Diskrete Mathematik für Informatik (SS 2025)

Martin Held

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12. Juni 2025



Personalia

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Formalia

LVA-URL (VO+PS): https://www.cosy.sbg.ac.at/~held/teaching/diskrete_

mathematik/dm.html.

Allg. Information: Basis-URL/for_students.html.

PLUSonline: Bitte melden Sie sich unbedingt im PLUSonline zu VO/PS an!

Abhaltezeit der VO: Donnerstag 7⁴⁵–11⁰⁰, mit etwa 20–25 Minuten Pause.

Abhalteort der VO: T01, FB Informatik, Jakob-Haringer Str. 2.

Abhaltezeit des PS: Freitag 11⁴⁰–13⁴⁰.

Abhalteort des PS: T01+T02+T03, Jakob-Haringer Str. 2.

Tutorium: Andreas Auer und Jatin Kumar:

Montag 16⁰⁰–18⁰⁰ (T06), Mittwoch 12³⁰–14³⁰ (T02);

FB Informatik, Jakob-Haringer Str. 2.

Achtung — das Proseminar ist prüfungsimmanent!



Electronic Slides and Online Material

In addition to these slides, you are encouraged to consult the WWW home page of this lecture:

https://www.cosy.sbg.ac.at/~held/teaching/diskrete_mathematik/dm.html.

In particular, this WWW page contains up-to-date information on the course, plus links to online notes, slides and (possibly) sample code.





A Few Words of Warning

I hope that these slides will serve as a practice-minded introduction to various aspects of discrete mathematics which are of importance for computer science. I would like to warn you explicitly not to regard these slides as the sole source of information on the topics of my course. It may and will happen that I'll use the lecture for talking about subtle details that need not be covered in these slides! In particular, the slides won't contain all sample calculations, proofs of theorems, demonstrations of algorithms, or solutions to problems posed during my lecture. That is, by making these slides available to you I do not intend to encourage you to attend the lecture on an irregular basis.



Acknowledgments

These slides are a revised and extended version of a draft prepared by Kamran Safdar. Included is material written by Christian Alt, Caroline Atzl, Michael Burian, Peter Gintner, Bernhard Guillon, Yvonne Höller, Stefan Huber, Sandra Huemer, Christian Lercher, Sebastian Stenger, Alexander Zrinyi. I also benefited from comments and suggestions made by Stefan Huber and Peter Palfrader.

This revision and extension was carried out by myself, and I am responsible for all errors.

Salzburg, February 2025

Martin Held



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Recommended Textbooks I

S. Maurer, A. Ralston.

Discrete Algorithmic Mathematics

A.K. Peters, 3rd edition, Jan 2005; ISBN 978-1-56881-166-6

K.H. Rosen.

Discrete Mathematics and Its Applications McGraw-Hill. 8th edition. 2019: ISBN 9781259676512

🐚 B. Kolman, R.C. Busby, S.C. Ross.

Discrete Mathematical Structures

Pearson India, 6th edition, 2017; ISBN 978-0134696447.

K.A. Ross, C.R.B. Wright.

Discrete Mathematics

Pearson Prentice Hall, 5th edition, Aug 2002; ISBN 9780130652478

C. Stein, R.L.S. Drysdale, K. Bogart. Discrete Mathematics for Computer Science

Addison-Wesley, March 2010; ISBN 978-0132122719.



Recommended Textbooks II

🕒 J. O'Donnell, C. Hall, R. Page.

Discrete Mathematics Using a Computer Springer, 2nd edition, 2006; ISBN 978-1-84628-241-6

N.L. Biggs.

Discrete Mathematics

Oxford University Press, 2nd edition, Feb 2003, reprinted (with corrections) 2008; ISBN 978-0-19-850717-8

M. Smid.

Discrete Structures for Computer Science: Counting, Recursion, and Probability http://cglab.ca/~michiel/DiscreteStructures, 2019

E. Lehman, F.T. Leighton, A.R. Meyer. Mathematics for Computer Science https://courses.csail.mit.edu/6.042, 2018

M.M. Fleck.

Building Blocks for Theoretical Computer Science http://mfleck.cs.illinois.edu/building-blocks/, 2017



Table of Content

- Introduction
- Propositional and Predicate Logic
- Definitions and Theorem Proving
- Numbers and Basics of Number Theory
- 5 Principles of Elementary Counting and Combinatorics
- **6** Complexity Analysis and Recurrence Relations
- Graph Theory
- 8 Cryptography



- Introduction
 - What is Discrete Mathematics?
 - Motivation



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- No universally accepted definition of the scope of DM exists . . .
- Typically, objects studied in DM can only assume discrete, separate values rather than values out of a continuum; sets of such objects are countable.



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 - Number theory,
 - Proofs and mathematical reasoning,
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Applications of Discrete Mathematics

 DM forms the mathematical language of computer science. It is at the very heart of several other parts of computer science.



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- We start with a set of sample problems; solutions for all problems will be worked out or, at least, sketched during this course.



- Introduction
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• An inspection of the numbers on the right-hand side might let us suspect that

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• An inspection of the numbers on the right-hand side might let us suspect that

$$1+2+3+\cdots+(n-1)+n=\frac{n(n+1)}{2}$$
.

But is this indeed correct? And, by the way, what do the dots in this equation really mean??

 An answer can be established by means of number theory (natural numbers, induction). And we get indeed

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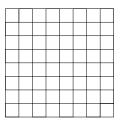
$$1+2+3+\cdots+(n-1)+n=\frac{n(n+1)}{2}$$

for all "natural numbers" n.

- Caution: Even after calculating this sum for all values of n between 1 and 500 one can not legitimately claim to know the sum for, say, n := 1000.
- Note: It would constitute a horrendous waste of CPU time to let a computer compute $1+2+3+\cdots+(n-1)+n$ by successively adding numbers if we could simply obtain the result by evaluating $\frac{n(n+1)}{2}$.

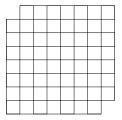


• Consider an 8 × 8 chessboard



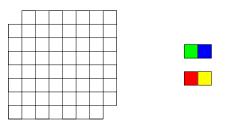


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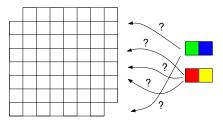


 Consider an 8 × 8 chessboard with the upper-left and lower-right cells removed, and assume that we are given red/yellow and green/blue domino blocks whose sizes match the size of two adjacent squares of the chessboard.



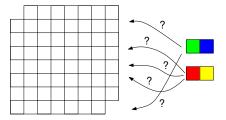


- Consider an 8 × 8 chessboard with the upper-left and lower-right cells removed, and assume that we are given red/yellow and green/blue domino blocks whose sizes match the size of two adjacent squares of the chessboard.
- Question: Can this chessboard be covered completely by 31 domino blocks of arbitrary color combinations?





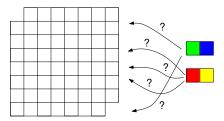
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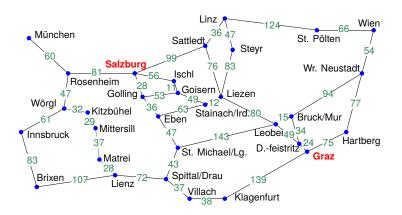


- We consult counting principles and obtain the answer: No!
- Caution: Simply trying out *all* possible placements of domino blocks hardly is an option for an 8×8 chessboard and definitely no option for an $n \times n$ board!



Sample Problem: Route Calculation

Question: What is the shortest route for driving from Salzburg to Graz?

