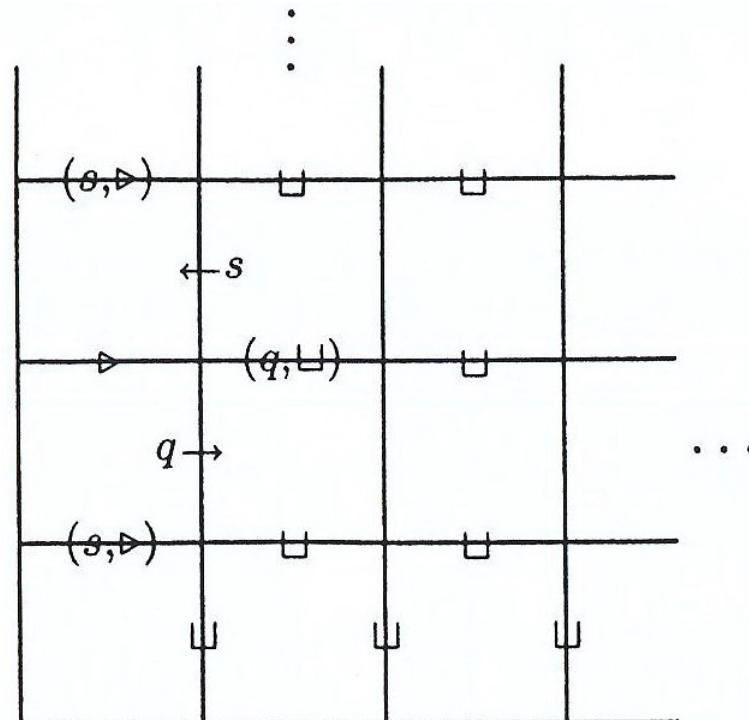
- if  $M$  never halts on the empty input, successive rows can be tiled ad infinitum
  - if  $M$  halts after  $k$  steps → it is impossible to tile more than  $k$  rows
- to ensure the relations  $H$  and  $V$  think of the edges of the tiles as being marked with certain information
  - tiles are allowed to adjoin each other only if the markings on the adjacent edges are identical
  - on the horizontal edges markings are either a member of  $\Sigma$  or  $K \times \Sigma$
  - on the vertical edges markings are absent or a member of  $K$  with a directional indicator

# An unsolvable tiling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- the markings on the horizontal edges between the  $n^{\text{th}}$  and  $(n+1)^{\text{st}}$  rows of tiles give the configuration of  $M$  after  $n-1$  steps of its computation



- Every TM  $M$  semidecides a unique language  $L(M)$ 
  - there are infinite many other TMs with the same language
    - e.g.: renumber the states of  $M$  or make some extra unnecessary head movements
- The complexity of this mapping between TMs and languages implies the next theorem

- Theorem (Rice's Theorem): suppose that  $C$  is a proper subset of the class of all recursively enumerable languages, the following problem is undecidable:
  - given a TM  $M$ , is  $L(M) \in C$ ?

- Corollary: the next problems are undecidable:
  - given a TM M, is  $L(M)$ 
    - finite?
    - regular?
    - context free?
    - recursive?
    - empty?
    - $\Sigma^*$ ?

- Every recursive language is recursively enumerable, but the two classes are not the same
  - language H (the halting problem) is recursively enumerable but not recursive

- Definition of Turing enumerable language  $L: \exists \text{ TM } M$  such that for some fixed state  $q \in K$  (not a halting state),

$$L = \{w : (s, \triangleright \underline{w}) \vdash_M^* (q, \triangleright \underline{w}w)\}$$

- we say  $M$  enumerates  $L$ 
  - $M$  starts from the blank tape

- after some computation M reaches state q
  - the current content of the tape is a member of L
  - we say M displays a string
- the computation continues, upon reaching q again a new string in L is generated again, and so on
- M may enumerate the strings of L in any order, possibly with repetitions

- Theorem: a language is recursively enumerable  $\leftrightarrow$  it is Turing enumerable

- Proof: if,  $\leftarrow$ 
  - suppose  $M$  enumerates  $L$
  - construct a 2-tape TM  $M'$  such that
    - $M'$  saves its input on tape 2
    - simulate  $M$  on tape 1
    - whenever the fixed state of  $M$  arises the two tapes are compared, if the tapes are
      - the same:  $M'$  halts
      - different:  $M'$  continues to simulate  $M$  on tape 1
  - $M'$  semidecides  $L$

- Definition of lexicographically Turing enumerable language  $L: \exists \text{ TM } M \text{ such that } M \text{ enumerates } L \text{ and if } (q, \triangleright \underline{w}) \vdash_M^+ (q, \triangleright \underline{w}'), q \text{ is the fixed state} \rightarrow w' \text{ comes lexicographically after } w$ 
  - we say  $M$  lexicographically enumerates  $L$ 
    - $M$  generates the strings of  $L$  in lexicographical order

- Theorem: a language is recursive  $\leftrightarrow$  it is lexicographically Turing enumerable

- Proof: only if →
  - suppose M decides L
  - construct M' such that
    - it generates one after the other all strings in the alphabet of L in lexicographic order
    - simulate M on the actual string
      - if accepted → display it by going into the fixed state
      - if rejected → continue with the next string
  - M' lexicographically enumerates L

- Proof: if  $\leftarrow$ 
  - suppose  $M$  lexicographically enumerates  $L$
  - if  $L$  finite  $\rightarrow L$  regular  $\rightarrow L$  recursive
    - $L = \{w_1 \cup w_2 \cup \dots\}$
  - suppose  $L$  is infinite
  - construct  $M'$  which works on  $w$  input such that
    - simulate  $M$  and check each displayed string  $x$ 
      - if  $x = w \rightarrow$  accept  $w$
      - if  $x < w \rightarrow$  go on simulating  $M$
      - if  $x > w \rightarrow$  reject  $w$
  - $M'$  decides  $L$  as there are finitely many  $x$ 's such that  $x < w$



# Introduction to the Theory of Computation

Lesson 10

6.1. The class P

6.2. Problems, problems

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- Definition of:
  - algorithm complexity
  - order of  $f$ ,  $O(f)$
  - relation  $f \approx g$
  - reflexive, transitive closure

- Classification of problems, those
  - which can be solved by algorithms
  - which can not be solved by algorithms
- Many problems in the first group can not be solved in any practical sense by computers because of the excessive time requirements

- Definition of the traveling salesman problem, TSP:  
given  $n$  cities and the distances in miles between them,  
minimize the total distance traversed to visit all the  
cities exactly once
  - $d_{i,j}$ : the distance between city  $i$  and  $j$ 
    - $d_{i,i} = 0$
    - $d_{i,j} = d_{j,i}$
  - a tour is defined by function  $\pi$ ,  $\pi(i)$  is the  $i$ -th city in  
the tour
  - the cost of a tour:  $c(\pi) = d_{\pi(1),\pi(2)} + d_{\pi(2),\pi(3)} + \dots$

- The number of possible itineraries:  $(n-1)!$ 
  - first  $n-1$  city can be selected, then  $n-2$ , ...
  - an algorithm can be easily designed that systematically examines all itineraries in order to find the shortest
    - the problem is solvable

- suppose there are 40 cities
  - number of itineraries:  $39!$ , larger than  $10^{45}$ 
    - if we could examine  $10^{15}$  tours/s which is quicker than any available or planned computer
    - required time: several billion lifetime of the universe

- The Church-Turing thesis defines something as algorithm if it can be represented by such a TM which halts on every input
- Similarly the practically feasible algorithm should be defined
- The algorithms (and the solvable problems) can be divided into two subgroups
  - practically feasible algorithms
  - practically infeasible algorithms

- The difficulty in the traveling salesman problem is the steep growth of the time or steps required by the algorithm on a given input
  - input: number of cities
  - the number of steps required has an exponential rate of growth
    - it is proportional with the number of itineraries  $(n-1)!$
    - grows even faster than  $2^n$
- If the number of steps required by an algorithm has a polynomial rate of growth then we tend to call it practically feasible

- It must be decided which computational device should be used as the basis of the practically feasible algorithm
  - TM is the obvious choice but it has a lot of variants
    - multi-tape
    - random access
    - ...

- these are equally adequate as they can be simulated by a classic TM which requires polynomial time
  - exception: nondeterministic TM
    - its simulation requires exponential time as each possible computation must be examined

- Definition of polynomially bounded deterministic TM M:  
 $\exists$  a polynomial  $p(|x|)$  such that there is no configuration C for any input x for which  $(s, > \underline{u}x) \vdash_M^{p(|x|)+1} C$   
– M always halts at most  $p(|x|)$  steps
- Definition of polynomially decidable language:  $\exists$  a polynomially bounded TM that decides it

- Definition of P: the class of all polynomially decidable languages
  - the quantitative analog of the class of recursive languages
- We regard something as a practically feasible algorithm if it belongs to P
- See: [http://qwiki.stanford.edu/wiki/Complexity\\_Zoo](http://qwiki.stanford.edu/wiki/Complexity_Zoo) for hierarchy of the problem classes

- Theorem: P is closed under complement
- Proof:
  - if language L is decidable by a polynomially bounded TM M
  - then  $L^C$  is decided by the version of M that inverts y and n states
    - the polynomial bound is unaffected

- Definition of language  $E: \{\rho(M)\rho(w) : M \text{ accepts input } w \text{ after at most } 2^{|w|} \text{ steps}\}$ 
  - the quantitative analog of the halting language  $H$
  - $E$  is recursive: a modified version of the universal TM,  $U$  can decide it
    - if after at most  $2^{|w|}$  steps  $U$  does not halt then  $w$  is rejected

- Theorem:  $E \notin P$ 
  - $E$  can be decided but only with a practically infeasible algorithm, i.e., with a TM which is not bounded polynomially based on the length of the input

- How well does class P capture the intuitive notion of practically feasible algorithm?
  - an algorithm with time
    - $n^{100}$  or  $10^{100}n^2$  belongs to P
    - $n^{\log \log n}$  does not belong to P
  - $\exists n_0$  such that for  $\forall n > n_0$ ,  $n^{100} < n^{\log \log n}$ 
    - but  $n_0$  is very big in this case
    - if in practice inputs smaller than  $n_0$  are important then  $n^{\log \log n}$  is preferable despite that it is not in P
- Example: simplex method for LP is exponential but widely used in practice

- Though theoretically possible, extreme time bounds rarely come up in practice
- In general
  - polynomial algorithms in computational practice have:
    - small exponents
    - small constant coefficients
  - nonpolynomial algorithms:
    - usually hopelessly exponential
    - quite limited use in practice

- The P class categorizes an algorithm based on its worst-case performance:
  - the largest running time over all possible input of lengths  $n$
- Some suggest an average-case approach would be better
  - it is not clear which distribution on the inputs should be adopted

# The reachability problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of the reachability problem: given a directed graph,  $G \subseteq V \times V$ ,  $V = \{v_1, \dots, v_n\}$  is there a path from  $v_i$  to  $v_j$ ?
  - the corresponding language:  $R = \{\rho(G)\rho(v_i)\rho(v_j) : \text{there is a path from } v_i \text{ to } v_j \text{ in } G\}$ 
    - $\rho(i)$ : some encoding of node  $v_i$  as strings
    - $\rho(G)$ : some encoding of graph  $G$  as strings
      - e.g.: adjacency matrix

# The reachability problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- A problem in general is:
  - a set of inputs, typically infinite
  - together with a yes-or-no question asked of each input
    - a property an input may or may not have

# The reachability problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- The corresponding language is the set of the encodings of those inputs which have the desired property
  - if a TM M decides a string that belongs to the language then we know the input have the desired property
- Problems and associated languages are interchangeable

- Theorem:  $R \in P$
- Proof:
  - compute the reflexive, transitive closure of  $G$ 
    - use the algorithm learned before
      - implement the algorithm with a random access TM
      - the algorithm needs  $O(n^3)$  time
      - random access TM can be simulated by ordinary TM in polynomial time
    - inspect if there is an edge between  $v_i$  and  $v_j$ 
      - if there is an edge then there is a path too

- Definition of the Euler cycle problem: given graph  $G$ , is there a closed path in  $G$  that uses each edge exactly once?
- A graph has Euler cycle if
  - for any pair of nodes  $(u, v)$ , neither of which is isolated,  $\exists$  path from  $u$  to  $v$
  - all nodes have equal number of incoming and outgoing arcs
- The Euler cycle problem is in  $P$

# Hamilton cycle problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of the Hamilton cycle problem: given graph  $G$ , is there a closed path in  $G$  that passes through each node exactly once?
- The Hamilton cycle problem is in NP

# The traveling salesman problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Optimization problems do not require a single yes or no answer so they can not be converted into language readily
- Modify the TSP: given a distance matrix and  $B > 0$  (budget), is there a tour  $\pi$  such that  $c(\pi) < B$  ?
- The TSP is in NP

- Definition of the partition problem: given a set of nonnegative integer  $a_1, \dots, a_n$  represented in binary form, is there a set  $P \subseteq \{1, \dots, n\}$  such that  $\sum_{i \in P} a_i = \sum_{i \notin P} a_i$ ?

