## Reductions and Completeness

## Degrees of Difficulty

- When is a problem more difficult than another?
- B reduces to A if there is a transformation $R$ which for every input $x$ of B yields an equivalent input $R(x)$ of A .
- The answer to $x$ for B is the same as the answer to $R(x)$ for A .
- There must be restrictions on the complexity of computing $R$.
- Otherwise, $R(x)$ might as well solve B .
* E.g., $R(x)=$ "yes" if and only if $x \in \mathrm{~B}$ !


## Degrees of Difficulty (concluded)

- We say problem A is at least as hard as problem B if B reduces to A.
- This makes intuitive sense: If A is able to solve your problem B after only a little bit of work $(R)$, then A must be at least as hard.
- If A were easy, it combined with $R$ (which is also easy) would make B easy, too. ${ }^{\text {a }}$
${ }^{a}$ Thanks to a lively class discussion on October 13, 2009.


## Reduction



Solving problem B by calling the algorithm for problem once and without further processing its answer.

## Comments ${ }^{\text {a }}$

- Suppose B reduces to A via a transformation $R$.
- The input $x$ is an instance of B.
- The output $R(x)$ is an instance of A.
- $R(x)$ may not span all possible instances of $A .{ }^{\mathrm{b}}$
- So some instances of A may never appear in the range of the reduction $R$.
${ }^{\text {a Contributed by Mr. Ming-Feng Tsai (D92922003) on October 29, }}$ 2003.
${ }^{\mathrm{b}} R(x)$ may not be onto; Mr. Alexandr Simak (D98922040) on October 13, 2009.


## Reduction between Languages

- Language $L_{1}$ is reducible to $L_{2}$ if there is a function $R$ computable by a deterministic TM in space $O(\log n)$.
- Furthermore, for all inputs $x, x \in L_{1}$ if and only if $R(x) \in L_{2}$.
- $R$ is said to be a (Karp) reduction from $L_{1}$ to $L_{2}$.
- Note that by Theorem 22 (p. 189), $R$ runs in polynomial time.
- Suppose $R$ is a reduction from $L_{1}$ to $L_{2}$.
- Then solving " $R(x) \in L_{2}$ " is an algorithm for solving " $x \in L_{1}$."


## A Paradox?

- Degree of difficulty is not defined in terms of absolute complexity.
- So a language $\mathrm{B} \in \operatorname{TIME}\left(n^{99}\right)$ may be "easier" than a language $\mathrm{A} \in \operatorname{TIME}\left(n^{3}\right)$.
- This happens when $B$ is reducible to $A$.
- But isn't this a contradiction if the best algorithm for B requires $n^{99}$ steps?
- That is, how can a problem requiring $n^{99}$ steps be reducible to a problem solvable in $n^{3}$ steps?


## A Paradox? (concluded)

- The so-called contradiction does not hold.
- When we solve the problem " $x \in \mathrm{~B}$ ?" via " $R(x) \in \mathrm{A}$ ?", we must consider the time spent by $R(x)$ and its length | $R(x) \mid$.
- If $|R(x)|=\Omega\left(n^{33}\right)$, then answering " $R(x) \in \mathrm{A}$ ?" takes $\Omega\left(\left(n^{33}\right)^{3}\right)=\Omega\left(n^{99}\right)$ steps, which is fine.
- Suppose, on the other hand, that $|R(x)|=o\left(n^{33}\right)$.
- Then $R(x)$ must run in time $\Omega\left(n^{99}\right)$ to make the overall time for answering " $R(x) \in \mathrm{A}$ ?" take $\Omega\left(n^{99}\right)$ steps.
- In either case, the contradiction disappears.


## HAMILTONIAN PATH

- A Hamiltonian path of a graph is a path that visits every node of the graph exactly once.
- Suppose graph $G$ has $n$ nodes: $1,2, \ldots, n$.
- A Hamiltonian path can be expressed as a permutation $\pi$ of $\{1,2, \ldots, n\}$ such that $-\pi(i)=j$ means the $i$ th position is occupied by node $j$. $-(\pi(i), \pi(i+1)) \in G$ for $i=1,2, \ldots, n-1$.
- hamiltonian path asks if a graph has a Hamiltonian path.