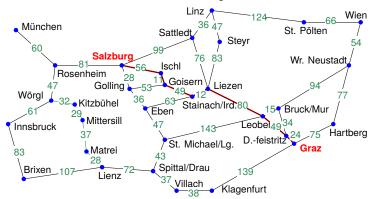




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Sample Problem: Route Calculation

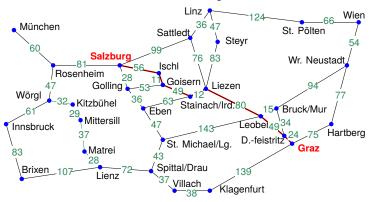
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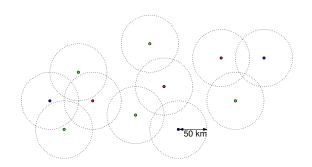


 Note: Simply trying all possible routes gets tedious! (How would you even guarantee that all possible routes have indeed been checked?)

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Sample Problem: Channel Assignment

Suppose that frequencies out of a set of m frequencies are to be assigned to n
broadcast stations within Austria. We are told that the area serviced by a station
lies within a disk with radius 50 kilometers. Obviously, no two different stations
whose broadcast areas overlap may use the same frequency.

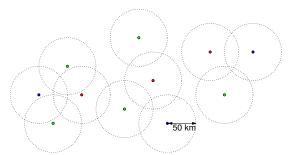




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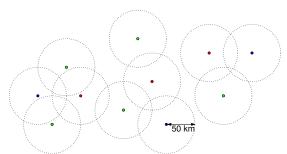
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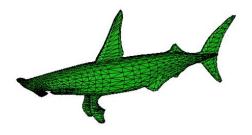
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 The solution can be obtained by using techniques of computational geometry combined with graph coloring.

Sample Problem: Memory Required for Storing a Polyhedron

• Suppose that a polyhedral model has *n* vertices. How many edges and faces can it have at most? What is the storage complexity relative to *n*?

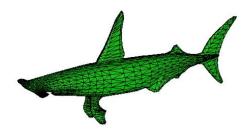




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Sample Problem: Memory Required for Storing a Polyhedron

• Suppose that a polyhedral model has *n* vertices. How many edges and faces can it have at most? What is the storage complexity relative to *n*?



• Answer provided by graph theory: A polyhedron with n vertices has at most 3n-6 edges and 2n-4 faces.





input:		10	Ю
10.0			



input:	100
after round 1:	75



input:	100
after round 1:	75
after round 2:	56



input:	100
after round 1:	75
after round 2:	56
after round 3:	42



Suppose that an algorithm is given n numbers as input and that it solves a
problem by proceeding as follows: During one round of computation, it performs
n computational steps. We know that during each round it discards at least 25%
of the numbers. The algorithm executes one round after the other until only one
number is left.

input:	100
after round 1:	75
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 Question: How many rounds does the algorithm run in the worst case (depending on the input size n)? How many computational steps are carried out in the worst case?

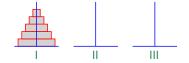


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- Question: How many rounds does the algorithm run in the worst case (depending on the input size n)? How many computational steps are carried out in the worst case?
- Answer provided by the theory of recurrence relations: The number of computational steps is linear in n, and the number of rounds is logarithmic in n.
- In asymptotic notation: O(n) and $O(\log n)$.

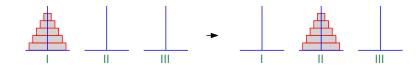


• Tower-of-Hanoi Problem (ToH): Given three pegs (labeled I,II,III) and a stack of *n* disks arranged on Peg I from largest at the bottom to smallest at the top,



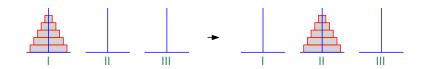


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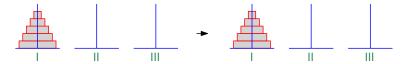


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- One can also prove: Every(!) algorithm that solves ToH needs at least 2ⁿ 1 moves.
- Thus, the solution achieved by the recursive algorithm is optimal as far as the number of moves is concerned.
- [Buneman&Levy (1980)]: There exists a simple iterative solution that avoids are exponential-sized stack!

 According to legend, the power of exponential growth was already known by the Brahmin Sissa ibn Dahir (ca. 300-400 AD): As a reward for the invention of the game of chess (or its Indian predecessor Chaturanga) he asked his king to place one grain of rice in the first square of a chessboard, two in the second, four in the third, and so on, doubling the amount of rice up to the 64-th square.



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- So, how many grains of rice did Sissa ask for?



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- [Sagan 1997]: "Exponentials can't go on forever, because they will gobble up everything".
- The "second half of the chessboard" is a phrase, coined by Kurzweil in 1999 to refer to the point where exponential growth begins to have a significant impact.

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 - Share in parts?!



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Sample Problem: Key Distribution and Message Encryption

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- Answer provided by cryptography: The Diffie-Hellman Algorithm provides a simple way to exchange a key via public communication channels.
- By the way, how could Alice and Bob encrypt or decrypt messages once they have exchanged their key?
- Answer: This is yet another application of number theory!



- 2 Propositional and Predicate Logic
 - Propositional Logic
 - Predicate Logic
 - Special Quantifiers



- 2
 - **Propositional and Predicate Logic**
 - Propositional Logic
 - Predicate Logic
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Propositional Logic

- Goal: specification of a language for formally expressing theorems and proofs.
- Aka: propositional calculus, logic of statements, statement logic;
- Dt.: Aussagenlogik.

Definition 1 (Proposition, Dt.: Aussage)

A *proposition* is a statement that is either true or false.

- Propositions can be atomic, like "The sun is shining", or compound,
 - like "The sun is shining and the temperature is high".
- In the latter case, the proposition is a composition of atomic or compound propositions by means of logical junctors. (Junctors are also known as connectives or operators.)



Language of Propositional Logic

Definition 2 (Propositional formula, Dt.: aussagenlogische Formel)

A propositional formula is constructed inductively from a set of

- propositional variables (typically p, q, r or $p_1, p_2, ...$);
- junctors: ¬, ∧, ∨, ⇒, ⇔;
- parentheses: (,);
- constants (truth values): \bot , \top (or F, T);

based on the following rules:

- A propositional variable is a propositional formula.
- ullet The constants \bot and \top are propositional formulas.
- If ϕ_1 and ϕ_2 are propositional formulas then so are the following:

$$(\neg \phi_1), (\phi_1 \land \phi_2), (\phi_1 \lor \phi_2), (\phi_1 \Rightarrow \phi_2), (\phi_1 \Leftrightarrow \phi_2).$$



Precedence Rules

 Precedence rules (Dt.: Vorrangregeln) are used frequently to avoid the burden of too many parentheses. From highest to lowest precedence, the following order is common.

$$\neg$$
, \wedge , \vee , \Leftrightarrow

- Unfortunately, different precedence rules tend to be used by different authors.
- Thus, make it clear which order you use, or in case of doubt, insert parentheses!
- It is common to represent the truth values of a proposition under all possible assignments to its variables by means of a *truth table*.
- In addition to the standard junctors we also define two other operators, Nand, denoted by ↑ (or sometimes by |), and Non, denoted by ↓.



Truth Tables

р	q	$\neg p$	$p \wedge q$	$p \vee q$	$p \Rightarrow q$	$p \Leftrightarrow q$	$p \uparrow q$	$p \downarrow q$
T	T	F	T	T	T	T	F	F
Τ	F	F	F	T	F	F	T	F
F	T	T	F	T	T	F	T	F
F	F	T	F	F	T	T	T	T

- Common names for the junctors in natural language:
 - ¬p: NoT, negation;
 - p ∧ q: And, conjunction;
 - p ∨ q: OR, disjunction;
 - p ⇒ q: IMPLIES, conditional, if p then q, q if p, p sufficient for q, q necessary for p;
 - p ⇔ q: IFF, equivalence, biconditional, p if and only if q, p necessary and sufficient for q.
- Note: The truth table (Dt.: Wahrheitstabelle) of a formula with n variables has 2ⁿ rows.