- Definition of the unary partition problem: similar to the partition problem but the input is encoded in unary form
- The partition problem can be solved
  - calculate  $B(i) = \{\sum_{j \in K} a_j : K \subseteq \{1, \dots, i\}\}$  for  $\forall i$ 
    - sums of all possible partition of the first  $i$  number
  - $H = \sum_{i \in \{1, \dots, n\}} a_i / 2$
  - if  $H \in B(n) \rightarrow$  the answer is yes

```
B(0) = { 0 }
for i = 1, . . . , n
    B(i) = B(i-1)
    for ∀ j in B(i-1)
        add j+ai to B(i)
```

- The number of steps required by the algorithm:  $n2^n$ 
  - the new B set is usually twice as big as the previous one

- The binary (or decimal or any other radix system) encoding of the numbers is powerful
  - to encode the number  $n$  we need  $\log_2(n)$  digits
    - $\log$  is the inverse of  $\exp$

- There are  $n$  numbers and each has value about  $2^n$ 
  - number of steps:  $n2^n$
  - input size:
    - binary encoding:  $n^*n \rightarrow$  complexity: exponential
      - the partition problem is in NP
    - unary encoding:  $n2^n \rightarrow$  complexity: polynomial
      - the unary partition problem is in P
- The details about encoding objects does not effect the membership of the corresponding language in P
  - exception: when binary encoding is used instead of unary

- Definition of Boolean variable: a statement which can be true or false
  - the negate of a Boolean variable  $x$  is denoted with  $x^C$

- Definition of the clause of conjunctive normal form: an expression from Boolean variables and 'or' operators
- Definition of conjunctive normal form: an expression from clauses and 'and' operators
  - e.g.:  $(x_1 \vee x_2 \vee x_4^C) \wedge (x_2 \vee x_3^C \vee x_5^C)$

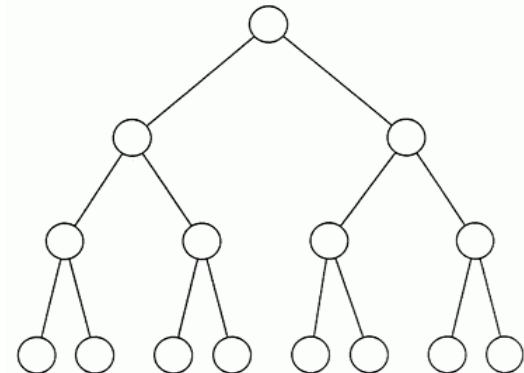
- Definition of truth assignment,  $T$ : a function which assign true or false value for each variable in formula  $F$ 
  - $T$  satisfies  $F$  if the value of  $F$  is true
  - $F$  is satisfiable if  $\exists$  a  $T$  which satisfies it
- Definition of the satisfiability problem: given a Boolean formula in conjunctive normal form, is it satisfiable ?
- The problem is in NP

- Definition of the 2 satisfiability problem: such a satisfiability problem whose formula contains only such clauses which have at most two variables
- The problem is in P
  - the solution method is described in the book

- Definition of polynomially bounded nondeterministic TM M:  $\exists$  a polynomial  $p(|x|)$  such that there is no configuration C for any input x for which  $(s, >[x]) \vdash_M^{p(|x|)+1} C$ 
  - M always halts at most  $p(|x|)$  steps

- Definition of NP: the class of languages decidable by a polynomially bounded nondeterministic TM
- Nondeterministic behavior is very strong quality
  - nondeterministic acceptance means for  $\forall w \in L \exists$  at least one computation which accepts  $w$  (all other configuration may reject the  $w$ )

- The operation of a nondeterministic TM can be represented by a tree
  - each node is a configuration
    - the root is the initial configuration
  - each node with 2 or more children represent a nondeterministic choice
  - with each new level
    - the length of the computation increases with just 1
    - the number of nodes may increase exponentially



- Theorem: the satisfiability problem is in NP
- Proof:
  - the nondeterministic TM which accepts the corresponding language works as follows
    - check if the input is really the encoding of a conjunctive normal form
    - count the variables in the formula
      - write 1 symbols on a second tape

- guess the values of the variables nondeterministically
- check the value of each clause
  - if at least one variable is true → the clause is true
- check the value of the formula
  - if all the clauses are true → the formula is true
- accept or reject based on the results
- all subtasks can be performed in polynomial time

- Similar proofs can be used to establish that other problems are in NP
  - use a nondeterministic TM to guess a string which helps to decide the input
    - the string is called certificate or witness
  - check deterministically if the certificate helps to accept or reject the input

- Example: TSP
  - guess the bijection  $\pi$
  - calculate the cost based on  $\pi$  and the distance matrix
  - check if the cost is smaller than the budget

- Definition of polynomial balanced language  $L'$ : if  $x; y \in L', L' \subseteq \Sigma^*; \Sigma^* \rightarrow \exists$  polynomial  $p(n)$  such that  $|y| < p(|x|)$ 
  - ;  $\notin \Sigma$
  - $x$  will be the input and  $y$  certificate

- Theorem: if  $L \in NP \leftrightarrow L' \in P$  and  $L'$  is polynomial balanced (the certificate is succinct)
  - $L' = \{x; y : y \text{ is a certificate for input } x\}$
  - the length of the certificates are polynomially bounded
  - the acceptance of the input can be determined in polynomial time from the certificate
- For the prime number problem the certificate is a pair of integer whose product is the input

- Definition of exponentially bounded deterministic TM  
 $M: \exists$  a polynomial  $p(|x|)$  such that there is no configuration  $C$  for any input  $x$  for which  
 $(s, >_x) \vdash_M^{\exp(p(|x|)+1)} C$ 
  - $M$  always halts at most  $\exp(p(|x|))$  steps

- Definition of EXP: the class of languages decidable by a exponentially bounded deterministic TM
- Theorem: if  $L \in NP \rightarrow L \in EXP$
- Proof: if  $L$  is decided by a nondeterministic TM  $M \rightarrow$  the deterministic version of  $M$  decides  $L$  in exponential time
  - see the equivalence of the nondeterministic and deterministic TM

- Theorem:  $P \subseteq NP \subseteq EXP$
- Proof:
  - see the previous proof for  $NP \subseteq EXP$
  - a deterministic TM can be seen as a special nondeterministic TM, so the definitions of P and NP imply  $P \subseteq NP$

- We know that  $E \notin P$ ,  $E \in EXP$ , so at least one of the previous inclusions are proper
  - $E = \{\rho(M)\rho(w) : M \text{ accepts input } w \text{ after at most } 2^{|w|} \text{ steps}\}$
- It is strongly supposed that both inclusions are proper but there is no proof for it
  - it is an unresolved question of the theory

- NP-complete problems have the following properties:
  - all problems in NP can be reduced to them via polynomial-time reductions
  - if there was a quick algorithm for any of these problems then all NP problems could be solved efficiently
- Unfortunately theory fails to prove that  $P \neq NP$ , but we strongly assume that
  - Clay Mathematics Institute has offered a \$1 million prize for the first correct proof

- Definition of polynomial-time computable function  $f: \exists a$  polynomially bounded TM that computes it
  - $f: \Sigma^* \rightarrow \Sigma^*$

- Definition of polynomial reduction from  $L_1$  to  $L_2$ ,  $\tau: \tau$  is a polynomial-time computable function,  $x \in L_1 \leftrightarrow \tau(x) \in L_2$ 
  - $L_1, L_2 \subseteq \Sigma^*$
  - $\tau: \Sigma^* \rightarrow \Sigma^*$
  - notation:  $L_1 \circ L_2$ 
    - $L_2$  is at least as hard as  $L_1$
  - it is analogue of the reductions we used to establish undecidability of problems with the help of  $H$

- Theorem: if  $L_2 \in P$ ,  $L_1 \circ L_2 \rightarrow L_1 \in P$
- Proof:
  - $L_2$  can be decided by  $M_2$  in polynomial time
  - $\tau$  can be computed in polynomial time
  - $L_1$  can also be decided by  $TM_2$  in polynomial time
- Theorem: if  $L_1 \in EXP$ ,  $L_1 \circ L_2 \rightarrow L_2 \in EXP$



# Introduction to the Theory of Computation

Lesson 11

## 7.1. Polynomial-time reductions

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**BEFEKTETÉS A JÖVŐBE**

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Theorem:  $\exists$  polynomial reduction from Hamilton cycle problem to Satisfiability

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Construction:
  - a G graph must be converted by  $\tau$  to a Boolean formula in conjunctive normal such a way that  $\tau(G)$  is satisfiable  $\leftrightarrow G$  has a Hamilton cycle
    - $G \subseteq V \times V, V = \{1, 2, \dots, n\}$
  - Hamilton cycle: a closed path visiting each node of G exactly once
  - formula  $\tau(G)$  will involve  $n^2$  Boolean variables:  $x_{i,j}$ ,  $1 \leq i, j \leq n$ 
    - $x_{i,j} = 1$  if node  $i$  of  $G$  is at the  $j$ -th position in the Hamilton cycle

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- the clauses of  $\tau(G)$  are the requirements for the existence of a Hamilton cycle
- requirement 1: exactly one node should appear  $j$ -th in the Hamilton cycle
  - at least one node should appear  $j$ -th in the Hamilton cycle
    - $(x_{1,j} \vee x_{2,j} \vee \dots \vee x_{n,j})$  for  $\forall j = 1, \dots, n$
  - two distinct nodes  $i$  and  $k$  should not appear  $j$ -th in the Hamilton cycle
    - $(x_{i,j}^C \vee x_{k,j}^C)$ , for  $i, j, k = 1, \dots, n$ , and  $j \neq k$

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- requirement 2: each node should appear in the Hamilton cycle
  - the j-th should appear at least once in the Hamilton cycle
    - $(x_{j,1} \vee x_{j,2} \vee \dots \vee x_{j,n})$  for  $\forall j = 1, \dots, n$
  - the j-th node should not appear in two distinct positions k and i in the Hamilton cycle
    - $(x_{i,j}^C \vee x_{i,k}^C)$ , for  $i, j, k = 1, \dots, n$ , and  $j \neq k$
- requirement 1 and 2 ensure that x represents a bijection between nodes and positions

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- requirement 3: if there is no  $(i, k)$  edge then  $i$  and  $k$  can not appear consecutively
  - $(x_{i,j}^C \vee x_{k,j+1}^C)$  for  $j = 1, \dots, n$ , for  $\forall (i, k)$  such that  $(i, k)$  is not an edge of  $G$ 
    - if  $j = n \rightarrow j+1 = 1$
- $\tau(G)$  has  $O(n^3)$  clauses, these clauses are simple
  - it is simple to construct a polynomial-time TM that computes  $\tau$

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Proof:  $G$  has a Hamilton cycle  $\leftrightarrow \tau(G)$  is satisfiable
  - suppose that  $G$  has a Hamilton cycle  $(\pi(1), \pi(2), \dots, \pi(n))$ 
    - let truth assignment  $T$  be the following:  $T(x_{i,j}) = \text{true} \leftrightarrow j = \pi(i)$
    - $T$  satisfies all clauses of  $\tau(G)$  because of the construction

# Polynomial reduction from Hamilton cycle problem to Satisfiability

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- suppose  $T$  is a truth assignment that satisfies  $\tau(G)$ 
  - requirement 1 and 2 ensure that  $T$  assigns one position for  $\forall$  node
  - denotes with  $\pi(i)$  the unique  $j$  for which  $T(x_{i,\pi(i)})$  is true
  - requirement 3 ensures that if  $i = \pi(j)$  and  $k = \pi(j+1)$  then  $(j, k)$  must be an edge of  $G$
- There is also a polynomial reduction from Satisfiability to Hamilton cycle problem

- Definition of the Two-machine scheduling problem:
  - $n$  tasks must be scheduled on two machines
  - both machines have the same speed
  - each task can be executed on either machine
  - there are no restrictions on the order of the execution

- given:
  - the execution times  $a_1, \dots, a_n$  of the tasks
  - deadline D (all in binary form)
- question: can you complete all these tasks on the two machines within the deadline?
  - $\exists$  partition of the numbers into two sets so that the numbers in each set add up to D or less?