## Reduction of hamiltonian path to Sat

- Given a graph $G$, we shall construct a CNF $R(G)$ such that $R(G)$ is satisfiable iff $G$ has a Hamiltonian path.
- $R(G)$ has $n^{2}$ boolean variables $x_{i j}, 1 \leq i, j \leq n$.
- $x_{i j}$ means
the $i$ th position in the Hamiltonian path is occupied by node $j$.



## The Clauses of $R(G)$ and Their Intended Meanings

1. Each node $j$ must appear in the path.

- $x_{1 j} \vee x_{2 j} \vee \cdots \vee x_{n j}$ for each $j$.

2. No node $j$ appears twice in the path.

- $\neg x_{i j} \vee \neg x_{k j}$ for all $i, j, k$ with $i \neq k$.

3. Every position $i$ on the path must be occupied.

- $x_{i 1} \vee x_{i 2} \vee \cdots \vee x_{i n}$ for each $i$.

4. No two nodes $j$ and $k$ occupy the same position in the path.

- $\neg x_{i j} \vee \neg x_{i k}$ for all $i, j, k$ with $j \neq k$.

5. Nonadjacent nodes $i$ and $j$ cannot be adjacent in the path.

- $\neg x_{k i} \vee \neg x_{k+1, j}$ for all $(i, j) \notin G$ and $k=1,2, \ldots, n-1$.


## The Proof

- $R(G)$ contains $O\left(n^{3}\right)$ clauses.
- $R(G)$ can be computed efficiently (simple exercise).
- Suppose $T \models R(G)$.
- From clauses of 1 and 2 , for each node $j$ there is a unique position $i$ such that $T \models x_{i j}$.
- From clauses of 3 and 4 , for each position $i$ there is a unique node $j$ such that $T \models x_{i j}$.
- So there is a permutation $\pi$ of the nodes such that $\pi(i)=j$ if and only if $T \models x_{i j}$.


## The Proof (concluded)

- Clauses of 5 furthermore guarantee that $(\pi(1), \pi(2), \ldots, \pi(n))$ is a Hamiltonian path.
- Conversely, suppose $G$ has a Hamiltonian path

$$
(\pi(1), \pi(2), \ldots, \pi(n))
$$

where $\pi$ is a permutation.

- Clearly, the truth assignment

$$
T\left(x_{i j}\right)=\text { true if and only if } \pi(i)=j
$$

satisfies all clauses of $R(G)$.

## A Comment ${ }^{\text {a }}$

- An answer to "Is $R(G)$ satisfiable?" does answer "Is $G$ Hamiltonian?"
- But a positive answer does not give a Hamiltonian path for $G$.
- Providing witness is not a requirement of reduction.
- A positive answer to "Is $R(G)$ satisfiable?" plus a satisfying truth assignment does provide us with a Hamiltonian path for $G$.

[^0]
## Reduction of REACHABILITY to CIRCUIT VALUE

- Note that both problems are in P.
- Given a graph $G=(V, E)$, we shall construct a variable-free circuit $R(G)$.
- The output of $R(G)$ is true if and only if there is a path from node 1 to node $n$ in $G$.
- Idea: the Floyd-Warshall algorithm.


## The Gates

- The gates are
- $g_{i j k}$ with $1 \leq i, j \leq n$ and $0 \leq k \leq n$.
- $h_{i j k}$ with $1 \leq i, j, k \leq n$.
- $g_{i j k}$ : There is a path from node $i$ to node $j$ without passing through a node bigger than $k$.
- $h_{i j k}$ : There is a path from node $i$ to node $j$ passing through $k$ but not any node bigger than $k$.
- Input gate $g_{i j 0}=$ true if and only if $i=j$ or $(i, j) \in E$.


## The Construction

- $h_{i j k}$ is an AND gate with predecessors $g_{i, k, k-1}$ and $g_{k, j, k-1}$, where $k=1,2, \ldots, n$.
- $g_{i j k}$ is an OR gate with predecessors $g_{i, j, k-1}$ and $h_{i, j, k}$, where $k=1,2, \ldots, n$.
- $g_{1 n n}$ is the output gate.
- Interestingly, $R(G)$ uses no $\neg$ gates: It is a monotone circuit.