Tautologies, Contradictions

Definition 3 (Tautology, Dt.: Tautologie)

A propositional formula is a *tautology* if it is true under all truth assignments to its variables.

Definition 4 (Contradiction, Dt.: Widerspruch)

A propositional formula is a *contradiction* if it is false under all truth assignments to its variables.

- Standard examples: $(p \vee \neg p)$ and $(p \wedge \neg p)$.
- Easy to prove: The negation of a tautology yields a contradiction, and vice versa.



Logical Equivalence

Definition 5 (Logical equivalence, Dt.: logische Äquivalenz)

Two propositional formulas are *logically equivalent* if they have the same truth table. Logical equivalence of formulas ϕ_1, ϕ_2 is commonly denoted by $\phi_1 \equiv \phi_2$.

Theorem 6

Two propositional formulas ϕ_1, ϕ_2 are logically equivalent iff $\phi_1 \Leftrightarrow \phi_2$ is a tautology.

Definition 7 (Complete set of junctors, Dt.: vollständige Junktorenmenge)

A set S of junctors is said to be *complete* (or truth-functionally adequate/complete) if, for any given propositional formula, a logically equivalent one can be written using only junctors of S.

Note: The sets {↑} and {↓} both are complete sets of junctors.



Laws for Logical Equivalence

Theorem 8

Let ϕ_1, ϕ_2 be propositional formulas. Then the following equivalences hold:

Identity:
$$\phi_1 \land T \equiv \phi_1$$
 $\phi_1 \lor F \equiv \phi_1$
Domination: $\phi_1 \lor T \equiv T$ $\phi_1 \land F \equiv F$

Idempotence:
$$\phi_1 \lor \phi_1 \equiv \phi_1$$
 $\phi_1 \land \phi_1 \equiv \phi_1$

Double negation:
$$\neg \neg \phi_1 \equiv \phi_1$$

Commutativity:
$$\phi_1 \wedge \phi_2 \equiv \phi_2 \wedge \phi_1$$
 $\phi_1 \vee \phi_2 \equiv \phi_2 \vee \phi_1$

$$\phi_1 \Leftrightarrow \phi_2 \equiv \phi_2 \Leftrightarrow \phi_1$$

Distributivity:
$$(\phi_1 \lor \phi_2) \land \phi_3 \equiv (\phi_1 \land \phi_3) \lor (\phi_2 \land \phi_3)$$

$$(\phi_1 \land \phi_2) \lor \phi_3 \equiv (\phi_1 \lor \phi_3) \land (\phi_2 \lor \phi_3)$$

Associativity:
$$(\phi_1 \lor \phi_2) \lor \phi_3 \equiv \phi_1 \lor (\phi_2 \lor \phi_3)$$

$$(\phi_1 \wedge \phi_2) \wedge \phi_3 \equiv \phi_1 \wedge (\phi_2 \wedge \phi_3)$$

De Morgan's laws:
$$\neg(\phi_1 \land \phi_2) \equiv \neg\phi_1 \lor \neg\phi_2$$

$$\neg(\phi_1 \lor \phi_2) \equiv \neg\phi_1 \land \neg\phi_2$$

Trivial tautology:
$$\phi_1 \lor \neg \phi_1 \equiv T$$

Trivial contradiction: $\phi_1 \land \neg \phi_1 \equiv F$

Contraposition:
$$\neg \phi_1 \Leftrightarrow \neg \phi_2 \equiv \phi_1 \Leftrightarrow \phi_2 \quad \neg \phi_2 \Rightarrow \neg \phi_1 \equiv \phi_1 \Rightarrow \phi_2$$

Implication as Disj.:
$$\phi_1 \Rightarrow \phi_2 \equiv \neg \phi_1 \lor \phi_2$$



Logical Implication and Proofs

Definition 9 (Logical implication, Dt.: logische Implikation)

A formula ϕ_1 logically implies ϕ_2 , denoted by $\phi_1 \models \phi_2$, if $\phi_1 \Rightarrow \phi_2$ is a tautology.

Definition 10 (Proof, Dt.: Beweis)

A *proof* of ψ based on premises ϕ_1,\ldots,ϕ_n is a finite sequence of propositions that ends in ψ such that each proposition is either a premise or a logical implication of the previous proposition.

- Note: Logical implication rather than logical equivalence!
- Thus,
 - note that it need not be possible to revert a proof!
 - pay close attention to which steps are actual equivalences if you intend to argue both ways!



Rules of Inference

- Aka: proof rules (Dt.: Schlußregeln).
- In addition to the following inference rules for propositional formulas ϕ_1, ϕ_2 , all the equivalence rules apply: Each equivalence can be written as two inference rules since they are valid in both directions.



Satisfiability

Definition 11 (Satisfiability, Dt.: Erfüllbarkeit)

A formula ϕ is *satisfiable* if there exists at least one truth assignment to the variables of ϕ that makes ϕ true.

Definition 12 (Satisfiability equivalent)

Two formulas are *satisfiability equivalent* if both formulas are either satisfiable or not satisfiable.



Conjunctive Normal Form

 In mathematics, normal forms are canonical representations of objects such that all equivalent objects have the same representation.

Definition 13 (Literal, Dt.: Literal)

A literal is a propositional variable or the negation of a propositional variable. A clause is a disjunction of literals.

• E.g., if p, q are variables then p and $\neg q$ are literals, and $(p \lor \neg q)$ is a clause.

Definition 14 (Conjunctive normal form, Dt.: konjunktive Normalform)

A propositional formula is in (general) conjunctive normal form (CNF) if it is a conjunction of clauses.

• E.g., $\neg p_1 \land (p_2 \lor p_5 \lor \neg p_6) \land (\neg p_3 \lor p_4 \lor \neg p_6)$ is a CNF formula.

Definition 15 (k-CNF)

A CNF formula is a k-CNF formula if every clause contains at most k literals.



Conjunctive Normal Form

- Note: Some textbooks demand *exactly k literals* rather than *at most k literals*.
- Note: It is common to demand that no variable may appear more than once in a clause.
- Note: For k ≥ 3, a general CNF formula can easily be converted in polynomial time (in the number of literals) into a k-CNF formula with exactly k literals per clause such that no variable appears more than once in a clause and such that the two formulas are satisfiability equivalent.



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Predicate Logic

Definition 16 (*n*-ary Relation, Dt.: *n*-stellige Relation)

Let A_1, A_2, \ldots, A_n be sets, for some $n \in \mathbb{N}$. An *n-ary relation* \mathcal{R} on A_1, A_2, \ldots, A_n is a subset of their Cartesian product, i.e., $\mathcal{R} \subseteq A_1 \times A_2 \times \cdots \times A_n$.

Definition 17 (*n*-ary Function, Dt.: *n*-stellige Funktion)

Let A_1,A_2,\ldots,A_n,B be sets, for some $n\in\mathbb{N}$. An n-ary function $\mathcal F$ from $A_1\times A_2\times\cdots\times A_n$ to B is an (n+1)-ary relation on A_1,A_2,\ldots,A_n,B such that for any $(a_1,a_2,\ldots,a_n)\in A_1\times A_2\times\cdots\times A_n$ there exists a unique $b\in B$ such that $(a_1,a_2,\ldots,a_n,b)\in\mathcal F$.

- It is common to write $y = \mathcal{F}(a_1, \ldots, a_n)$ for "pick y such that $(a_1, \ldots, a_n, y) \in \mathcal{F}$ ".
- The set $A_1 \times A_2 \times \cdots \times A_n$ is called the *domain* and the set B is called the codomain of \mathcal{F} .
- An n-ary relation/function over a set A is a relation/function where $A_1 = A_2 = \ldots = A_n = A$, i.e., $A_1 \times A_2 \times \cdots \times A_n = A^n$. It is also called an n-place relation/function.
- A 1-ary relation/function is called *unary*, and a 2-ary relation/function is called *binary*.