- Definition of Knapsack problem:
  - given:
    - a set  $S = \{a_1, \dots, a_n\}$  of nonnegative integers
    - an integer  $K$  (all in binary form)
  - question:  $\exists$  a subset  $P \subseteq S$  such that  $\sum_{a \in P} a = K$  ?
- Between two-machine scheduling problem, Knapsack problem and partition problem there are six polynomial reductions reducing any of these problems to any other
- Example:  
reduction from Knapsack to Partition

# Polynomial reduction from Knapsack problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem:  $\exists$  polynomial reduction from Knapsack problem to Partition problem
- Construction:
  - given an instance of Knapsack problem, with integers  $a_1, \dots, a_n$ , and  $K$
  - if  $K = H = \frac{1}{2}\sum_{i=1}^n a_i \rightarrow$  the reduction is to erase  $K$  from the input
  - if  $K \neq H \rightarrow$  erase  $K$  from the input and add two new integers  $a_{n+1} = 2H + 2K$ ,  $a_{n+2} = 4H$  to the set of  $a_i$ 's

# Polynomial reduction from Knapsack problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Proof:
  - in partition the new integers must be on opposite sides

$$\sum_{i=1 \dots n} a_i < a_{n+1} + a_{n+2} = 2H + 2K + 4H$$

- write the partition of the numbers with set P

$$4H + \sum_{a \in P} a = 2H + 2K + \sum_{a \in S-P} a$$

- $P \subseteq \{1, \dots, n\}$

- add  $\sum_{a \in P} a$  to both sides

$$4H + 2 \sum_{a \in P} a = 4H + 2K$$

$$\sum_{a \in P} a = K$$

# Polynomial reduction from Knapsack problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- the resulting instance of Partition problem has a solution  $\leftrightarrow$  the original instance of Knapsack problem has one
  - the Knapsack corresponds to P
- the reduction can be carried out in polynomial time

# Polynomial reduction from Knapsack problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Theorem:  $\exists$  polynomial reduction from Partition problem to Knapsack problem
  - Partition problem is a special case of Knapsack problem
  - the reduction calculates bound  $K = \frac{1}{2} \sum_{i=1}^n a_i$ 
    - if  $K$  is not an integer  $\rightarrow$  the given instance of Knapsack and of Partition have no solution

# Reduction from Partition problem to Two-machine scheduling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Theorem:  $\exists$  polynomial reduction from Partition problem to Two-machine scheduling problem
- Construction:
  - given an instance of Partition problem with integers  $a_1, \dots, a_n$
  - the reduction keeps  $a_1, \dots, a_n$  as execution times and calculates deadline  $D = \frac{1}{2} \sum_{i=1}^n a_i$ 
    - if  $D$  is not an integer  $\rightarrow$  the given instance of Two-machine scheduling and Partition problem have no solution

# Reduction from Partition problem to Two-machine scheduling problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Proof:
  - if the sums of the two sets are exactly D (Partition problem) → at most D is also true (Two-machine scheduling problem)
  - if the sums of the two sets are at most D (Two-machine scheduling problem) and D is the half-sum → the sums are exactly D (Partition problem)

# Reduction from Two-machine scheduling problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Theorem:  $\exists$  polynomial reduction from Two-machine scheduling problem to Partition problem

# Reduction from Two-machine scheduling problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Construction:
  - given an instance of two-machine scheduling with task lengths  $a_1, \dots, a_n$  and deadline D
  - idle time:  $I = 2D - \sum_{i=1}^n a_i$ 
    - the sum of times while the machines do not work
  - add some new numbers to the set of lengths of tasks
    - the sum of these new numbers is I
    - any number between 0 and I can be created by adding a subset of these new numbers

# Reduction from Two-machine scheduling problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

– determine the new numbers

- 1<sup>st</sup> solution: we could add  $I$  copies of 1, but this is not a polynomial-time reduction
  - it would represent a kind of unary encoding of  $I$  while all other numbers are given binary
- 2<sup>nd</sup> solution:
  - powers of 2 if they are smaller than  $I/2$
  - another integer to bring the sum to  $I$
  - e.g.: if  $I = 56 \rightarrow$  new integers: 1, 2, 4, 8, 16 and 25

# Reduction from Two-machine scheduling problem to Partition problem

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- Proof:
  - if the sums of the two sets are at most D (Two-machine scheduling problem) → adding the new integers to the problem helps to keep busy the machines in the idle time
    - if there is no idle time → we have a Partition problem
  - if the sums of the two sets are exactly D (Partition problem) → at most D is also true (Two-machine scheduling problem)

- Theorem: if  $\tau_1$  is a polynomial reduction from  $L_1$  to  $L_2$  and  $\tau_2$  is a polynomial reduction from  $L_2$  to  $L_3 \rightarrow \tau_1 \circ \tau_2$  is a polynomial reduction from  $L_1$  to  $L_3$

- Proof:
  - $\tau_1$  is computed by TM  $M_1$  in time  $p_1$  (polynomial)
  - $\tau_2$  is computed by TM  $M_2$  in time  $p_2$  (polynomial)
  - $\tau_1 \circ \tau_2$  can be computed by  $M_1 M_2$  in time  $p_2(p_1(|x|))$ 
    - the input of  $M_2$  is longer than the input of  $M_1$
    - $|\tau_1(x)| \leq p_1(|x|)$  in each step  $M_1$  can write at most 1 symbol to the tape
  - $x \in L_1 \leftrightarrow \tau_2(x) \in L_2 \leftrightarrow \tau_1 \circ \tau_2(x) \in L_3$

- Definition: A language  $L \subseteq \Sigma^*$  is called  $\text{NP}$ -complete if
  - $L \in \text{NP}$
  - and for every language  $L' \in \text{NP}$ ,  $\exists$  polynomial reduction from  $L'$  to  $L$

- Theorem: let  $L$  be an  $N\Pi$ -complete language,  
 $\Pi = N\Pi \leftrightarrow L \in \Pi$
- Proof: only if →
  - $L \in N\Pi$ -complete  $\rightarrow L \in \Pi$
  - $L \in \Pi$ ,  $N\Pi = \Pi$ ,  $\rightarrow L \in \Pi$
- Proof: if ←
  - $L \in \Pi$ -complete is decided by a deterministic TM  $M_1$  in time  $p_1(n)$ , a polynomial

- let  $L' \in \text{N}\Pi$
- $\exists$  a polynomial reduction  $\tau$  from  $L'$  to  $L$  according to the definition of  $\text{N}\Pi$ -complete
  - $\tau$  is computed by TM  $M_2$  in time  $p_2(n)$ , a polynomial
- $M_2 M_1$  decides  $L'$  in time  $p_1(p_2(|x|))$ , a polynomial
  - the input of  $M_1$  is longer than the input of  $M_2$
  - $|\tau(x)| \leq p_2(|x|)$ , in each step  $M_2$  can write at most 1 symbol to the tape
- if  $\forall L' \in \text{N}\Pi$  can be decided in polynomial time  $\rightarrow \Pi = \text{N}\Pi$

- NP-hard, co-NP, co-NP-complete with examples



# Introduction to the Theory of Computation

Lesson 12

7.2. Cook's Theorem

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**BEFEKTETÉS A JÖVŐBE**

- Theorem: if  $L_1 \in \text{NP-complete}$ ,  $\exists$  a polynomial reduction from  $L_1$  to  $L_2 \rightarrow L_2 \in \text{NP-complete}$
- Proof:
  - $\exists$  polynomial reduction from any  $L' \in \text{NP}$  to  $L_1$  because  $L_1 \in \text{NP-complete}$
  - $\exists$  polynomial reduction from  $L_1$  to  $L_2$
  - $\exists$  polynomial reduction from any  $L' \in \text{NP}$  to  $L_2$  because of the transitivity of the polynomial reductions

- Satisfiability is the first problem shown to be  $\text{NP}$ -complete (Stephen A. Cook, 1971)
- We prove that the Bounded tiling problem is  $\text{NP}$ -complete using the definition of  $\text{NP}$ -complete

- Definition of Bounded tiling problem:
  - given:
    - a tiling system ( $D, H, V$ )
      - $D$  set of tiles
      - $H$  relation of horizontally joinable tiles
      - $V$  relation of vertically joinable tiles
    - the length of the square to be tiled:  $s \in \mathbb{Z}^+$
    - the first row of tiles as a function  $f_0 : \{0, \dots, s-1\} \rightarrow D$ 
      - in tiling problem only the  $(0, 0)$  tile is given
    - $\exists$  tiling of an  $s \times s$  square? (and not of a whole quadrant)

– tiling is a function  $f : \{0, 1, \dots, s-1\} \times \{0, 1, \dots, s-1\} \rightarrow D$ , such that

- $f(m, 0) = f_0(m), \forall m < s$
- $(f(m, n), f(m+1, n)) \in H, \forall m < s - 1, n < s$
- $(f(m, n), f(m, n+1)) \in V, \forall m < s, n < s - 1$

- Theorem: Bounded tiling problem is  $\text{NP}$ -complete
- Proof:
  - certificate of the problem: a complete listing of the  $s^2$  values of the tiling function  $f$

- the certificate is succinct:
  - the certificate contains  $s^2$  encoding of tiles
  - the input contains  $s$  encoding of tiles in  $f_0$ 
    - if only the corner was given → the certificate is not succinct
- we can check in polynomial time if  $f$  is really a tiling function
- the problem has a succinct certificate, checkable in polynomial time so it is in  $\text{N}\Pi$

- we must show  $\forall L \in N\Pi$  can be reduced via a polynomial reduction to the Bounded tiling problem
- consider any  $L \in N\Pi$ 
  - we only know that  $\exists$  a nondeterministic TM  $M$  deciding  $L$  in polynomial time,  $p(|x|)$
- $\tau$ , a polynomial-time computable function
  - calculates  $s = p(|x|) + 2$
  - defines  $(D, H, V)$  based on  $M$ 
    - tiles are defined in the same way they were defined at the Tiling problem (not detailed now)
      - » based on  $K, \Sigma, \delta$  of  $M$

- the edges of the tiles are signed with the elements of the previous sets
- tiles with the same signs can be joined
- calculates  $f_0$  the 0<sup>th</sup> row of the  $s \times s$  square
  - $f_0(0)$  is a tile with upper edge marking >
  - $f_0(1)$  is a tile with upper edge marking  $(s, \int)$
  - $f_0(i+1)$  for  $i = 1, \dots, |x|$  is a tile with upper edge marking  $x_i$
  - $f_0(i)$  for all  $i > |x| + 1$ , is a tile with upper edge marking  $\int$

- the markings of the horizontal edges between rows  $t$  and  $t+1$  represent the tape content of  $M$  after the  $t^{\text{th}}$  step (starting with input  $x$ )
  - can be proved with induction
  - the horizontal edge markings between the 0<sup>th</sup> and the 1<sup>st</sup> rows will spell the initial configuration of  $M(s, > \sqsupset x)$  on input  $x$

- there is a tile not used before with  $(y, a)$  lower and upper markings
  - $y$  is the accepting state,  $a \in \Sigma$
  - it makes possible to repeat the last row if the input is accepted

- $x \in L \leftrightarrow M \text{ accepts } x \leftrightarrow T(x) \text{ has a tiling}$ 
  - suppose  $\exists T(x)$  has a tiling
    - the last row contains tiles too
  - the computation of  $M$  on input  $x$  can not continue for more than  $p(|x|) = s-2$  steps
    - the upper markings of the  $s-2^{\text{nd}}$  row must contain one of the symbols  $y$  and  $n$
  - since there are tiles in the  $s-1^{\text{st}}$  (last) row too, and there is no tile with lower marking  $n$ , symbol  $y$  must appear in the upper markings of the  $s-2^{\text{nd}}$  row
    - the computation is accepting

- conversely if  $M$  accepts  $x$  then it can be simulated by a tiling
  - we know the horizontal edge markings correspond to tape contents
  - the last row containing an accepting state may be repeated several times to complete the tiling

- Theorem: (Cook's theorem) Satisfiability is NΠ-complete
  - we have already proved that Satisfiability is in NΠ
  - reducing the Bounded tiling problem to Satisfiability implies that the latter problem is also NΠ-complete
  - given a tiling system  $(D, H, V)$ , side  $s$ , bottom row  $f_0$ , where  $D = \{d_1, \dots, d_k\}$
  - construct a Boolean formula  $\tau(D, H, V, s, f_0)$ , such that  $\exists$  an  $s \times s$  tiling  $f \leftrightarrow \tau(D, H, V, s, f_0)$  is satisfiable

- Construction:
  - introduce  $x_{m,n,d}$  for  $\forall 0 \leq m, n < s, d \in D$ 
    - $x_{m,n,d} = \text{true} \leftrightarrow f(m, n) = d$
  - the clauses of  $\tau$  are the requirements for the existence of a tiling
  - requirement 1: exactly one tile should appear in each position
    - each position has at least one tile
      - $(x_{m,n,d_1} \vee x_{m,n,d_2} \vee \dots \vee x_{m,n,d_k}), \forall m, n < s$
    - a position can not have more than one tile
      - $(x_{m,n,d}^C \vee x_{m,n,d'}^C), \forall m, n < s, d \neq d'$

- requirement 2: the first row is given by  $f_0$ 
  - $(x_{i,0,f_0(i)}), i = 0, \dots, s-1$
- requirement 3: two tiles can not be next to each other if relation  $H$  does not allow it
  - $(x_{m,n,d}^C \vee x_{m+1, n, d'}^C), \forall n < s, m < s-1, (d, d') \notin H$ 
    - either  $d$  can not appear in  $(m, n)$  or  $d'$  can not appear in  $(m+1, n)$  if  $(d, d') \notin H$

- requirement 4: two tiles can not be one on another if relation  $V$  does not allow it
  - $(x_{m,n,d}^C \vee x_{m, n+1, d'}^C), \forall n < s-1, m < s, (d, d') \notin V$ 
    - either  $d$  can not appear in  $(m, n)$  or  $d'$  can not appear in  $(m, n+1)$  if  $(d, d') \notin V$

- Proof: only if →
  - suppose  $\tau(D, H, V, s, f_0)$  is satisfiable by the truth assignment  $T$
  - since the clauses of requirement 1 are true  $T$  represents a function
$$f : \{0, \dots, s - 1\} \times \{0, \dots, s - 1\} \rightarrow D$$
  - since the clauses of requirement 2 are true  $f$  extends  $f_0$
  - since the clauses of requirement 3 and 4 are true  $f$  represents a legal tiling