## Reduction of CIRCUIT SAT to SAT

- Given a circuit $C$, we will construct a boolean expression $R(C)$ such that $R(C)$ is satisfiable iff $C$ is. $-R(C)$ will turn out to be a CNF.
$-R(C)$ is a depth- 2 circuit; furthermore, each gate has out-degree 1.
- The variables of $R(C)$ are those of $C$ plus $g$ for each gate $g$ of $C$.
- g's propagate the truth values for the CNF.
- Each gate of $C$ will be turned into equivalent clauses.
- Recall that clauses are $\wedge$-ed together by definition.


## The Clauses of $R(C)$

$g$ is a variable gate $x$ : Add clauses $(\neg g \vee x)$ and $(g \vee \neg x)$.

- Meaning: $g \Leftrightarrow x$.
$g$ is a true gate: Add clause $(g)$.
- Meaning: $g$ must be true to make $R(C)$ true.
$g$ is a false gate: Add clause $(\neg g)$.
- Meaning: $g$ must be false to make $R(C)$ true.
$g$ is a $\neg$ gate with predecessor gate $h$ : Add clauses $(\neg g \vee \neg h)$ and $(g \vee h)$.
- Meaning: $g \Leftrightarrow \neg h$.


## The Clauses of $R(C)$ (concluded)

$g$ is a $\vee$ gate with predecessor gates $h$ and $h^{\prime}$ : Add clauses $(\neg h \vee g),\left(\neg h^{\prime} \vee g\right)$, and $\left(h \vee h^{\prime} \vee \neg g\right)$.

- Meaning: $g \Leftrightarrow\left(h \vee h^{\prime}\right)$.
$g$ is a $\wedge$ gate with predecessor gates $h$ and $h^{\prime}$ : Add clauses $(\neg g \vee h),\left(\neg g \vee h^{\prime}\right)$, and $\left(\neg h \vee \neg h^{\prime} \vee g\right)$.
- Meaning: $g \Leftrightarrow\left(h \wedge h^{\prime}\right)$.
$g$ is the output gate: Add clause $(g)$.
- Meaning: $g$ must be true to make $R(C)$ true.

Note: If gate $g$ feeds gates $h_{1}, h_{2}, \ldots$, then variable $g$ appears in the clauses for $h_{1}, h_{2}, \ldots$ in $R(C)$.

## An Example

$$
\begin{aligned}
& \text { ( } \\
& \qquad\left[h_{1} \Leftrightarrow x_{1}\right) \wedge\left(h_{2} \Leftrightarrow x_{2}\right) \wedge\left(h_{3} \Leftrightarrow x_{3}\right) \wedge\left(h_{4} \Leftrightarrow x_{4}\right) \\
& \wedge \quad\left[g_{1} \Leftrightarrow\left(h_{1} \wedge h_{2}\right)\right] \wedge\left[g_{2} \Leftrightarrow\left(h_{3} \vee h_{4}\right)\right] \\
& \left.\left.\wedge \quad\left[g_{5} \Leftrightarrow\left(g_{3} \vee g_{2}\right)\right] \wedge\left(g_{4}\right)\right] \wedge \neg g_{5}\right)
\end{aligned}
$$

## An Example (concluded)

- In general, the result is a CNF.
- The CNF has size proportional to the circuit's number of gates.
- The CNF adds new variables to the circuit's original input variables.


## Composition of Reductions

Proposition 25 If $R_{12}$ is a reduction from $L_{1}$ to $L_{2}$ and $R_{23}$ is a reduction from $L_{2}$ to $L_{3}$, then the composition $R_{12} \circ R_{23}$ is a reduction from $L_{1}$ to $L_{3}$.

- Clearly $x \in L_{1}$ if and only if $R_{23}\left(R_{12}(x)\right) \in L_{3}$.
- How to compute $R_{12} \circ R_{23}$ in space $O(\log n)$, as required by the definition of reduction?