Predicate Logic: Predicates

Definition 18 (Predicate, Dt.: Prädikat)

For an *n*-ary relation \mathcal{R} over A, an *n*-ary predicate over A is the *n*-ary function $f_{\mathcal{R}}: A^n \to \{T, F\}$, where

$$f_{\mathcal{R}}(a_1,\ldots,a_n):=\left\{egin{array}{ll} T & \mbox{if } (a_1,\ldots,a_n)\in\mathcal{R}, \\ F & \mbox{otherwise}. \end{array}
ight.$$

- Thus, a predicate is a Boolean function.
- Note: This is a slight abuse of notation since the symbols ":" and "→" in " $f: M \to N$ " actually form already a 3-ary predicate!
- An 1-ary predicate is called unary, and a 2-ary predicate is called binary.
- ullet A sample unary predicate on $\mathbb R$ is

"x is non-negative" :=
$$\left\{ \begin{array}{ll} T & \text{if } x \geqslant 0, \\ F & \text{otherwise.} \end{array} \right.$$

Dt.: Prädikatenlogik.



Language of Predicate Logic

Definition 19 (Predicate vocabulary, Dt.: Symbolmenge)

A predicate vocabulary consists of

- a set C of constant symbols,
- a set F of function symbols,
- a set V of variables, typically $\{x_1, x_2, \ldots\}$ or $\{a, b, \ldots\}$,
- a set P of predicate symbols, including the 0-ary predicate symbols (truth values)
 ⊥, ⊤ or F, T,

together with

- logical junctors \neg , \land , \lor , \Rightarrow , \Leftrightarrow ,
- quantifiers ∃, ∀,
- parentheses.



Language of Predicate Logic

Definition 20 (Term)

A *term* over (C, V, F) is defined inductively as follows:

- Every constant $c \in C$ is a term.
- Every variable $x \in \mathcal{V}$ is a term.
- If t_1, \ldots, t_n are terms and f is an n-ary function symbol then $f(t_1, \ldots, t_n)$ is a term.
- ullet Note: Constants can be thought of as 0-ary function symbols. Thus, a set $\mathcal C$ of constants need not be considered when defining the language of predicate logic.



Language of Predicate Logic

Definition 21 (Formulas)

The set of *formulas* over (C, V, F, P) is defined inductively as follows:

- \bullet \perp and \top are formulas.
- If t_1, \ldots, t_n are terms and $P \in \mathcal{P}$ is an n-ary predicate, then $P(t_1, \ldots, t_n)$ is a (so-called atomic) formula.
- If ϕ and ψ are formulas then $(\neg \phi)$, $(\phi \land \psi)$, $(\phi \lor \psi)$, $(\phi \Rightarrow \psi)$ and $(\phi \Leftrightarrow \psi)$ are formulas.
- If ϕ is a formula then $(\forall x \ \phi)$ and $(\exists x \ \phi)$ are formulas. In both cases, the *scope* of the quantifier is given by the formula ϕ to which the quantifier is applied.

Definition 22 (Quantifier-free formula, Dt.: quantorenfreie Formel)

A quantifier-free formula is a formula which does not contain a quantifier.



Quantifiers

Definition 23 (Universe of discourse, Dt.: Wertebereich, Universum)

The *universe of discourse* specifies the set of values that the variable x may assume in $(\forall x \ \phi)$ and $(\exists x \ \phi)$.

Definition 24 (Universal quantifier, Dt.: Allquantor)

 $(\forall x \ P(x))$ is the statement

"P(x) is true for all x (in the universe of discourse)".

Definition 25 (Existential quantifier, Dt.: Existenzquantor)

 $(\exists x \ P(x))$ is the statement "there exists x (in the universe of discourse) such that P(x) is true".

- The notation $(\exists !x \ P(x))$ is a convenience short-hand for "there exists exactly one x such that P(x) is true",
 - i.e., for denoting existence and uniqueness of a suitable \boldsymbol{x} .



Precedence Rules for Quantified Formulas

- No universally accepted precedence rule exists.
- Thus, you have to make your specific order very clear.
- Even better, use parentheses or (significant!) spaces between coherent parts of the expression.
- First-order logic versus higher-order logic: In first-order predicate logic, predicate quantifiers or function quantifiers are not permitted, and variables are the only objects that may be quantified. Also, predicates are not allowed to have predicates as arguments.



Definition 26 (Free variables, Dt.: freie Variable)

The *free variables* of a formula ϕ or a term t, denoted by $FV(\phi)$ and FV(t), are defined inductively as follows:

```
FV(c) := \{\};
For a constant c \in C:
For a variable x \in \mathcal{V}:
                                                         FV(x) := \{x\}:
                                 FV(f(t_1,\ldots,t_n)) := FV(t_1) \cup \ldots \cup FV(t_n);
For a term f(t_1, \ldots, t_n):
For a formula P(t_1, \ldots, t_n): FV(P(t_1, \ldots, t_n)) := FV(t_1) \cup \ldots \cup FV(t_n);
Also.
                                                        FV(\bot) := \{\}.
                                                        FV(\top) := \{\};
                                                    FV((\neg \phi)) := FV(\phi),
For formulas \phi and \psi:
                                                FV((\phi \wedge \psi)) := FV(\phi) \cup FV(\psi),
                                                FV((\phi \vee \psi)) := FV(\phi) \cup FV(\psi),
                                               FV((\phi \Rightarrow \psi)) := FV(\phi) \cup FV(\psi),
                                               FV((\phi \Leftrightarrow \psi)) := FV(\phi) \cup FV(\psi):
                                                 FV((\forall x \ \phi)) := FV(\phi) \setminus \{x\},
For a formula \phi:
                                                 FV((\exists x \ \phi)) := FV(\phi) \setminus \{x\}.
```

Definition 27 (Bound variables, Dt.: gebundene Variable)

The bound variables of a formula ϕ or a term t, denoted by $BV(\phi)$ and BV(t), are defined inductively as follows:

```
For a constant c \in \mathcal{C}:
                                                        BV(c) := \{\}:
For a variable x \in \mathcal{V}:
                                                        BV(x) := \{\}:
For a term f(t_1, ..., t_n): BV(f(t_1, ..., t_n)) := \{\};
For a formula P(t_1, \ldots, t_n): BV(P(t_1, \ldots, t_n)) := \{\};
Also.
                                                       BV(\bot) := \{\},\
                                                       BV(\top) := \{\};
For formulas \phi and \psi:
                                                  BV((\neg \phi)) := BV(\phi),
                                               BV((\phi \wedge \psi)) := BV(\phi) \cup BV(\psi),
                                               BV((\phi \vee \psi)) := BV(\phi) \cup BV(\psi),
                                              BV((\phi \Rightarrow \psi)) := BV(\phi) \cup BV(\psi).
                                              BV((\phi \Leftrightarrow \psi)) := BV(\phi) \cup BV(\psi):
                                                BV((\forall x \ \phi)) := BV(\phi) \cup \{x\},\
For a formula \phi:
                                                BV((\exists x \ \phi)) := BV(\phi) \cup \{x\}.
```

Free and Bound Variables

- Note: Technically speaking, one variable symbol may denote both a free and a bound variable of a formula!
- However, common sense dictates to use a different symbol if a different variable is meant, even if not required by the syntax of predicate logic:
 - Do not use the same symbol for bound and free variables! E.g.,

$$(P(x) \Rightarrow (\forall x \ Q(x)))$$

is syntactically correct but extremely difficult to parse for a human.