- Proof: if  $\leftarrow$ 
  - suppose  $f$  is a legal tiling
  - it is easy to see that all clauses of  $\tau$  are true

- Theorem: 3 - Satisfiability is  $\text{NP}$ -complete
  - special case of Satisfiability in which all clauses involve 3 or fewer literals
  - it is in  $\text{NP}$ , as it is a special case of Satisfiability known to be in  $\text{NP}$
  - we reduce Satisfiability to 3 - Satisfiability

- Construction:
  - leave the short clauses unchanged
  - for  $\forall$  long clause  $(\lambda_1 \vee \lambda_2 \vee \dots \vee \lambda_k)$  in  $F$ ,  $k > 3$ 
    - introduce new Boolean variables:  $y_1, \dots, y_{k-3}$ 
      - $y_i = \lambda_{i+2} \vee \dots \vee \lambda_k$
    - replace the long clause with
$$(\lambda_1 \vee \lambda_2 \vee y_1), (y_1^C \vee \lambda_3 \vee y_2), (y_2^C \vee \lambda_4 \vee y_3), \dots,$$
$$(y_{k-4}^C \vee \lambda_{k-2} \vee y_{k-3}), (y_{k-3}^C \vee \lambda_{k-1} \vee \lambda_k)$$
      - $(y_i^C \vee \lambda_{i+2} \vee y_{i+1})$  states if  $y_i = \text{true} \rightarrow \lambda_{i+2} = \text{true}$  or  $y_{i+1} = \text{true}$
  - $T$  can be carried out in polynomial time

- Proof: only if →
  - suppose that truth assignment  $T$  satisfies  $\tau(F)$
  - $T$  satisfies the short clauses of  $F$  as they are the same as in  $\tau(F)$
  - regard any long clause  $C$  of  $F$ , at least one of  $\lambda_i$  of  $C$  is true
    - suppose  $\forall \lambda_i$  of  $C$  is false
    - the clauses corresponding to  $C$  in  $\tau(F)$  are true because of the  $y_i$  variables
      - $y_1 = \text{true}$ , because of  $(\lambda_1 \vee \lambda_2 \vee y_1)$
      - $y_2 = \text{true}$ , because of  $(y_1^C \vee \lambda_3 \vee y_2)$
      - ...

- $y_{k-3} = \text{true}$ , because of  $(y_{k-4}^C \vee \lambda_{k-2} \vee y_{k-3})$
- $(y_{k-3}^C \vee \lambda_{k-1} \vee \lambda_k) = \text{false}$ , but we know  $T$  satisfies  $T(F)$  so the assumption " $\forall \lambda_i$  of  $C$  is false" is not valid
- $T$  satisfies both the short and long clauses of  $F$
- Proof: if  $\leftarrow$ 
  - suppose that truth assignment  $T$  satisfies  $F$
  - let  $T'$  the extension of  $T$  for the  $y_i$  values
  - for  $\forall$  long clause  $C = (\lambda_1 \vee \lambda_2 \vee \dots \vee \lambda_k)$  of  $F$ 
    - $j = \operatorname{argmin}\{T(\lambda_j) = \text{true}\}$ 
      - the smallest index for which  $T(\lambda_j) = \text{true}$
      - $j$  exists because  $T$  satisfies  $F$

- $T'(y_i) = \text{true}$ ,  $i \leq j - 2$
- $T'(y_i) = \text{false}$ , otherwise
- it is easy to see that  $T'$  satisfies the clauses corresponding to  $C$  (if  $C$  is satisfiable)
  - $(\lambda_1 \vee \lambda_2 \vee y_1)$  states  $\lambda_1 = \text{true}$  or  $\lambda_2 = \text{true}$  or the rest of  $\lambda$  is true
  - $(y_{k-3}^C \vee \lambda_{k-1} \vee \lambda_k)$  states if  $y_{k-3} = \text{true} \rightarrow \lambda_{k-1} = \text{true}$  or  $\lambda_k = \text{true}$

- Definition of Max set problem: given a set  $F$  of clauses, and an integer  $K$ ,  $\exists$  a truth assignment that satisfies at least  $K$  of the clauses?
- Theorem: Max set is  $N\Pi$ -complete
- Proof:
  - Max set problem is a generalization of Satisfiability
    - an instance of Satisfiability can be thought of as an instance of Max set problem with special  $K$ 
      - $K =$  the number of clauses in  $F$

- Construction:
  - given an instance  $F$  of Satisfiability with  $m$  clauses
  - create  $(F, m)$ , the input of the Max set problem, by appending to  $F$  the parameter  $K = m$
- Proof:
  - suppose  $\exists$  truth assignment  $T$  satisfying  $F$
  - $T$  also satisfies  $(F, m)$
  - suppose  $\exists$  truth assignment  $T'$  satisfying  $(F, m)$
  - $T'$  also satisfies  $F$



# Introduction to the Theory of Computation

Lesson 13

7.3. More NP-complete problems

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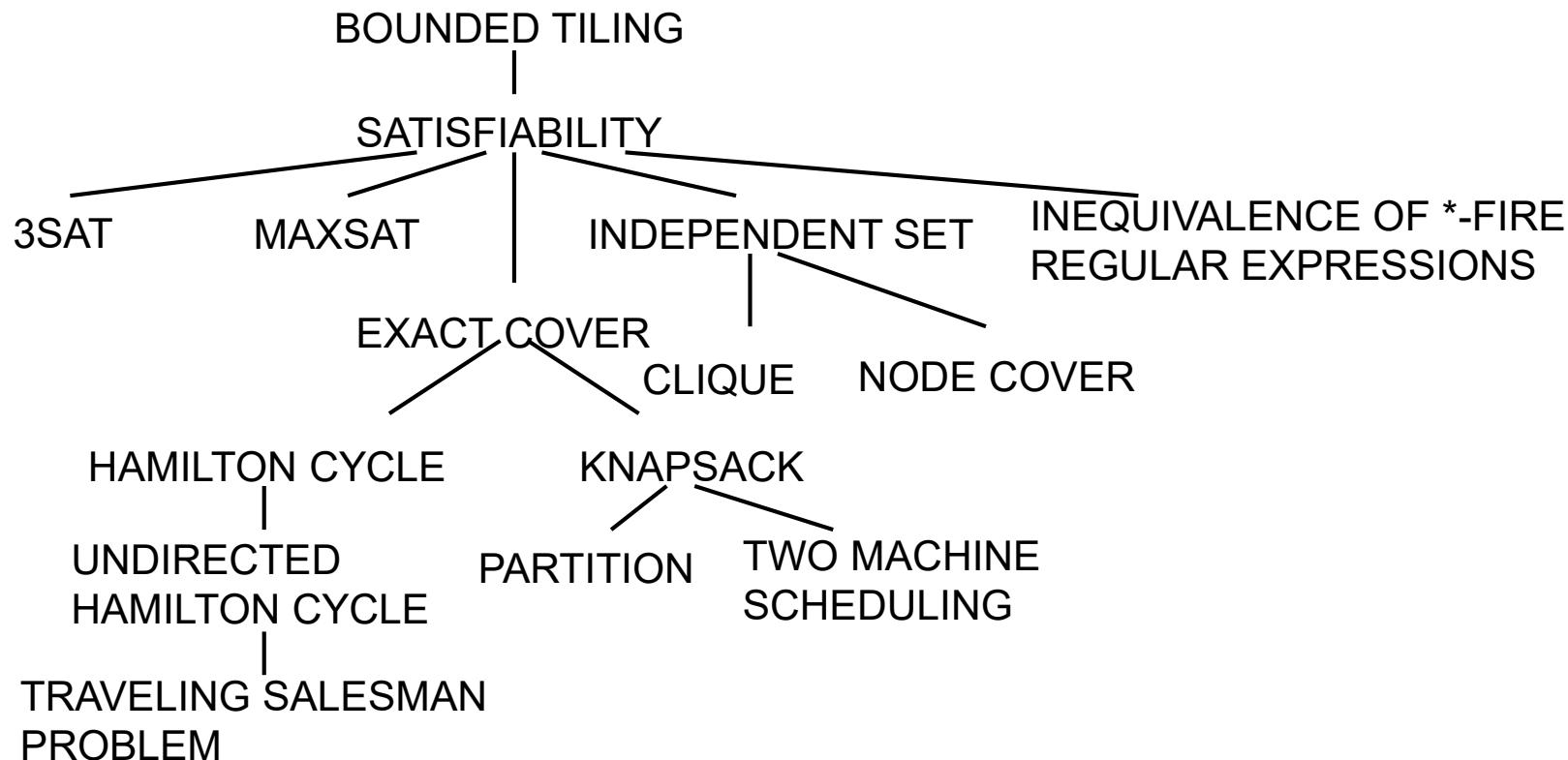
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- NΠ-complete problems are important because they
  - come up frequently under various guises in applications
    - do not try to devise an effective algorithm for an NΠ-complete problem
  - have been studied so much
  - are often handy starting points for NΠ-complete reductions
- Satisfiability is important for all three reasons

# More NΠ-complete problems

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
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- Definition of the Exact cover problem:
  - given:
    - a finite set  $U = \{u_1, \dots, u_n\}$ , called universe
    - $F = \{S_1, \dots, S_m\}$ , called family,  $S_i \subseteq U, i = 1, 2 \dots m$
  - is there an exact cover  $C$  called subfamily, such that
    - $C \subseteq F$
    - the sets in  $C$  are disjoint
    - $\cup C = U$

- Example:
  - universe  $U = \{u_1, u_2, u_3, u_4, u_5, u_6\}$
  - family  $F = \{\{u_1, u_3\}, \{u_2, u_3, u_6\}, \{u_1, u_5\}, \{u_2, u_3, u_4\}, \{u_5, u_6\}, \{u_2, u_4\}\}$
  - exact cover  $C = \{\{u_1, u_3\}, \{u_5, u_6\}, \{u_2, u_4\}\}$

- Theorem: Exact cover is  $N\Pi$ -complete
- Proof:
  - Exact cover is in  $N\Pi$ 
    - $C$  is a valid certificate of  $(U, F)$
    - $C$  is polynomially concise
      - $C$  is part of  $(U, F)$
      - can be checked in polynomial time if  $C$  is really an exact cover

- Construction:
  - reduce Satisfiability to Exact cover
  - given a Boolean formula  $F$  with clauses  $(C_1, \dots, C_l)$  over the Boolean variables  $x_1, \dots, x_n$ 
    - beware  $F$  has two meaning now
  - we show how to construct in polynomial time an equivalent instance  $\tau(F)$  of the Exact cover problem
  - let  $\lambda_{j,k}$  the  $k^{\text{th}}$  literal of clause  $C_j$ 
    - $C_j$  has  $m_j$  literal

- the universe of  $\tau(F)$ :  $U = \{x_i : 1 \leq i \leq n\} \cup \{C_j : j = 1, \dots, l\} \cup \{p_{j,k} : 1 \leq j \leq l, k = 1, \dots, m_j\}$ 
  - one element for each Boolean variable
  - one for each clause
  - one for each position in each clause
- the family of  $\tau(F)$ :
  - $\{p_{j,k}\}$  for  $\forall p_{j,k}$ 
    - $p_{j,k}$ 's are easy to cover
  - $T_{i,\text{true}} = \{x_i\} \cup \{p_{j,k} : \lambda_{j,k} = x_i^C\}$ 
    - for  $\forall x_i$
    - it contains the position for all negative occurrences of  $x_i$

- $T_{i,\text{false}} = \{x_i\} \cup \{p_{j,k} : \lambda_{j,k} = x_i\}$ 
  - for  $\forall x_i$
  - it contains the position for all positive occurrences of  $x_i$
- $\{C_j, p_{j,k}\}$  for  $\forall$  clause  $C_i$ , for  $\forall$  literal in it

- Lemma:  $\tau(F)$  has an exact cover  $\leftrightarrow F$  is satisfiable
- Proof,  $\rightarrow$ :
  - suppose that an exact cover  $C$  exists
    - beware  $C$  has two meanings now
  - construct truth assignment  $T$  as:
    - if  $T_{i,true} \in C \rightarrow T(x_i) = \text{true}$
    - if  $T_{i,false} \in C \rightarrow T(x_i) = \text{false}$ 
      - beware  $T$  has two meaning now
    - since  $x_i$  must be covered exactly once, either  $T_{i,true} \in C$  or  $T_{i,false} \in C$

- $T$  satisfies  $F$ 
  - consider clause  $C_j$
  - $C_j$  must be covered by one such set  $\{C_j, p_{j,k}\}$
  - $p_{j,k}$  must appear in either  $T_{i,true}$  or  $T_{i,false}$  for some  $i$ 
    - see the definition of  $T_{i,true}$  and  $T_{i,false}$

- suppose that the  $k^{\text{th}}$  literal in  $C_j$  is  $x_i$ 
  - then  $T_{i,\text{false}}$  contains  $p_{j,k}$
  - exact cover  $C$  must contain  $T_{i,\text{true}}$ 
    - » which does not contain  $p_{j,k}$
  - according to the truth assignment  $T(x_i) = \text{true}$
  - $T$  makes true the  $k^{\text{th}}$  literal of  $C_j$
  - if you suppose that the  $k^{\text{th}}$  literal in  $C_j$  is  $x_i^C$   
then  $T(x_i) = \text{false}$ , so  $C_j$  is true again
- the previous is true for  $\forall C_j$ , hence  $F$  is satisfiable