## The Proof (continued)

- An obvious way is to generate $R_{12}(x)$ first and then feeding it to $R_{23}$.
- This takes polynomial time. ${ }^{\text {a }}$
- It takes polynomial time to produce $R_{12}(x)$ of polynomial length.
- It also takes polynomial time to produce $R_{23}\left(R_{12}(x)\right)$.
- Trouble is $R_{12}(x)$ may consume up to polynomial space, much more than the logarithmic space required.

[^1]
## The Proof (concluded)

- The trick is to let $R_{23}$ drive the computation.
- It asks $R_{12}$ to deliver each bit of $R_{12}(x)$ when needed.
- When $R_{23}$ wants to read the $i$ th bit, $R_{12}(x)$ will be simulated until the $i$ th bit is available.
- The initial $i-1$ bits should not be written to the string.
- This is feasible as $R_{12}(x)$ is produced in a write-only manner.
- The $i$ th output bit of $R_{12}(x)$ is well-defined because once it is written, it will never be overwritten by $R_{12}$.


## Completeness ${ }^{\text {a }}$

- As reducibility is transitive, problems can be ordered with respect to their difficulty.
- Is there a maximal element?
- It is not altogether obvious that there should be a maximal element.
- Many infinite structures (such as integers and real numbers) do not have maximal elements.
- Hence it may surprise you that most of the complexity classes that we have seen so far have maximal elements.
${ }^{\mathrm{a}}$ Cook (1971) and Levin (1971).


## Completeness (concluded)

- Let $\mathcal{C}$ be a complexity class and $L \in \mathcal{C}$.
- $L$ is $\mathcal{C}$-complete if every $L^{\prime} \in \mathcal{C}$ can be reduced to $L$.
- Most complexity classes we have seen so far have complete problems!
- Complete problems capture the difficulty of a class because they are the hardest problems in the class.


## Hardness

- Let $\mathcal{C}$ be a complexity class.
- $L$ is $\mathcal{C}$-hard if every $L^{\prime} \in \mathcal{C}$ can be reduced to $L$.
- It is not required that $L \in \mathcal{C}$.
- If $L$ is $\mathcal{C}$-hard, then by definition, every $\mathcal{C}$-complete problem can be reduced to $L .^{\text {a }}$
${ }^{\text {a }}$ Contributed by Mr. Ming-Feng Tsai (D92922003) on October 15, 2003.

Illustration of Completeness and Hardness


## Closedness under Reductions

- A class $\mathcal{C}$ is closed under reductions if whenever $L$ is reducible to $L^{\prime}$ and $L^{\prime} \in \mathcal{C}$, then $L \in \mathcal{C}$.
- P, NP, coNP, L, NL, PSPACE, and EXP are all closed under reductions.


## Complete Problems and Complexity Classes

Proposition 26 Let $\mathcal{C}^{\prime}$ and $\mathcal{C}$ be two complexity classes such that $\mathcal{C}^{\prime} \subseteq \mathcal{C}$. Assume $\mathcal{C}^{\prime}$ is closed under reductions and $L$ is $\mathcal{C}$-complete. Then $\mathcal{C}=\mathcal{C}^{\prime}$ if and only if $L \in \mathcal{C}^{\prime}$.

- Suppose $L \in \mathcal{C}^{\prime}$ first.
- Every language $A \in \mathcal{C}$ reduces to $L \in \mathcal{C}^{\prime}$.
- Because $\mathcal{C}^{\prime}$ is closed under reductions, $A \in \mathcal{C}^{\prime}$.
- Hence $\mathcal{C} \subseteq \mathcal{C}^{\prime}$.
- As $\mathcal{C}^{\prime} \subseteq \mathcal{C}$, we conclude that $\mathcal{C}=\mathcal{C}^{\prime}$.


## The Proof (concluded)

- On the other hand, suppose $\mathcal{C}=\mathcal{C}^{\prime}$.
- As $L$ is $\mathcal{C}$-complete, $L \in \mathcal{C}$.
- Thus, trivially, $L \in \mathcal{C}^{\prime}$.