Also, do not re-use symbols of bound variables inside nested quantifiers!
 E.g.,

$$(\forall x \ (P(x) \Rightarrow (\forall x \ Q(x))))$$

is syntactically correct but horrible to parse.

Definition 28 (Sentence, Dt.: geschlossener Ausdruck)

A formula ϕ is a *sentence* if $FV(\phi) = \{\}$.



Substitutions

Definition 29 (Substitution, Dt.: Ersetzung)

For a formula ϕ , variable x and term t, we obtain the *substitution* of x by t, denoted as $\phi[t/x]$, by replacing each free occurrence of x in ϕ by t.

Definition 30 (Valid substitution, Dt.: gültige Ersetzung)

A substitution of t for x in a formula ϕ is valid if and only if no free variable of t ends up being bound in $\phi[t/x]$.

- Not a valid substitution of x: $\phi \equiv (\exists y \in \mathbb{N} \ y > 10 \ \land \ x < y)$ and t := 2y + 5.
- Again, it is very poor practice to substitute x by t if t contains any variable that also is a bound variable of ϕ !

$$\phi \equiv (\forall z \in \mathbb{N} \ z^2 > 0) \quad \lor \quad (\exists y \in \mathbb{N} \ y > 10 \quad \land \quad x < y) \text{ and } t := 2z + 5.$$



Equivalence Rules

Theorem 31

Let x be a variable, and ϕ and ψ be formulas which normally contain x as a free variable. Then the following equivalences hold:

De Morgan's laws: $(\neg(\forall x \ \phi)) \equiv (\exists x \ (\neg\phi))$

$$(\neg(\exists x \ \phi)) \equiv (\forall x \ (\neg\phi))$$

Trivial conjunction: $(\forall x \ (\phi \land \psi)) \equiv ((\forall x \ \phi) \land (\forall x \ \psi))$

Only if
$$x \notin FV(\psi)$$
: $(\forall x \ (\phi \land \psi)) \equiv ((\forall x \ \phi) \land \psi)$

$$(\forall x \ (\phi \lor \psi)) \equiv ((\forall x \ \phi) \lor \psi)$$



Rules of Inference

• Let x, y be variables and ϕ, ψ be propositional formulas. The following inference rules allow to deduce new formulas.

$$\frac{((\forall x \ \phi) \lor (\forall x \ \psi))}{(\forall x \ (\phi \lor \psi))} \qquad \frac{(\exists x \ (\phi \land \psi))}{(\exists x \ \phi) \land (\exists x \ \psi)} \qquad \frac{(\exists x \ (\forall y \ \phi))}{(\forall y \ (\exists x \ \phi))}$$

- Note that the other direction does not hold for any of these inference rules!
- In addition to these three inference rules all the equivalence rules apply: Each
 equivalence can be written as two inference rules since they are valid in both
 directions.



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Special Quantifiers

What is the syntactical meaning of

$$\sum_{i=m}^{n} f(i) ?$$

Apparently, this is the common short-hand notation for

$$\sum_{i=m}^{n} f(i) = \sum_{m \leq i \leq n} f(i) = \sum_{P(i,m,n)} f(i) = f(m) + f(m+1) + \cdots + f(n-1) + f(n),$$

where f(i) is a term with the free variable i and $(m \le i \le n)$ is a formula with free variables i, m, n, and $P(i, m, n) :\Leftrightarrow [(i \ge m) \land (i \le n)].$



Special Quantifiers

 Thus, the ∑-quantifier takes a predicate, P(i, m, n), and and a term, f(i), and converts it to the new term

$$(f(m) + f(m+1) + f(m+2) + \cdots + f(n-1) + f(n)),$$

By convention, the variable i is bound inside of $\sum_{i=m}^{n} f(i)$, while m and n remain free.

Similarly,

$$\prod_{i=m}^{n} f(i) := f(m) \cdot f(m+1) \cdot f(m+2) \cdot \ldots \cdot f(n-1) \cdot f(n).$$

Again, by convention, if n < m then

$$\sum_{i=m}^{n} f(i) := 0$$
 and $\prod_{i=m}^{n} f(i) := 1$.

• Union (\cup) and intersection (\cap) of several sets are further examples of special quantifiers: $\bigcup_{i=1}^{n} A_i$.



Special Quantifiers: Sets

- Standard notation for a set with a finite number of elements: { , ,..., };
 e.g., {1,2,3,4}.
- Obvious disadvantage: explicit enumeration of all elements of a set allows to specify only finite sets!
- Infinite sets require us to give a statement A to specify a characteristic property of the set:

$$S := \{x : A\}$$
 or $S := \{f(x) : A\},$

where S shall contain those elements x, or those terms f(x), for some universe of discourse, for which the statement A holds.

- Typically, *x* will be a free variable of *A*.
- Thus, the three symbols "{" and ":" and "}" together act as a quantifier that binds
 x.



Convenient Short-Hand Notations

 The following short-hand notations are convenient for using the predicate x ∈ X in conjunction with sets or quantifiers:

```
\{x \in X: A(x)\} is a short-hand notation for \{x: x \in X \land A(x)\}
(\forall x \in X \ A(x)) is a short-hand notation for (\forall x \ (x \in X \Rightarrow A(x)))
```

$$(\exists x \in X \ A(x))$$
 is a short-hand notation for $(\exists x \ (x \in X \land A(x)))$

If x is a typed variable — e.g., a real number — and P is a "simple" unary predicate — e.g., P(x): ⇔ (x > 3) — then the following notations are also used commonly:

$$(\forall P(x) \ A(x))$$
 is a short-hand notation for $(\forall x \ (P(x) \Rightarrow A(x)))$

$$(\exists P(x) \ A(x))$$
 is a short-hand notation for $(\exists x \ (P(x) \land A(x)))$

• Another wide-spread notation is to drop the parentheses:

$$\forall x \ P(x)$$
 instead of $(\forall x \ P(x))$



- Definitions and Theorem Proving
 - Need for Rigorous Analysis
 - Definitions
 - Syntactical Proof Techniques
 - Types of Proofs



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- How many compatible numbers can you pick for n := 20?



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- How many compatible numbers can you pick for n := 20?
- Intuition: Start at 1 and scan the integers from 1 to 20, successively picking those integers which are compatible with all integers picked previously:

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 is indeed the maximum number of compatible integers within {1, 2, 3, ..., 19, 20}.
 Right?



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 - We get 9 compatible integers. Our selection scheme makes it plausible that this is indeed the maximum number of compatible integers within $\{1, 2, 3, \ldots, 19, 20\}$. Right?
 - Well, what about the following 10 integers?
 - 1 3 5 7 9 11 13 15 17 19
 - Oops! Why should we believe that we can't find 11 or more compatible integers within $\{1, 2, 3, \ldots, 19, 20\}$?
 - The answer is provided by the pigeonhole principle (Thm. 147): Every subset of compatible integers of {1, 2, 3, ..., 19, 20} can contain at most one of each of the following 10 pairs:



Lesson Learned

• An intuitively appealing argument or approach is no substitute for a formal proof: Intuition might be wrong!



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Even though proofs and a rigorous formal analysis might seem boring (difficult, mind-boggling, mind-numbing, unnecessary, . . .) there is just no way around them if we want to be sure that our findings are correct!



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Even though proofs and a rigorous formal analysis might seem boring (difficult, mind-boggling, mind-numbing, unnecessary, \dots) there is just no way around them if we want to be sure that our findings are correct!

So, be prepared for at least some boring (difficult, mind-boggling, mind-numbing, unnecessary, ...) proofs!