- Proof,  $\leftarrow$ :
  - suppose that truth assignment  $T$  satisfies  $F$
  - construct an exact cover  $C$  by selecting sets:
    - $T_{i,\text{true}}$  if  $T(x_i) = \text{true}$
    - $T_{i,\text{false}}$  if  $T(x_i) = \text{false}$ 
      - for  $\forall x_i$
    - $\{C_j, p_{j,k}\}$ ,  $k$  should be such that the  $k^{\text{th}}$  literal of  $C_j$  is made true by  $T$  (there may be several such literal) and  $p_{j,k}$  have not occurred yet in  $C$ 
      - for  $\forall C_j$
    - $\{p_{j,k}\}$  for such  $p$ 's which has not occurred yet
  - $C$  covers  $\tau(F)$  as the  $x$ ,  $C$ , and  $p$  variables are covered

# Exact cover: Example

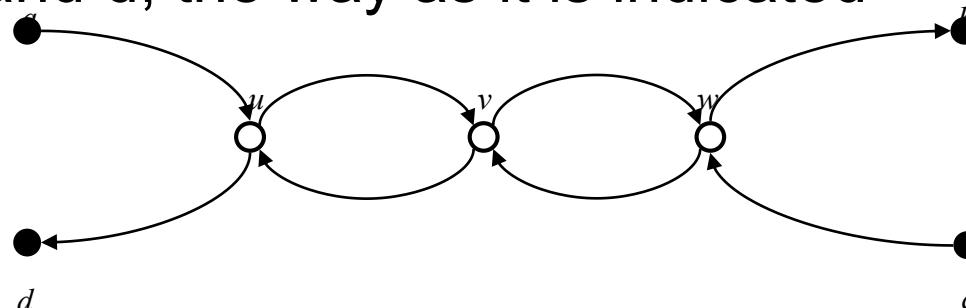
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- Formulae  $F = C_1 \wedge C_2 \wedge C_3 \wedge C_4$ , where  $C_1 = (x_1 \vee x_2^C)$ ,  $C_2 = (x_1^C \vee x_2 \vee x_3)$ ,  $C_3 = (x_2)$ ,  $C_4 = (x_2^C \vee x_3^C)$
- Universe  $U = \{x_1, x_2, x_3, x_4, C_1, C_2, C_3, C_4, p_{1,1}, p_{1,2}, p_{2,1}, p_{2,2}, p_{2,3}, p_{3,1}, p_{4,1}, p_{4,2}\}$
- Family  $F = \{$ 
  - $\{p_{1,1}\}, \{p_{1,2}\}, \{p_{2,1}\}, \{p_{2,2}\}, \{p_{2,3}\}, \{p_{3,1}\}, \{p_{4,1}\}, \{p_{4,2}\}$
  - $T_{1,true} = \{x_1, p_{2,1}\}, T_{1,false} = \{x_1, p_{1,1}\}$
  - $T_{2,true} = \{x_2, p_{1,2}, p_{4,1}\}, T_{2,false} = \{x_2, p_{2,2}, p_{3,1}\}$
  - $T_{3,true} = \{x_3, p_{4,2}\}, T_{3,false} = \{x_3, p_{2,3}\}$
  - $\{C_1, p_{1,1}\}, \{C_1, p_{1,2}\}, \{C_2, p_{2,1}\}, \{C_2, p_{2,2}\}, \{C_2, p_{2,3}\}, \{C_3, p_{3,1}\}, \{C_4, p_{4,1}\}, \{C_4, p_{4,2}\} \}$

- T truth assignment:  $T(x_1) = \text{true}$ ,  $T(x_2) = \text{true}$ ,  $T(x_3) = \text{false}$
- Exact cover  $C = \{T_{1,\text{true}}, T_{2,\text{true}}, T_{3,\text{false}}, \{C_1, p_{1,1}\}, \{C_2, p_{2,2}\}, \{C_3, p_{3,1}\}, \{C_4, p_{4,2}\}, \{p_{1,2}\}, \{p_{2,1}\}, \{p_{2,3}\}, \{p_{4,1}\}\}$

- Definition of gadget: a graph displayed with abbreviation
  - consider the next sub-graph of graph G
    - there is a series of node: u, v, w, ...
    - there are forward and backward arcs between (u, v), (v, w), ...
  - only the first and the last nodes in the series are connected to the rest of the graph through nodes a, b c, and d, the way as it is indicated



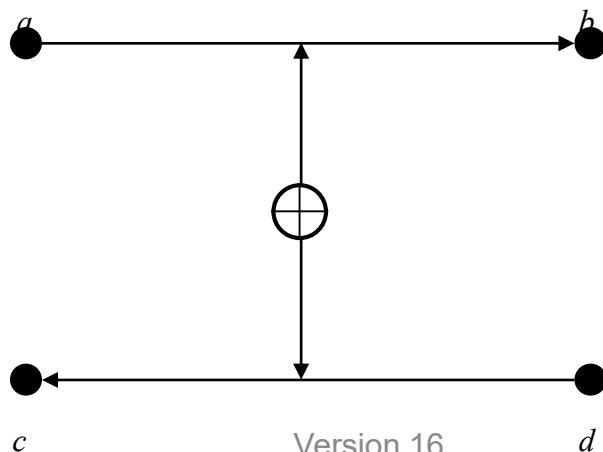
- search for the Hamilton cycle of G
  - a cycle which traverses each node exactly once
- via which edges is the Hamilton cycle going to traverse the middle nodes: u, v, w?
  - (a, u), (u, v), (v, w), (w, b) or
  - (c, w), (w, v), (v, u), (u, d)
  - all other possibilities omit some nodes or traverse through a node more than once

# Hamilton cycle

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
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- sequence  $(a, u), (u, v), (v, w), (w, b)$  and  $(c, w), (w, v), (v, u), (u, d)$  can be replaced by  $(a, b)$  and  $(c, d)$  respectively
- in any Hamilton cycle of  $G$  either  $(a, b)$  is traversed or  $(c, d)$  is, but not both (exclusive or)
- this situation is depicted in the figure where the two edges are connected by an exclusive or by a sign



- Theorem: Hamilton cycle is  $\text{NP}\Pi$ -complete
- Proof:
  - we already know that it is  $\text{NP}\Pi$
  - reduce Exact cover to Hamilton cycle with a polynomial algorithm

- Construction:
  - given an instance  $(U, F)$  of Exact cover
    - $U = \{u_1, \dots, u_n\}$  and  $F = \{S_1, \dots, S_m\}$
  - we shall describe a polynomial-time algorithm which produces a directed graph  $G = \tau(U, F)$
  - graph  $G$  has nodes:  $u_0, u_1, \dots, u_n$  and  $S_0, S_1, \dots, S_m$ 
    - one for each element of the universe, and for each set in the family, plus two more nodes

- graph G has two  $(S_{i-1}, S_i)$  arcs,  $i = 1, \dots, m$ 
  - two different edges connecting the same pair of nodes are needed because of the exclusive or
  - one of them is called the long edge, the other is the short edge
- graph G has  $k (u_{j-1}, u_j)$  arcs,  $j = 1, \dots, n$ , where  $k$  is the number of appearance of  $u_j$  in the  $S$  sets
- graph G has arcs  $(u_n, S_0)$  and  $(S_m, u_0)$

- join  $\forall$  copy of edge  $(u_{j-1}, u_j)$  with the long edge  $(S_{i-1}, S_i)$  by an exclusive or such that  $u_j \in S_i$ 
  - as each copy of edge  $(u_{j-1}, u_j)$  corresponds to an appearance of  $u_j$  in some  $S_i$

- Lemma:  $\tau(U, F)$  has a Hamilton cycle  $\leftrightarrow (U, F)$  has an exact cover
- Proof,  $\rightarrow$ :
  - suppose that  $\tau(U, F)$  has a Hamilton cycle
    - do not forget that the exclusive or constraints must hold
      - if we ignore them it is easy to give Hamilton cycle
    - the cycle contains the nodes in this order:  $S_0, S_1, \dots, S_m, U_0, U_1, \dots, U_n, S_0$ 
      - there are no backward arcs and the starting point is irrelevant

- let  $C$  be the set of all such  $S_i$  sets where the short edge  $(S_{i-1}, S_i)$  is traversed by the Hamilton cycle
- the Hamilton cycle traverses exactly one of the copies of  $(u_{j-1}, u_j)$  for  $\forall u_j \in U \rightarrow \forall u_j$  is contained in exactly one  $S_i$  set in  $C$ 
  - there may be several  $S_i$  sets which contains  $u_j$
  - there is one  $S_i$  such that  $u_j \in S_i \in C$ 
    - if the Hamilton cycle traverses a copy of  $(u_{j-1}, u_j) \leftrightarrow$  it does not traverse the long edge  $(S_{i-1}, S_i)$  connected to the previous edge through the exclusive or

- if the Hamilton cycle does not traverse the long edge  $(S_{i-1}, S_i) \leftrightarrow$  it do traverses the short edge  $(S_{i-1}, S_i)$  because there are just two arcs between  $S_{i-1}, S_i$
- if the Hamilton cycle traverses the short edge  $(S_{i-1}, S_i) \rightarrow S_i \in C$  by definition

- there is no more than one  $S_i$  such that  $u_j \in S_i \in C$ 
  - if the Hamilton cycle traverses a copy of  $(u_{j-1}, u_j)$   
 $\leftrightarrow$  it does not traverse the other copies
  - if the Hamilton cycle does not traverse the other copies of  $(u_{j-1}, u_j)$   $\leftrightarrow$  it traverses the long edges  $(S_{i-1}, S_i)$  connected to them through the exclusive or
  - if the Hamilton cycle traverses the long edges  $(S_{i-1}, S_i)$   $\leftrightarrow$  it does not traverse the short edges  $(S_{i-1}, S_i)$
  - if the Hamilton cycle does not traverse the short edges  $(S_{i-1}, S_i)$   $\leftrightarrow$  these  $S_i$  are not in  $C$  by definition

- if  $\forall u_j$  is contained in exactly one  $S$  set in  $C \rightarrow C$  is an exact cover by definition
  - the union of the  $S$  sets in  $C$  results in  $U$  without overlap

- conversely, suppose that an exact cover  $C$  exists then a Hamilton cycle in the graph  $\tau(U, F)$  can be constructed as follows:
  - traverse the short copies of all edges  $(S_{j-1}, S_j)$  where  $S_j \in C$ , and the long edges for all other sets
  - for each element  $u_i$  traverse the copy of the edge  $(u_{i-1}, u_i)$ , that corresponds to the unique set in  $C$  that contains  $u_i$
  - complete the Hamilton cycle by the edges  $(u_n, S_0)$  and  $(S_m, u_0)$

# More NΠ-complete problems: Undirected Hamilton cycle

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Research, in connection to NΠ-complete problems, often is focused on whether the instance being much less complex, may not be solvable by an efficient algorithm
- Alternatively, it can be often shown that even substantially restricted versions of the problem remain NΠ-complete
- Example:  
Undirected Hamilton cycle - the Hamilton cycle problem restricted to graphs that are undirected, that is, symmetric without self-loops

# More NΠ-complete problems: Undirected Hamilton cycle

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: Undirected Hamilton cycle is NΠ-complete
- Proof:
  - we reduce Hamilton cycle to it
  - given a graph  $G \subseteq V \times V$ , and we shall construct a symmetric graph  $G' \subseteq V' \times V'$  without self-loops
  - $G$  has a Hamilton cycle if and only if  $G'$  has one
  - $V' = \{v_0, v_1, v_2 : v \in V\}$  that is,  $G'$  has three nodes  $v_0, v_1, v_2$  for each node  $v$  of  $G$ 
    - $v_0$  is the entry and  $v_2$  is the exit node

# More NPI-complete problems: Undirected Hamilton cycle

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- thus the edges in  $G'$ :  $\{(u_2, v_0), (v_0, u_2) : (u, v) \in G\} \cup \{(v_0, v_1), (v_1, v_0), (v_1, v_2), (v_2, v_1) : v \in V\}$ 
  - that is the nodes  $v_0, v_1, v_2$  are connected by a path in this order, and there is an undirected edge between  $u_2$  and  $v_0$  whenever  $(u, v) \in G$
- now we must prove that  $G'$  has a Hamilton cycle if and only if  $G$  has one
  - suppose that a Hamilton cycle of  $G'$  arrives at node  $v_0$  from an edge of  $(u_2, v_0)$

# More NPI-complete problems: Undirected Hamilton cycle

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- if Hamilton cycle leaves  $v_0$  through an edge other than  $(v_0, v_1)$  then it cannot pick up node  $v_1$  (and it won't be a Hamilton cycle then)
- thus, edge  $(v_0, v_1)$  must be a part of the cycle, and so is  $(v_1, v_2)$
- then the cycle has to continue through one of the edges  $(v_2, w_0)$  where  $(v, w) \in G$
- and then to  $(w_1, w_2)$ , to some  $z_0$ , where  $(w, z) \in G$  etc.