## Two Important Corollaries

Proposition 26 implies the following.
Corollary $27 P=N P$ if and only if an NP-complete problem in $P$.

Corollary $28 L=P$ if and only if a $P$-complete problem is in $L$.

## Complete Problems and Complexity Classes

Proposition 29 Let $\mathcal{C}^{\prime}$ and $\mathcal{C}$ be two complexity classes closed under reductions. If $L$ is complete for both $\mathcal{C}$ and $\mathcal{C}^{\prime}$, then $\mathcal{C}=\mathcal{C}^{\prime}$.

- All languages $\mathcal{L} \in \mathcal{C}$ reduce to $L \in \mathcal{C}^{\prime}$.
- Since $\mathcal{C}^{\prime}$ is closed under reductions, $\mathcal{L} \in \mathcal{C}^{\prime}$.
- Hence $\mathcal{C} \subseteq \mathcal{C}^{\prime}$.
- The proof for $\mathcal{C}^{\prime} \subseteq \mathcal{C}$ is symmetric.


## Table of Computation

- Let $M=(K, \Sigma, \delta, s)$ be a single-string polynomial-time deterministic TM deciding $L$.
- Its computation on input $x$ can be thought of as a $|x|^{k} \times|x|^{k}$ table, where $|x|^{k}$ is the time bound.
- It is a sequence of configurations.
- Rows correspond to time steps 0 to $|x|^{k}-1$.
- Columns are positions in the string of $M$.
- The $(i, j)$ th table entry represents the contents of position $j$ of the string after $i$ steps of computation.


## Some Conventions To Simplify the Table

- $M$ halts after at most $|x|^{k}-2$ steps.
- The string length hence never exceeds $|x|^{k}$.
- Assume a large enough $k$ to make it true for $|x| \geq 2$.
- Pad the table with $\bigsqcup$ s so that each row has length $|x|^{k}$.
- The computation will never reach the right end of the table for lack of time.
- If the cursor scans the $j$ th position at time $i$ when $M$ is at state $q$ and the symbol is $\sigma$, then the $(i, j)$ th entry is a new symbol $\sigma_{q}$.


## Some Conventions To Simplify the Table (continued)

- If $q$ is "yes" or "no," simply use "yes" or "no" instead of $\sigma_{q}$.
- Modify $M$ so that the cursor starts not at $\triangleright$ but at the first symbol of the input.
- The cursor never visits the leftmost $\triangleright$ by telescoping two moves of $M$ each time the cursor is about to move to the leftmost $\triangleright$.
- So the first symbol in every row is a $\triangleright$ and not a $\triangleright_{q}$.


## Some Conventions To Simplify the Table (concluded)

- Suppose $M$ has halted before its time bound of $|x|^{k}$, so that "yes" or "no" appears at a row before the last.
- Then all subsequent rows will be identical to that row.
- $M$ accepts $x$ if and only if the $\left(|x|^{k}-1, j\right)$ th entry is "yes" for some position $j$.


## Comments

- Each row is essentially a configuration.
- If the input $x=010001$, then the first row is

- A typical row may look like



## Comments (concluded)

- The last rows must look like

- Three out of the table's 4 borders are known:



## A P-Complete Problem

Theorem 30 (Ladner (1975)) CIRCUIT VALUE is $P$-complete.

- It is easy to see that circuit value $\in \mathrm{P}$.
- For any $L \in \mathrm{P}$, we will construct a reduction $R$ from $L$ to CIRCUIT VALUE.
- Given any input $x, R(x)$ is a variable-free circuit such that $x \in L$ if and only if $R(x)$ evaluates to true.
- Let $M$ decide $L$ in time $n^{k}$.
- Let $T$ be the computation table of $M$ on $x$.


## The Proof (continued)

- When $i=0$, or $j=0$, or $j=|x|^{k}-1$, then the value of $T_{i j}$ is known.
- The $j$ th symbol of $x$ or $\bigsqcup$, a $\triangleright$, and a $\bigsqcup$, respectively.
- Recall that three out of $T$ 's 4 borders are known.