Definitions and Theorem Proving

- Need for Rigorous Analysis
- Definitions
 - Basics of Definitions
 - Recursive Definitions
 - Fibonacci, Factorial, Sum, Product
 - Words
 - Caveats
- Syntactical Proof Techniques
- Types of Proofs



How to Deal with Formal Statements ...

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- Hence, prior to diving into other areas of Discrete Mathematics, we start with taking a practical look at the formal nuts and bolts of mathematical reasoning.
- In the following slides on definitions and theorem proving we pre-suppose an "intuitive" understanding of natural numbers, integers, reals, etc.; e.g., as taught in school.
- We will later on put these number systems on slightly more formal grounds.



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- An explicit definition relates an entity that is to be specified ("definiendum") to an already known entity ("definiens").



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$$f(x_1, x_2, \ldots, x_n) := t,$$

where the term t (normally) contains x_1, x_2, \dots, x_n as free variables.



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Warning

The definiendum does not occur in the definiens of an explicit definition of a function f or predicate P! That is, the symbols f and P do not appear on the right-hand side.

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- It is common to use the special symbols := and :⇔ for definitions, where the symbol ":" appears on the side of the definiendum.
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- Using =: and ⇔: is very good practice since
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 - it forces the author to decide whether or not something is a consequence of prior knowledge or some newly introduced entity.
- However, if ":=" or ":⇒" are used once in a text then they have to be used for absolutely all definitions in that text!!



 Poster seen in a tutoring institute at Salzburg:

$$x^{2}+\rho x+q=0$$

$$x_{1/2}=\frac{-\rho}{2}\pm\sqrt{\frac{\rho^{2}}{4}-q}$$

$$D=\frac{\rho^{2}}{4}-q$$

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- Better formalism:
 - If x_1, x_2 are the roots of the second-degree polynomial equation $x^2 + px + q = 0$, with $p, q \in \mathbb{R}$ and unknown $x \in \mathbb{R}$, then

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• With $D := p^2 - 4q$ we get

$$D\left\{\begin{array}{l} > \\ = \\ < \end{array}\right\}0: \left\{\begin{array}{l} \text{2 distinct real roots,} \\ \text{1 real root,} \\ \text{0 real roots.} \end{array}\right.$$



- Aka: Inductive definition.
- How can we state

x is ancestor of y if x is parent of y, or if x is parent of parent of y, or if x is parent of parent of y, or if x.

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- Recursive definitions (typically) consist of two parts:
 - a basis in which the definiendum does not occur in the definiens, and
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Warning

To avoid infinite circles, the definiendum must not occur in the basis!



Recursive Definitions: Sum and Product

Definition 32 (Sum and product)

Consider k real numbers $a_1, a_2, \ldots, a_k \in \mathbb{R}$, together with some $m, n \in \mathbb{N}$ such that $1 \leq m, n \leq k$.



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$$\sum_{i=m}^n a_i \ := \ \left\{ \begin{array}{ccc} 0 & \text{if} & n < m, \\ a_m & \text{if} & n = m, \\ (\sum_{i=m}^{n-1} a_i) + a_n & \text{if} & n > m, \end{array} \right.$$



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 The definitions for n < m are convenience settings that have turned out to be useful in practice.



Definition 33 (Factorial, Dt.: Fakultät, Faktorielle)

For
$$n \in \mathbb{N}_0$$
,

$$n! := \begin{cases} 1 & \text{if } n \leq 1, \\ n \cdot (n-1)! & \text{if } n > 1. \end{cases}$$



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Definition 34 (Fibonacci numbers)

For $n \in \mathbb{N}_0$,

$$F_n := \begin{cases} 0 & \text{if } n = 0, \\ 1 & \text{if } n = 1, \\ F_{n-1} + F_{n-2} & \text{if } n \ge 2. \end{cases}$$



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Fibonacci Numbers

- The Fibonacci numbers are named after Leonardo da Pisa (1180?–1241?), aka "figlio di Bonaccio".
- The Fibonacci numbers have been studied extensively; they exhibit lots of interesting mathematical properties. For instance,

$$\lim_{n\to\infty}\frac{F_{n+1}}{F_n}=\phi,\quad \text{where }\phi:=\frac{1+\sqrt{5}}{2}\text{ is known as } \text{golden ratio.}$$

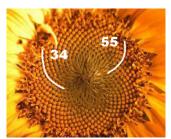


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 The Fibonacci numbers are also found in nature: E.g., the numbers of CW/CCW spirals of sunflower heads are given by subsequent Fibonacci numbers.



[Image credit: Wikipedia.]



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Definition 35 (Word)

Let Σ be a finite set. The set Σ^* of *words* over Σ is defined follows:

① Base clause: The empty word, denoted by the Greek letter ϵ , belongs to Σ^* .



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Definition 35 (Word)

- **①** Base clause: The empty word, denoted by the Greek letter ϵ , belongs to Σ^* .
- **2** Recursion clause: For all $a \in \Sigma$ and all $\sigma \in \Sigma^*$, the ordered pair (a, σ) belongs to Σ^* .



- Consider an arbitrary (but fixed) finite set Σ. We call it the alphabet; the individual elements of Σ are called symbols or characters.
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- **3** Extremal clause: A word is in Σ^* if it is ϵ or if it can be constructed from ϵ via a finite number of applications of the recursion clause.



- Consider an arbitrary (but fixed) finite set Σ. We call it the alphabet; the individual elements of Σ are called *symbols* or *characters*.
- E.g., $\Sigma := \{a, b, c, \dots, x, y, z\}$ or $\Sigma := \{0, 1\}$.

Definition 35 (Word)

- Base clause: The empty word, denoted by the Greek letter ε, belongs to Σ*.
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- **3** Extremal clause: A word is in Σ^* if it is ϵ or if it can be constructed from ϵ via a finite number of applications of the recursion clause.
 - Aka string (Dt. Zeichenkette). The set Σ^* of all words over Σ is known as Kleene closure of Σ



- Consider an arbitrary (but fixed) finite set Σ. We call it the alphabet; the individual elements of Σ are called symbols or characters.
- E.g., $\Sigma := \{a, b, c, \dots, x, y, z\}$ or $\Sigma := \{0, 1\}$.

Definition 35 (Word)

- **①** Base clause: The empty word, denoted by the Greek letter ϵ , belongs to Σ^* .
- **2** Recursion clause: For all $a \in \Sigma$ and all $\sigma \in \Sigma^*$, the ordered pair (a, σ) belongs to Σ^* .
- **3** Extremal clause: A word is in Σ^* if it is ϵ or if it can be constructed from ϵ via a finite number of applications of the recursion clause.
- Aka string (Dt. Zeichenkette). The set Σ^* of all words over Σ is known as Kleene closure of Σ .
- Of course, in order to avoid confusion, ϵ is not allowed to be a character of Σ .



- Consider an arbitrary (but fixed) finite set Σ . We call it the *alphabet*; the individual elements of Σ are called *symbols* or *characters*.
- E.g., $\Sigma := \{a, b, c, \dots, x, y, z\}$ or $\Sigma := \{0, 1\}$.

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- Of course, in order to avoid confusion, ϵ is not allowed to be a character of Σ .
- It is important to note that every element of Σ^* is a finite sequence of zero or more characters (if we disregard the parentheses and commas) but that Σ^* itself is an infinite set containing words of every possible finite length.