# More NPI-complete problems: Undirected Hamilton cycle

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- therefore the edges of  $(u_0, v_2)$  in the Hamilton cycle of  $G'$  constitute in fact a Hamilton cycle of  $G$
- conversely, any Hamilton cycle  $(v^1, v^2, \dots, v^{|V|})$  of  $G$  can be converted into a Hamilton cycle of  $G'$  as:  
 $(v_0^1, v_1^1, v_2^1, v_0^2, v_2^2, \dots, v_0^{|V|}, v_1^{|V|}, v_2^{|V|})$

# More NΠ-complete problems: Traveling Salesman

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: The Traveling Salesman Problem is NΠ-complete
- Proof:
  - we already know that it is in NΠ
  - we reduce Undirected Hamilton cycle to it
  - given a symmetric graph G
  - we construct an instance of the Traveling Salesman

# More NPI-complete problems: Traveling Salesman

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $n$ , the number of cities,  $|V|$ , the distance  $d_{ij}$  between any two cities  $i$  and  $j$  is
  - $d_{jj} = 0$ , if  $i = j$
  - $d_{jj} = 1$ , if  $(v_i, v_j) \in G$
  - $d_{jj} = 2$ , otherwise
- since  $G$  is a symmetric graph without loops,  $d_{ij} = d_{ji}$  for all cities  $i$  and  $j$
- budget  $B = n$

# More NII-complete problems: Traveling Salesman

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- any tour of the cities has cost equal to the number of  $n$  plus the number of the intercity distances traversed that are not edges of  $G$
- thus, a tour of cost  $B$  or less exists if and only if the tour is a Hamilton cycle of  $G$

# More NΠ-complete problems: Knapsack

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: Knapsack is NΠ-complete
- Proof:
  - we know that Knapsack is in NΠ
  - we reduce Exact Cover to Knapsack
  - given an instance of Knapsack  $a_1, \dots, a_n$ , and  $K$ , a subset  $P$  of  $\{1, \dots, n\}$  such that  $\sum_{i \in P} a_i = K$  can serve as a certificate that the answer to the given instance is yes
  - we are given a universe  $U = \{u_1, \dots, u_n\}$ , and a family  $F = \{S_1, \dots, S_m\}$  of subsets of  $U$

# More NΠ-complete problems: Knapsack

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- we shall construct an instance  $\tau(U, F)$  of Knapsack, that is
  - nonnegative integers  $a_1, \dots, a_k$  and another  $K$  such that,
  - there is a subset  $P \subseteq \{1, \dots, k\}$  with  $\sum_{i \in P} a_i = K$
  - if and only if there is a set of sets  $C \subseteq F$  that are disjoint and collectively cover all of  $U$
- this construction relies on an unexpected relationship between set union and integer addition
  - subsets of a set of  $n$  elements can be represented as strings over  $\{0, 1\}^n$

# More NΠ-complete problems: Knapsack

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- such strings can be interpreted as integers between 0 and  $2^n - 1$ , written in binary
- the union of such sets is the same as adding the corresponding integers
  - the question whether the disjoint union of the teams makes up the whole  $U$  (Exact Cover), seems to be the same as whether there are integers among the given ones that add up to  $K = 1 + 2 + 4 + \dots + 2^{n-1}$
  - the binary number with  $n$  ones

# More NΠ-complete problems: Knapsack

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

$$\begin{aligned}S_1 &= \{u_3, u_4\} \\S_2 &= \{u_2, u_3, u_4\} \\S_3 &= \{u_1, u_2\}\end{aligned}$$

$$\begin{array}{rcl}S_1 &= 0011 & 0011 \\S_2 &= 0111 & 0111 \\S_3 &= 1100 & \underline{1100} \\&& 1111\end{array}$$

- there is only one problem with this correspondence, because in integer addition we may have a carry bit
- e.g. the sum  $11 + 13 + 15 + 24 = 63$ , in binary  $001011 + 001101 + 001111 + 011000 = 111111$   
if we translate back to subsets of  $\{u_1, \dots, u_6\}$ , the sets  $\{u_3, u_5, u_6\}$ ,  $\{u_3, u_4, u_6\}$ ,  $\{u_3, u_4, u_5, u_6\}$  and  $\{u_2, u_3\}$  are neither disjoint, nor do they cover all of  $U$

# More NΠ-complete problems: Knapsack

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- but this problem can be resolved as considering the strings in  $\{0, 1\}^n$  as integers not in binary, but in m-ary
  - m is the number of sets in F
  - we have m integers  $a_1, \dots, a_m$ , where  $a_i = \sum_{uj \in S_i} m^{j-1}$
  - we ask whether there is a subset that adds up to  $K = \sum_{j=1} m^{j-1}$

# More NΠ-complete problems: Knapsack

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének együttes javítása a Pannon Egyetemen

- this way carry is not a problem, because the addition of fewer than  $m$  digits in  $m$ -ary, with each of the digits either 0 or 1, can never result in carry
  - so the resulting instance of Knapsack has a solution if and only if the original instance of Exact Cover has a solution

# More NΠ-complete problems: Two-Machine Scheduling

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Corollary: Partition and Two-Machine Scheduling are NΠ-complete
- Proof:
  - there are polynomial reductions from Knapsack to both of these problems

# More NΠ-complete problems: Independent Set

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: Independent Set is NΠ-complete
- Proof:
  - we know that it is in NΠ
  - we reduce 3-Satisfiability to it
    - we are given a Boolean formula F with clauses  $C_1, \dots, C_m$ , each with at most three literals  
(if a clause has only one or two literals we allow a literal to be repeated to make them three)
    - we construct an undirected graph G and an integer K, such that there is a set of K nodes in G with no edges between them if and only if F is satisfiable

# More NPI-complete problems: Independent Set

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- for each one of the clauses  $C_1, \dots, C_m$  of  $F$ , we have three nodes in  $G$  ( $c_{i1}, c_{i2}, c_{i3}$ ), connected by edges so that they form a triangle
  - these are all the nodes of  $G$  - a total of  $3m$  nodes
  - the goal is  $K = m$
  - for defining the remaining edges of  $G$ , node  $c_{ij}$  is identified with the  $j^{\text{th}}$  literal of clause  $C_i$
  - two nodes are joined by an edge if and only if their literals are the negation of one another
- suppose there is an independent set  $I$  in  $G$  with  $K = m$  nodes

# More NPI-complete problems: Independent Set

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- there is exactly one node in  $I$  from each triangle (since any two nodes from the same triangle are connected by an edge)
- a node in  $I$  means that the corresponding literal is  $T$
- there are no edges between nodes in  $I$  means, that no two such literals are the negation of one another, therefore they can be the basis of a true assignment  $T$
- $T$  may not be fully defined on all variables, for the set of nodes in  $I$  may fail to involve them all
- $T$  may take any truth value on such missing variables, so the resulting truth assignment  $T$  satisfies all clauses

# More NII-complete problems: Independent Set

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- conversely, given any truth assignment satisfying  $F$ , we can obtain an independent set of size  $m$  by picking for each clause a node corresponding to a satisfied literal

# More NΠ-complete problems: Clique and Node Cover

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: Clique and Node Cover are NΠ-complete
- Proof:
  - they are both in NΠ
  - Clique: requires that all edges between any two nodes in the set be present

# More NII-complete problems: Clique and Node Cover

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- we reduce Independent Set to Clique
  - given an instance  $(G, K)$  of Independent Set
    - $G \subseteq V \times V$  is an undirected graph
    - $K \geq 2$  is the goal
  - create  $(G', K')$  Clique
    - $G' = V \times V - \{(i, i) : i \in V\} - G$
    - $K' = K$ , we keep the same goal
  - the maximum independent set of  $G$  is precisely the maximum clique in  $G'$

# More NPI-complete problems: Clique and Node Cover

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

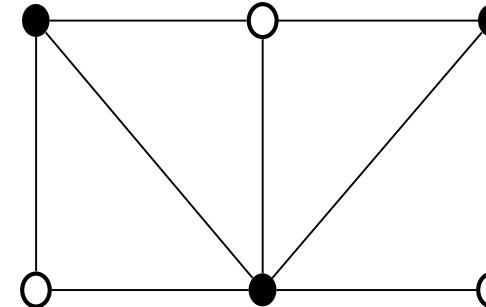
- we reduce Independent Set to Node Cover
- since the nodes in a node cover  $N \subseteq V$  hit all edges, the set  $V - N$  must have no edges between its elements, and thus is an independent set
  - $N \subseteq V$  is a node cover of  $G \leftrightarrow V - N$  is an independent set of  $G$ 
    - the maximum independent set of  $G$  has size  $K$  or more  $\leftrightarrow$  the minimum node cover of  $G$  has size  $|V| - K$  or less
  - the reduction leaves the graph the same, simply replaces  $K$  by  $|V| - K$

# More NPI-complete problems: Clique and Node Cover

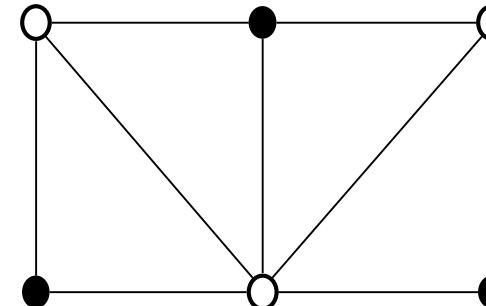
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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Node Cover:
  - the black dots belongs to it



- Independent Set:



# More NII-complete problems: Inequivalence of regular expressions

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A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Remaining question: are two regular expressions (or two NFA) are equivalent?
  - consider the complementary problem
- Given two regular expressions  $R_1$  and  $R_2$ ,  
 $L(R_1) \neq L(R_2)$ ?
  - $w \in (L(R_1) - L(R_2)) \cup (L(R_2) - L(R_1))$  is a certificate
  - the certificate should be polynomial succinct
    - in the length of the two expressions  $|R_1| + |R_2|$
    - Kleene star causes problems

# More NPI-complete problems: Clique and Node Cover

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Definition of \*-free regular expressions: regular expressions over union and concatenation, not containing any occurrences of Kleene star
  - e.g.:  $R = (0 \cup 1)00(0 \cup 1) \cup 010(0 \cup 1)0$
  - if  $x \in L(R) \rightarrow |x| \leq |R|$

# More NΠ-complete problems: Inequivalence of \*-Free Regular Expressions

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Theorem: Inequivalence of \*-Free Regular Expressions is NΠ-complete
- Proof:
  - we already know that the problem is in NΠ
  - we reduce Satisfiability to it
  - given any Boolean formula  $F$  with variables  $x_1, \dots, x_n$  and clauses  $C_1, \dots, C_m$

# More NII-complete problems: Inequialence of \*-Free Regular Expressions

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- we produce  $*$ -free regular expressions  $R_1$  and  $R_2$ , such that:  $L(R_1) \neq L(R_2) \leftrightarrow$  the given Boolean formula is satisfiable
- $R_2$ :  $(0 \cup 1)(0 \cup 1)\dots(0 \cup 1)$ 
  - $(0 \cup 1)$  repeated n times
  - $L(R_2) = \{0, 1\}^n$ , the set of all binary strings of length n

# More NPI-complete problems: Inequialence of \*-Free Regular Expressions

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $R_1: \alpha_1 \cup \alpha_2 \cup \dots \cup \alpha_m$ 
  - $\alpha_i = \alpha_{i,1}\alpha_{i,2} \dots \alpha_{i,n}$ , where
    - $a_{i,j} = 0$ , if  $x_j$  is a literal of  $C_i$
    - $a_{i,j} = 1$ , if  $x_j^C$  is a literal of  $C_i$
    - $a_{i,j} = (0 \cup 1)$ , otherwise
  - strings in  $\{0, 1\}^n$  can be thought of as truth assignments to the Boolean variables  $\{x_1, \dots, x_n\}$

# More NII-complete problems: Inequialence of \*-Free Regular Expressions

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- $L(\alpha_i)$  is precisely the set of all truth assignments that fails to satisfy  $C_i$
- $L(R_1)$  is the set of all truth assignments that fails to satisfy at least one of the clauses of the given Boolean formula
- the set of all truth assignments that satisfies the given Boolean formula is  $L(R_1)^C$
- the given Boolean formula is satisfiable  $\leftrightarrow L(R_1)^C \neq \emptyset \leftrightarrow L(R_1) \neq \{0, 1\}^n = L(R_2)$
- Corollary: the general equivalence problem for RE or NFA can only be harder than the special case

# More NII-complete problems: DFA with minimal number of states

EFOP-3.4.3-16-2016-00009

A felsőfokú oktatás minőségének és hozzáférhetőségének  
együttes javítása a Pannon Egyetemen

- Corollary: Unless  $\Pi = \text{N}\Pi$ , there is no algorithm which (given a RE or a NFA) constructs the minimum-state equivalent DFA in time that is polynomial in the input and in the output