## The Proof (continued)

- Consider other entries $T_{i j}$.
- $T_{i j}$ depends on only $T_{i-1, j-1}, T_{i-1, j}$, and $T_{i-1, j+1}$.

| $T_{i-1, j-1}$ | $T_{i-1, j}$ | $T_{i-1, j+1}$ |
| :---: | :---: | :---: |
|  | $T_{i j}$ |  |

- Let $\Gamma$ denote the set of all symbols that can appear on the table: $\Gamma=\Sigma \cup\left\{\sigma_{q}: \sigma \in \Sigma, q \in K\right\}$.
- Encode each symbol of $\Gamma$ as an $m$-bit number, where ${ }^{\text {a }}$

$$
m=\left\lceil\log _{2}|\Gamma|\right\rceil .
$$

${ }^{\text {a }}$ State assignment in circuit design.

## The Proof (continued)

- Let the $m$-bit binary string $S_{i j 1} S_{i j 2} \cdots S_{i j m}$ encode $T_{i j}$.
- We may treat them interchangeably without ambiguity.
- The computation table is now a table of binary entries $S_{i j \ell}$, where

$$
\begin{aligned}
& 0 \leq i \leq n^{k}-1, \\
& 0 \leq j \leq n^{k}-1, \\
& 1 \leq \ell \leq m .
\end{aligned}
$$

## The Proof (continued)

- Each bit $S_{i j \ell}$ depends on only $3 m$ other bits:

$$
\begin{array}{lllll}
T_{i-1, j-1}: & S_{i-1, j-1,1} & S_{i-1, j-1,2} & \cdots & S_{i-1, j-1, m} \\
T_{i-1, j}: & S_{i-1, j, 1} & S_{i-1, j, 2} & \cdots & S_{i-1, j, m} \\
T_{i-1, j+1}: & S_{i-1, j+1,1} & S_{i-1, j+1,2} & \cdots & S_{i-1, j+1, m}
\end{array}
$$

- There is a boolean function $F_{\ell}$ with $3 m$ inputs such that

$$
\begin{aligned}
S_{i j \ell}= & F_{\ell}\left(S_{i-1, j-1,1}, S_{i-1, j-1,2}, \ldots, S_{i-1, j-1, m}\right. \\
& S_{i-1, j, 1}, S_{i-1, j, 2}, \ldots, S_{i-1, j, m} \\
& \left.S_{i-1, j+1,1}, S_{i-1, j+1,2}, \ldots, S_{i-1, j+1, m}\right)
\end{aligned}
$$

where for all $i, j>0$ and $1 \leq \ell \leq m$.

## The Proof (continued)

- These $F_{i}$ 's depend only on $M$ 's specification, not on $x$.
- Their sizes are fixed.
- These boolean functions can be turned into boolean circuits.
- Compose these $m$ circuits in parallel to obtain circuit $C$ with $3 m$-bit inputs and $m$-bit outputs.
- Schematically, $C\left(T_{i-1, j-1}, T_{i-1, j}, T_{i-1, j+1}\right)=T_{i j} .{ }^{\text {a }}$
${ }^{\text {a }} C$ is like an ASIC (application-specific IC) chip.


## Circuit $C$



## The Proof (concluded)

- A copy of circuit $C$ is placed at each entry of the table.
- Exceptions are the top row and the two extreme columns.
- $R(x)$ consists of $\left(|x|^{k}-1\right)\left(|x|^{k}-2\right)$ copies of circuit $C$.
- Without loss of generality, assume the output "yes"/"no" appear at position $\left(|x|^{k}-1,1\right)$.
- Encode "yes" as 1 and "no" as 0 .



## A Corollary

The construction in the above proof yields the following, more general result.

Corollary 31 If $L \in \operatorname{TIME}(T(n))$, then a circuit with $O\left(T^{2}(n)\right)$ gates can decide if $x \in L$ for $|x|=n$.


[^0]:    ${ }^{\text {a }}$ Contributed by Ms. Amy Liu (J94922016) on May 29, 2006.

[^1]:    ${ }^{\text {a }}$ Hence our concern below disappears had we required reductions to be in P instead of L .