Words: Length and Concatenation

Definition 36 (Length of a word)

Let Σ be a finite set. The *length* of a word σ over Σ is defined as follows:

$$|\sigma| := \begin{cases} 0 & \text{if } \sigma = \epsilon, \\ 1 + |\sigma'| & \text{if } \sigma = (a, \sigma') \text{ for some } a \in \Sigma \text{ and } \sigma' \in \Sigma^*. \end{cases}$$



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$$\sigma_1 \bullet \sigma_2 := \begin{cases} \sigma_2 & \text{if } \sigma_1 = \epsilon, \\ (a, \sigma_1' \bullet \sigma_2) & \text{if } \sigma_1 = (a, \sigma_1') \text{ for some } a \in \Sigma \text{ and } \sigma_1' \in \Sigma^*. \end{cases}$$



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- In practice it is a convention to drop the ordered-pair notation and to write $a\sigma$ rather than (a, σ) . E.g., word rather than $(w, (o, (r, (d, \epsilon))))$.
- Similarly, one writes word rather than wo rd. (This simplification is justified by the fact that the binary operator • is associative.)

Definitions like

$$P(x, y, z) :\Leftrightarrow (x < 2y)$$
 or $P(x) :\Leftrightarrow (x < 2y)$

can be seen as syntactically correct but they are semantically problematic!



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$$\frac{m}{n}\sharp\frac{p}{q}:=\frac{m+p}{n+q}.$$

- Then $\frac{1}{1}\sharp \frac{2}{3} = \frac{3}{4}$, but $\frac{2}{2}\sharp \frac{2}{3} = \frac{4}{5}$.
- Since $\frac{1}{1} = \frac{2}{2}$, we conclude $\frac{4}{5} = \frac{3}{4}$, and, thus, 0 = 1. Yikes!



- Definitions and Theorem Proving
 - Need for Rigorous Analysis
 - Definitions
 - Syntactical Proof Techniques
 - Syntax and Proofs
 - Equivalence Transformations
 - Types of Proofs



Definition 38 (Proof, Dt.: Beweis)

To *prove* a statement means to derive it from axioms (or postulates) and other previously established theorems by means of rules of logic.



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A statement is a *theorem* if it has been proved. If the statement is of the form $H \Rightarrow C$ then we call H the *hypothesis* and C the *conclusion*.



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- Depending on the importance of the result, terms like *lemma* (Dt.: Lemma, Hilfssatz) or *corollary* (Dt.: Korollar) are also used instead of "theorem".



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- Depending on the importance of the result, terms like lemma (Dt.: Lemma, Hilfssatz) or corollary (Dt.: Korollar) are also used instead of "theorem".
- A conjecture is a statement which has not yet been proved or disproved.
- The status of a conjecture may remain unknown for decades or even centuries:
 Fermat's Last Theorem was stated by Pierre de Fermat in 1637 and proved by Andrew Wiles (with the help of Richard Taylor) in 1993–1995.

- Syntactical proof techniques are proof techniques that are based on the analysis
 of the syntactical structure of a statement.
- Syntactical proof techniques allow us to reason about statements and to simplify statements with no or very little "understanding" of their mathematical meaning.
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$H \Rightarrow C \dots$

 \dots is true if either H is false (and C arbitrary) or if C is true for H being true.



- If conclusion *C* is of the form $(A \wedge B)$:
 - Prove A under the assumption H; and
 - Prove B under the assumption H.



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Warning

In all the rules on this slide, *A* and *B* must not be part of a quantified formula. (Otherwise, get rid of the quantifier first!)

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- If conclusion *C* is of the form $(\forall x \ A)$:
 - Proof technique: Let x_0 be arbitrary but fixed (Dt.: "beliebig aber fix"). From now on, x_0 can be treated as a constant!
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- We are not allowed to make any assumptions on x₀ except for those that hold for all x in the universe of discourse.



- If conclusion *C* is of the form $(\exists x \ A)$:
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Proof (existential): We have p(2) = 5 > 0 and p(0) = -1 < 0. Since p is continuous on the closed interval [0,2], the Intermediate Value Theorem (Dt.: Zwischenwertsatz) tells us that there exists a real number x strictly between 0 and 2 such that p(x) = 0.



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- If conclusion C is of the form $(\exists!x \ A)$:
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 - Prove its uniqueness.
- If hypothesis H is of the form $(\exists x \ A)$:
 - Let x_0 such that $A[x_0/x]$.
 - Add $A[x_0/x]$ to knowledge.
 - Again: x₀ must not occur anywhere else in H or C!



Natural-Language Synonyms of Formal Terms

- On many occasions a conjecture will not be stated in formal terms but by using a natural language.
- Then one has to decode the natural-language formulation and translate it into formal terms!



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 - A implies B, A impliziert B,
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 - B if A,
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 - A if and only if B, A genau dann wenn B,
 - A is necessary and sufficient for B,
 A ist notwendig und hinreichend für B.



Equivalence Transformations

• First attempt to prove $(\forall n \in \mathbb{N} \quad \frac{2n+1}{n+1} \geqslant \frac{3}{2})$:

$$\frac{2n+1}{n+1} \geqslant \frac{3}{2}$$

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• Correct proof of $(\forall n \in \mathbb{N} \mid \frac{2n+1}{n+1} \geqslant \frac{3}{2})$: Let $n \in \mathbb{N}$ be arbitrary but fixed. Then:

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a &= b \\
\Rightarrow & a^2 &= ab
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\Leftrightarrow & (a+b) & = b & | a := b \\
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• And here comes a "proof" of 4 = 5:



• Let $a, b \in \mathbb{N}$ be equal natural numbers. We "prove" that 1 = 2:

$$x + y = 9$$



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\end{array}$$



• Let $a, b \in \mathbb{N}$ be equal natural numbers. We "prove" that 1 = 2:

$$\begin{array}{ccccc}
 & a = b & | \cdot a \\
 & a^2 = ab & | -b^2 \\
 & \Leftrightarrow & a^2 - b^2 = ab - b^2 \\
 & \Leftrightarrow & (a-b) \cdot (a+b) = b \cdot (a-b) & | \div (a-b) \\
 & \Leftrightarrow & (a+b) = b & | a := b \\
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$$x + y = 9$$

$$x^{2} - y^{2} = 9x - 9y$$

$$x^{2} - 9x + \frac{81}{4} = y^{2} - 9y + \frac{81}{4}$$

$$(x - \frac{9}{2})^{2} = (y - \frac{9}{2})^{2}$$

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$$x = y$$

$$4 = 5$$

$$| \cdot (x - y)$$

$$| + \frac{81}{4} - 9x + y^{2}$$

$$| \sqrt{}$$

$$| + \frac{9}{2}$$



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Advice

• In general, a relation $a \circ b$ may only be replaced by a new relation $a' \circ b'$ if one can argue that $(a \circ b) \Leftrightarrow (a' \circ b')$.



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Advice

- In general, a relation $a \circ b$ may only be replaced by a new relation $a' \circ b'$ if one can argue that $(a \circ b) \Leftrightarrow (a' \circ b')$.
- It is advisable to prove $a \circ b$, where $o \in \{=, <, >, \leqslant, \ge\}$, by constructing a chain $a_0 \circ a_1 \circ a_2 \circ \ldots \circ a_n$, with $a_0 = a$ and $a_n = b$, for some $n \in \mathbb{N}$.



3 Definitions and Theorem Proving

- Need for Rigorous Analysis
- Definitions
- Syntactical Proof Techniques
- Types of Proofs
 - Without Loss of Generality
 - Direct Enumeration
 - Case Analysis
 - Direct Proof
 - Proof by Contrapositive
 - Proof by Contradiction
 - Indirect Proof
 - Disproving Conjectures



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- Dt.: O.B.d.A. ("Ohne Beschränkung der Allgemeinheit").



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Warning

Do not use "w.l.o.g." unless *you could* indeed *explain* explicitly and in full detail how to carry on without that assumption!



Types of Proofs: Direct Enumeration

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 - E.g.: The conjecture
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 - Note: Direct enumeration only works if the set given is finite!