# More NII-complete problems: DFA with minimal number of states

EFOP-3.4.3-16-2016-00009

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- Proof:
  - let  $M_n$  denote the simple  $n + 1$  - state DFA accepting  $\{0, 1\}^n$ 
    - there is no DFA with less state to accept  $\{0, 1\}^n$
  - the formula  $F$  of the previous proof with  $n$  variables is unsatisfiable  $\leftrightarrow L(R_1) = \{0, 1\}^n$ 
    - the minimum-state DFA equivalent to  $R_1$  is exactly  $M_n$
  - indirection: suppose that an algorithm as described in the corollary exists, with a time bound of the form  $p(|x| + |y|)$ , where

# More NII-complete problems: DFA with minimal number of states

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- $p$  is a polynomial
  - $x$  is the input of the algorithm
  - $y$  is the output of the algorithm
- we could solve Satisfiability in polynomial-time:
- given any Boolean formula  $F$  with  $n$  variables
  - first, perform the reduction to obtain a regular expression  $R_1$
  - second, run on input  $R_1$ , the purported algorithm for  $p(|R_1| + |M_n|)$  steps
  - if the algorithm stops within time answer
    - no to Satisfiability, if the output is  $M_n$
    - yes, otherwise

# More NΠ-complete problems: DFA with minimal number of states

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- if the algorithm does not stop within time  $p(|R_1| + |M_n|) \rightarrow$  answer yes
  - the output can not be  $M_n$  any more
- since Satisfiability is NΠ-complete, we must conclude that  $\Pi = N\Pi$ , if the algorithm does really exist

- Once our problem has been shown to be NΠ-complete, what we can do?
  - examine only special cases
  - settle with approximation solutions
  - wait patiently for a solution
    - use stronger PC
- Special cases: examine if we really need to solve this problem in the full generality
  - e.g.: though Satisfiability is NΠ-complete, its special case 2-Satisfiability can be solved efficiently
  - often the special case turns out to be itself NΠ-complete, e.g. 3-Satisfiability



# Introduction to the Theory of Computation

## Lesson 14

### 7.4. Coping with NP-completeness

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- Many NΠ-complete graph problems become trivial when the graph is a tree (it has no cycles)
- The trees have hierarchical structure:
  - pick an arbitrary root
  - let:
    - $T(u)$  the subtree of node  $u$
    - $C(u)$  the children of node  $u$
    - $G(u)$  the grandchildren of node  $u$

- The Independent Set problem is easy when the graph is a tree
  - the size of the largest independent set of the tree can now be found by computing
    - $I(u) = \max \{\sum_{v \in C(u)} I(v), 1 + \sum_{v \in G(u)} I(v)\}$
  - this equation says that, in designing the largest independent set of  $T(u)$  we have two choices:
    - we do not put  $u$  into the independent set
      - we can put together all maximum independent sets in the subtrees of its children

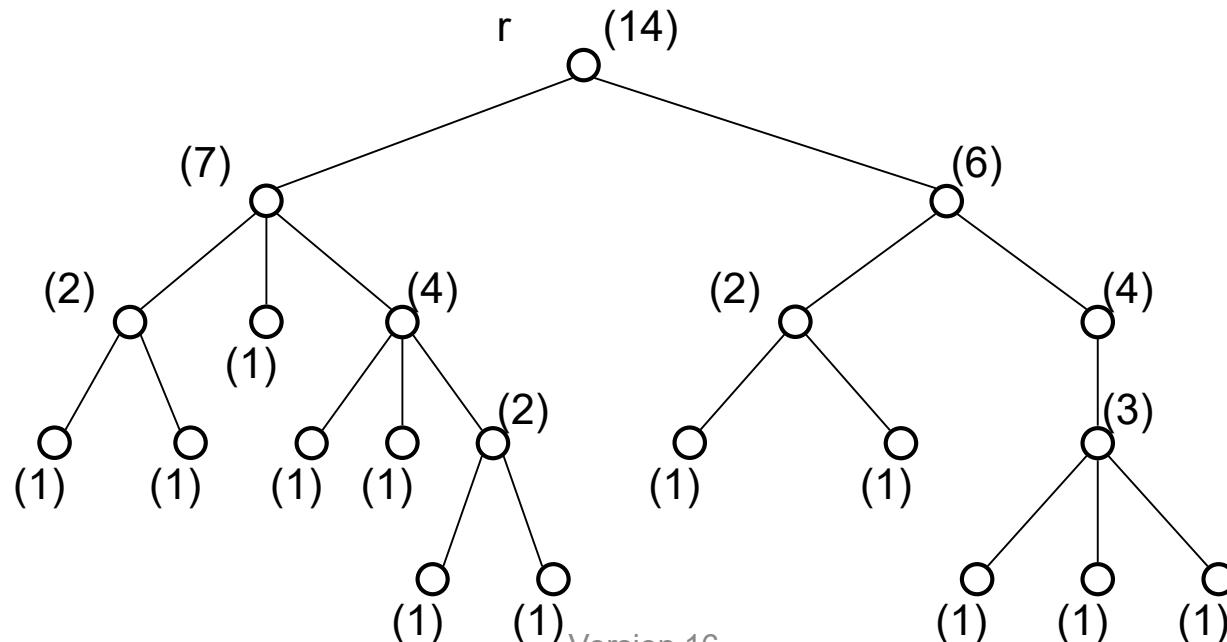
- we put  $u$  in the independent set but we must omit all its children, and assemble the maximum independent sets of the subtrees of all its grandchildren
- dynamic programming can calculate  $I(u)$ 
  - start at the leaves (where  $I(u)$  is trivially one)
  - compute  $I(u)$  for larger and larger subtrees
    - reuse previously calculated  $I(u)$  values
  - the value of  $I$  at the root is the size of the maximum independent set of the tree
- the algorithm is polynomial

# Example

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- The root, denoted  $r$ , has two children and five grandchildren
  - the values of  $I(u)$  are shown in parentheses, the largest independent set of the tree has size 14



- Approximation solutions: use more simple algorithms producing solutions guaranteed to be close to the optimum
- Definition of  $\varepsilon$ -approximation algorithm A:
  - given an NΠ-complete optimization problem
  - for each instance  $x$  the optimal solution is  $\text{opt}(x)$ 
    - assume that  $\text{opt}(x)$  is always a positive integer
  - A is a polynomial algorithm, returning  $A(x)$
  - the following inequality must hold:
$$\frac{|\text{opt}(x) - A(x)|}{\text{opt}(x)} \leq \varepsilon$$

- All NΠ-complete optimization problems can be divided into groups:
  - fully approximable:  $\exists$  algorithm for all  $\varepsilon > 0$ 
    - e.g.: Two Machine Scheduling
  - partly approximable:  $\exists$  algorithm for some range of  $\varepsilon$ 's, but not down to zero
    - e.g.: Node Cover, Max Sat
  - inapproximable: there is no algorithm
    - e.g.: Traveling Salesman, Clique, Independent Set

- Theorem:  $\exists$  1-approximation algorithm for Node Cover

- Proof:

$C := \emptyset$

```
while  $\exists$  an edge  $[u, v]$  in  $G$  do
    add  $u, v$  to  $C$ , delete them from
     $G$ 
```

- it returns a node cover that is at most twice the optimum size

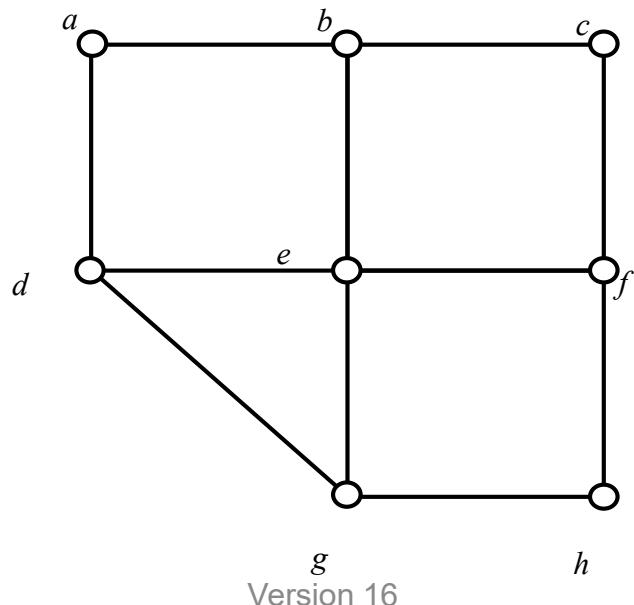
- let  $C$  the node cover returned by the algorithm and  $C'$  be the optimum node cover
- the number of the edges chosen by the algorithm is not larger than  $|C'|$ 
  - at least one end point of each selected edge is in  $C'$  (see the definition of Node Cover)
  - if we remove  $|C'|$  edges from  $G$ , no full edge left
- $|\text{selected edges}| \leq |C'| \rightarrow |C| \leq 2 \cdot |C'|$ 
  - $|C|$  is exactly twice the number of the edges chosen by the algorithm
- it is indeed a 1-approximation algorithm
  - $|C'| = 3, |C| = 6 \rightarrow |3 - 6| / 3 = 1$

# Coping with NΠ-completeness

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- Select
  - (a, b)
  - (e, f)
  - (g, h)



- Theorem: Two-Machine Scheduling is fully approximable
- Proof:
  - the Partition problem can be solved in  $O(nS)$  time
    - $n$  - the number of integers,  $S$  - their sum
  - adapt the algorithm to solve Two-Machine Scheduling (for finding the smallest  $D$ )
    - the last cycle should go until  $S$  (not just up to  $H = \frac{1}{2}S$ )
    - the smallest sum in  $B(n)$  that is greater than or equal to  $\frac{1}{2}S$  is the desired minimum  $D$

- Acceleration:
  - round up these task lengths to the next hundred
  - normalize them by 100
  - we lost a little accuracy but we have decreased the time requirements a hundredfold
- Example:
  - task lengths:
    - 45362, 134537, 85879, 56390, 145627, 197342, 83625, 126789, 38562, 75402 with  $n = 10$ , and  $S \approx 10^6$
  - solving it with the unaccelerated algorithm takes  $10^7$  steps

- rounded numbers: 45400, 134600, 85900, 56400, 145700, 197400, 83700, 126800, 38600, 75500
- normalized numbers: 454, 1346, 859, 564, 1457, 1974, 837, 1268, 386, 755
- thus we can solve this instance in about  $10^5$  steps
- If we round up to the next  $k^{\text{th}}$  power of ten, the absolute error is no more than  $n10^k$ 
  - to calculate the relative error, this quantity must be divided by the optimum, which can be not less than  $S/2$
  - we have a  $2n10^k/S$ -approximation algorithm, with running time  $O(nS/10^k)$
  - if we say  $2n10^k/S = \text{any desirable } \varepsilon > 0$ , we will get running time  $O(n^2/\varepsilon)$  - which is polynomial

- Theorem: Travelling Salesman is inapproximable
- Proof:
  - revisit the reduction from the Hamilton Cycle,  $x$ , to Travelling Salesman,  $\tau(x)$
  - modify the reduction:
    - the distances  $d_{i,j}$  are now the following:
      - $d_{i,j} = 0$ , if  $i = j$
      - $d_{i,j} = 1$ , if  $(v_i, v_j) \in G$
      - $d_{i,j} = 2 + n\epsilon$  otherwise

– this construction means:

- if  $x$  has a Hamilton cycle, then the optimum cost of  $\tau(x)$  is  $n$
- else the optimum cost is greater than  $n - 1 + 2 + n\epsilon + 1 = n(1 + \epsilon)$ 
  - at least one distance  $2 + n\epsilon$  must be traversed, in addition to at least  $n - 1$  others of cost at least 1
  - $+1$  comes from the greater

- assume  $\exists$  a polynomial-time  $\varepsilon$ -approximation algorithm A for the Travelling Salesman problem
- run A on the previously constructed  $\tau(x)$ 
  - if  $A(\tau(x)) \geq n(1 + \varepsilon) + 1 \rightarrow$  the optimum of  $\tau(x)$  can not be n, so x has no Hamilton cycle
    - the relative error of A would be at least:  
$$|n(1 + \varepsilon) + 1 - n| / n > \varepsilon$$
  - if  $A(\tau(x)) \leq n(1 + \varepsilon) \rightarrow$  the optimum solution of  $\tau(x)$  must be n, so x has a Hamilton cycle
    - $\tau(x)$  is designed not to have a tour of cost between  $n + 1$  and  $n(1 + \varepsilon)$

- by applying A on  $\tau(x)$  we can tell whether  $x$  has a Hamilton Cycle or not in polynomial time
- since Hamilton Cycle is NΠ-complete, the assumption is false
- The Travelling Salesman Problem, in which
  - distances  $d_{i,j}$  satisfy the triangle inequality is partly approximable
    - $d_{i,j} \leq d_{i,k} + d_{k,j}$  for  $\forall i, j, k$
    - the best known error bound is  $\frac{1}{2}$
  - cities are restricted to be points on the plane with the usual Euclidean distances, is fully approximable

- All NΠ-complete problems are solvable by polynomially bounded nondeterministic Turing machines
  - we only know exponential methods to stimulate such machines
  - though we can not beat the exponential behavior in the worst case, we can create such algorithms which work more efficiently in average case

- A typical NΠ-complete problem asks, if any member of a large set  $S_0$  of candidate certificates (e.g. permutation of nodes) satisfies certain constraints (e.g. is a Hamilton path?)
  - $|S_0|$  depends exponentially on the size of the problem (e.g. the number of nodes in the graph)

- The solving of an NΠ-complete problem by a nondeterministic TM produces a tree of configurations
  - corresponds some  $S$  to each configurations, a subset of  $S_0$ 
    - $S_0$  belongs to the initial configuration
    - we call these subsets subproblems
    - if the children of configuration  $C$  are  $C_1, C_2, \dots$  and to these configurations correspond  $S, S_1, S_2, \dots$  respectively  $\rightarrow S = S_1 \cup S_2 \cup \dots$
  - to tell if  $S$  has a valid certificate is a problem similar to tell if  $S_0$  has one, but easier as  $|S| < |S_0|$ 
    - this is called self-reducibility

- If we can not tell anything to the original problem → replace it with its subproblems
- This suggests a genre of algorithms for solving NΠ-complete problems:
  - we always maintain a set of active subproblems, A
    - initially  $A = \{S_0\}$
  - at each point we
    - choose a subproblem from A
    - replace it with its smaller subproblems
    - this is called branching

- each newly generated subproblem is submitted to a quick heuristic test, the answers can be:
  - empty: the subproblem has no solution, so it can be omitted
    - this is called backtracking
  - a solution of the original problem contained in the current subproblem → terminates
  - ?: we can not prove either of the previous two
    - add the subproblem in hand to A

```
A := {S0}
while A is not empty do
    choose a subproblem S
    delete S from A
    choose a way of branching out of S, say
        to subproblems S1, ..., Sr
    for each subproblem Si in this list
        do
            if test(Si) find solution then
                halt
            else if test(Si) returns "?" add Si
                to A
    return "no solution"
```

- The backtracking algorithm terminates because, in the end, the subproblems will contain just one candidate solution (a leaf of the tree)
- The effectiveness of the backtracking algorithm depends on three important design decisions:
  - how does one choose the next subproblem?
  - how is the chosen subproblem further split into smaller subproblems
  - which test is used?

- Backtracking algorithm for Satisfiability
  - how to create subproblems: choose a variable  $x$  and fix its value
    - subproblem 1:  $T(x) = \text{true}$ 
      - the clauses in which  $x$  appears are omitted because the whole clause is true
      - $x^C$  is omitted from the clauses in which it appears because its value is fixed
    - subproblem 2:  $T(x) = \text{false}$ 
      - the clauses in which  $x^C$  appears are omitted
      - $x$  is omitted from the clauses in which it appears

- how to choose the variable  $x$  on which to branch:  
select a variable that appears in the smallest clause
  - may soon lead to backtracking
- how to choose the next subproblem: select the subproblem that contains the smallest clause
- how to test a subproblem:
  - if  $\exists$  an empty clause return "subproblem is empty"
  - if there are no clauses return "solution found"
  - otherwise return "?"

- Original problem:

$$(x \vee y \vee z), (x^C \vee y), (y^C \vee z), (z^C \vee x), (x^C \vee y^C \vee z^C)$$

–  $T(x) = \text{true}$ :  $(y), (y^C \vee z), (y^C \vee z^C)$

- $T(y) = \text{true}$ :  $(z), (z^C)$

- $T(z) = \text{true}$ : ()

- $T(z) = \text{false}$ : ()

- $T(y) = \text{false}$ : ()

–  $T(x) = \text{false}$ :  $(y \vee z), (y^C \vee z), (z^C)$

- $T(z) = \text{true}$ : ()

- $T(z) = \text{false}$ :  $(y), (y^C)$

- $T(y) = \text{true}$ : ()

- $T(y) = \text{false}$ : ()

- The problem is unsatisfiable

- Designing a backtracking algorithm for Hamilton cycle
  - we already have: a path with endpoints 'a' and b, going through a set of nodes  $T \subseteq V - \{a, b\}$
  - we are looking for: a Hamilton path from 'a' to b through the remaining nodes in V, to close the Hamilton cycle
  - branching: choose  $c \notin T$  to extend the path by edge  $[a, c]$
  - we do not specify how to choose the subproblem

- the test:
  - if  $G - T - \{a, b\}$  is disconnected or  $G - T$  has a degree-one node other than a or b, return → "subproblem is empty"
  - if  $G - T$  is a path from a to b → return "solution"
  - otherwise return "?"
- at a specific graph with 8 nodes the backtrack algorithm examines only 19 subproblem while total enumeration consider all the leaves  $(n - 1)! = 5040$ 
  - there is no guarantee for the time in the worst case it is still exponential
  - efficiency depends on the branching and on the selection strategies

- An interesting variant of backtracking for optimization problems is called branch and bound
  - each solution has a cost associated with it, and we wish to find the candidate solution in  $S_0$  with the smallest cost

- The branch and bound algorithm:

$A := \{S_0\}$ , bestsofar := inf

while  $A$  is not empty do

    choose a subproblem  $S$ , delete it from  
     $A$

    branch out of  $S$  to subproblems  $S_1, \dots$   
     $S_r$

    for each subproblem  $S_i$  do

        if  $|S_i| = 1 \rightarrow$  update bestsofar

        else if  $\text{lowerbound}(S_i) < \text{bestsofar}$   
             $\rightarrow$  add  $S_i$  to  $A$

return bestsofar

- Operation
  - bestsofar: the smallest cost of any solution seen so far
    - initially infinite
    - it is an upper bound of the optimal solution of the original solution
  - if  $|S_i| = 1 \rightarrow$  the optimal solution of the subproblem is the cost of the single solution in it
    - if this value is smaller than bestsofar  $\rightarrow$  bestsofar is updated
  - lowerbound: a method for obtaining a value which is smaller than the cost of any solution in  $S$ 
    - the sharper (bigger) the bound is the better

- Main idea: if  $\text{lowerbound}(S_i) \geq \text{bestsofar} \rightarrow S_i$  is disregarded because it can not contain such solution whose cost is better than bestsofar

- Main idea: allow a solution of an optimization problem to change a little, and adopt the new solution if it has improved cost
  - let  $S_0$  be the set of candidate solutions in an minimization problem
  - we define a neighbourhood relation  $N \subseteq S_0 \times S_0$ 
    - it captures the intuitive notion changing a little
  - for  $s \in S_0$ , the set  $\{s' : (s, s') \in N\}$  is called the neighbourhood of  $s$

- The algorithm:

```
s := initialsolution
```

```
while ∃ a solution s' such that
```

```
(s, s') ∈ N and cost(s') < cost(s) do:
```

```
    s := s'
```

```
return s
```

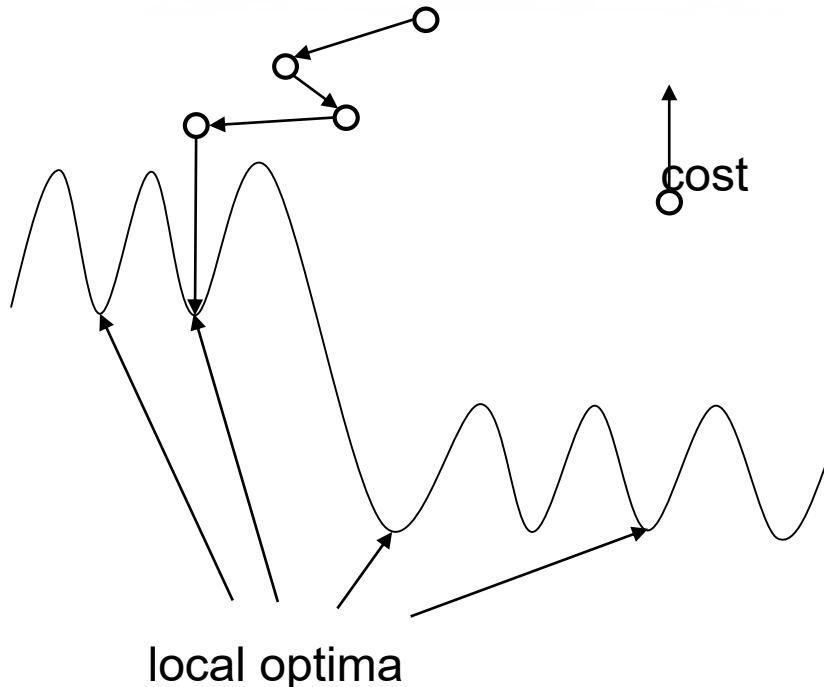
- Operation:
  - the algorithm keeps improving  $s$  by replacing with a neighbour  $s'$  with a better cost, until there is no more
  - we find only a local optimum
  - trade off: larger neighbourhoods
    - probably better local optimum
    - more iteration, so slower algorithm

- decision points:
  - the method used in finding  $s'$ 
    - adopt the first better solution we find in the neighbourhood of  $s$
    - wait to find the best neighbour
  - the procedure initialsolution
    - experience shows it should be randomized
    - after restart a different local optima may be obtained

# Local improvements

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- An optimization problem can be pictured as an energy landscape
  - local optima are depicted as valleys
  - local improvement heuristics jump from solution to solution, until a local optimum is found

- Local improvement algorithm for the Travelling Salesman problem
  - tours are neighboring when they share all but very few links
    - 2-change algorithm: two tours related by  $N \leftrightarrow$  they differ in just two links
    - better results can be achieved by adopting the 3-change
    - 4-change does not seem to be advantageous
  - best heuristic algorithm currently known: Lin-Kernighan algorithm, it relies on  $\lambda$ -change

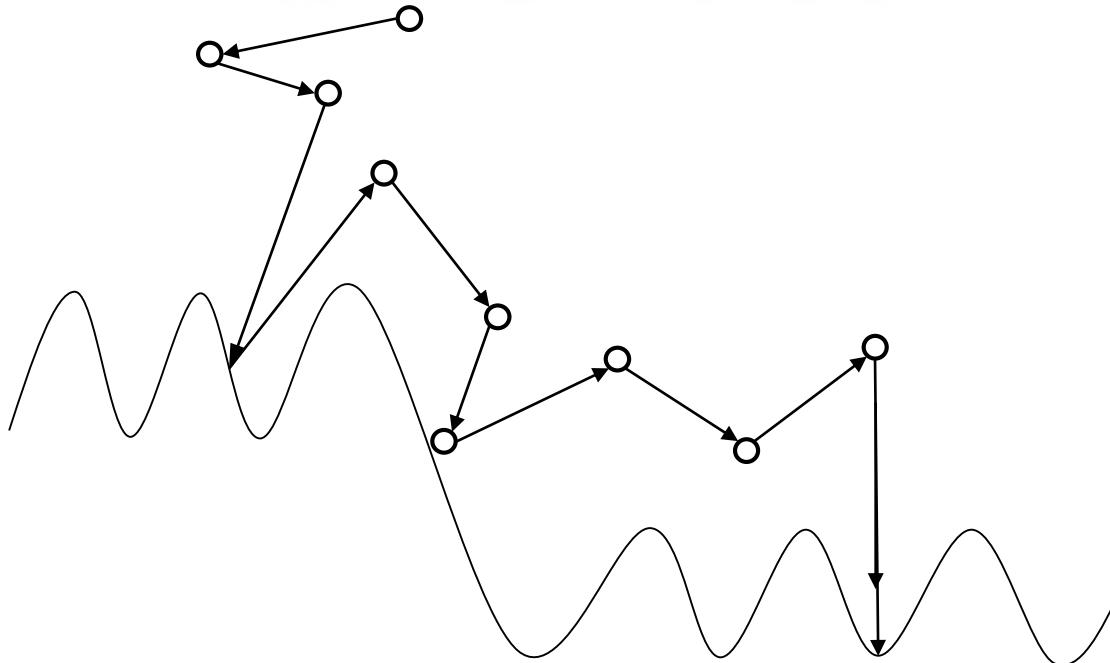
- Local improvement algorithm for Max Set:
  - truth assignments are neighboring when they differ in the value of a single variable
  - experience shows that it is advantageous
    - to adopt as  $s'$  the best neighbour of  $s$  (instead of the first one found that is better than  $s$ )
    - to make lateral moves, that is to adopt a solution even if the new cost is the same as (and not less than) the old one

- Main idea: allow the algorithm to escape from bad local optima by performing occasional cost-increasing changes
  - the inspiration comes from the physics of cooling solids

# Simulated annealing

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- Advantage over the basic local improvement algorithm: because its occasional cost-increasing moves help it avoid early convergence in a bad local optimum
  - this often comes at loss of efficiency

- The algorithm:

```
s := initialsolution, T := T0
```

```
repeat
```

```
    generate randomly s' such that (s,  
    s') ∈ N, let Δ := cost(s') - cost(s)
```

```
    if Δ ≤ 0 then s := s' else
```

```
    s := s' with probability e-Δ/T
```

```
    update(T)
```

```
until T = 0
```

```
return the best solution seen
```

- The probability that a cost increasing change will be adopted depends on
  - the amount of the cost increase  $\Delta$
  - the temperature  $T$ 
    - as time passes,  $T$  decreases as well as the probability to accept a worse solution

- Other local improvement methods:
  - genetic algorithms
  - neural networks
- Local improvement algorithms:
  - they do not in general return the global optimum
    - there is no guaranty for the quality of the returned solution
  - they have exponential worst-case complexity
- But for many NΠ-complete problems, in practice they turn out to be the ones that perform best



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Thank you for your attention!



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