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- Aka Proof by Exhaustion. Dt.: Fallunterscheidung.
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 $(H \wedge A_k) \Rightarrow C.$

Warning

It is essential to guarantee that $A_1 \vee A_2 \vee \ldots \vee A_k$ holds, i.e., that no case is missing!



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Proof: Factoring $n^7 - n$ yields

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Let n:=7q+r with $q,r\in\mathbb{N}_0$ and $0\leqslant r\leqslant 6$. We consider seven cases, depending on whether r=0,1,2,3,4,5 or 6.



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Case n = 7q + 6: Then n + 1 = 7q + 7 is divisible by 7.



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- We want to prove $H \Rightarrow C$:
 - We build a chain of reasoning that starts at H and ends in C.
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$$x_0^2 = x_0 \cdot x_0 < y_0 \cdot x_0 < y_0 \cdot y_0 = y_0^2$$



Types of Proofs: Proof by Contrapositive

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which implies $y_0 \le x_0$ since we may divide by the positive number $x_0 + y_0$.



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Warning

Make sure that the statements are negated correctly!



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 - We assume $(H \land \neg C)$ as new hypothesis and prove $\neg H$.
 - This approach is correct since $(H \Rightarrow C) \equiv ((H \land \neg C) \Rightarrow \neg H)$.
- Warning: As when proving the contrapositive it is essential to check twice that the statements are indeed negated correctly!



- Aka Reductio ad absurdum.
- Dt.: indirekter Beweis.
- We want to prove $H \Rightarrow C$.
 - Consider a statement R that is known to be true, like $0 \neq 1$.



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Note

Since an indirect proof is similar to a proof by contradiction, many textbooks treat it as one proof technique, or use the terms "reductio ad absurdum", "indirect proof", and "proof by contradiction" as synonyms.



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$$0 = \frac{p^3}{q^3} + \frac{p}{q} + 1$$
 and, thus, $0 = p^3 + pq^2 + q^3$.

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Case p,q even: This is not possible since we assumed (rightfully) that $\frac{p}{q}$ is irreducible.



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- If the conjecture is of the form $(\forall x \ A)$:
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- If, however, the conjecture is of the form $(\exists x \ A)$:
 - Then a counterexample does not suffice!
 - Rather, to disprove this conjecture, we'd have to prove formally $(\forall x \neg A)$.



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 - Real Numbers
 - More Proof Techniques



- 4 Nu
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 - Field
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- The axioms tell us the properties of the operations.
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- E.g., we have $(\sqrt{\pi} + 1) \sqrt{\pi} = 1$ because

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- Algebraic structures get their names based on the type of operations and axioms supported.
- Well-known structures include group, ring, field, and vector space. (Many more algebraic structures are studied in abstract algebra, though!)



Definition 40 (n-ary Operation, Dt.: n-stellig Verknüpfung)

Let n be a fixed non-negative integer and X_1, X_2, \ldots, X_n be non-empty sets. An n-ary operation from X_1, X_2, \ldots, X_n to another set Y is a function $\omega : X_1 \times X_2 \times \cdots \times X_n \to Y$.



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 An operation on a set X is also called an internal operation (Dt.: innere Verknüpfung).



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- Note: The standard division ÷ is a binary operation neither on the natural numbers nor on the rational numbers.



So, a binary operation on a set X is a function

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- For binary operations it is customary to use symbols like $\star, \circ, +, \cdot, \div$ rather than letters like ω .
- Furthermore, for binary operations it is common to use the infix notation

$$X_1 \star X_2$$
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• The symbol – tends to be used both for an unary and a binary operation.

Definition 41 (Composition, Dt.: Hintereinanderausführung)

Consider two operations $f:A\to B$ and $g:B\to C$. The composition (Dt.: Komposition, Hintereinanderausführung) $g\circ f$ of f and g is defined as

$$(g \circ f)(x) := g(f(x))$$
 for all $x \in A$.



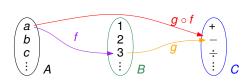
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 for all $x \in A$.

• That is, the standard interpretation of $g \circ f$ is "carry out f followed by g".





Definition 41 (Composition, Dt.: Hintereinanderausführung)

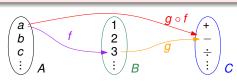
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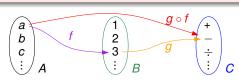
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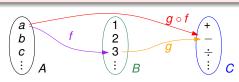
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- We will use the symbol o exclusively for denoting compositions of operations.

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Not all authors stick to the convention $(g \circ f)(x) := g(f(x)) \dots$





Definition 42 (Associativity, Dt.: Assoziativität)

A binary operation \star on a (non-empty) set G is associative if

$$\forall a, b, c \in G \quad (a \star b) \star c = a \star (b \star c).$$



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- Associativity means that the order in which consecutive operations are applied does not change the result.
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Properties of Operations: Distributivity

Definition 44 (Distributivity, Dt.: Distributivität)

A binary operation · on a (non-empty) set G is distributive over a binary operation + on G if

$$\forall a, b, c \in G \quad a \cdot (b + c) = (a \cdot b) + (a \cdot c),$$

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- Note that addition does not distribute over multiplication (over \mathbb{R}).



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- Note that addition does not distribute over multiplication (over \mathbb{R}).
- Some textbooks prefer to split up the conditions of Def. 44 and say that · is left-distributive if

$$\forall a, b, c \in G \quad a \cdot (b + c) = (a \cdot b) + (a \cdot c),$$

and right-distributive if

$$\forall a, b, c \in G \quad (a+b) \cdot c = (a \cdot c) + (b \cdot c).$$



Definition 45 (Neutral element, Dt.: neutrales Element)

The element $n \in G$ is a *neutral element* (aka zero element, identity element) of a binary operation \star on a (non-empty) set G if

$$\forall a \in G \quad a \star n = a = n \star a.$$



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Definition 46 (Inverse element, Dt.: inverses Element)

The element $b \in G$ is an *inverse element* of the element $a \in G$ for the binary operation \star on a (non-empty) set G if

$$a \star b = n = b \star a$$
.

where n denotes the neutral element of \star on G.



Properties of Operations: Uniqueness of Neutral Element

Lemma 47

A binary operation \star on a (non-empty) set G has at most one neutral element.



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Proof: Assume that $n_1, n_2 \in G$ are neutral elements of \star on G. By Def. 45,

$$\forall a \in G \quad a \star n_1 = a = n_1 \star a \quad \text{and} \quad \forall a \in G \quad a \star n_2 = a = n_2 \star a.$$



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These identities hold for all $a \in G$. Hence, in particular, they have to hold if $a := n_1$ and $a := n_2$:

$$n_2 \star n_1 = n_2 = n_1 \star n_2$$
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Corollary 48

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If a binary operation \star on a (non-empty) set G has a neutral element then it is unique.

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An element $a \in G$ has at most one inverse element $b \in G$ for an associative binary operation \star on G.



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Definition 51 (Group, Dt.: Gruppe)

A set G together with a binary operation \star on G defines a *group* if the following properties hold:



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- Note that a ⋆ b = b ⋆ a is not required for all a, b ∈ G. That is, commutativity need not hold!



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 - Sample (Abelian) groups: the integers Z under addition, non-zero rational numbers $\mathbb{Q}\setminus\{0\}$ under multiplication.
 - Not a group: The integers under multiplication.



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- Multiplication tables for groups of orders two and three:

*	n	а
n	n	а
а	а	n

*	n	а	b
n	n	а	b
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b	b	n	а



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- Again up to renaming, there are only two possible multiplication tables for groups with four elements, i.e., only two different groups.



• The dihedral group (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:





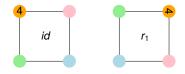
Finite Group: Dihedral Group D₄

- The dihedral group (Dt.: Diedergruppe) D₄ is formed by the clockwise rotations and reflections of a square which map the square onto itself:
 - id,



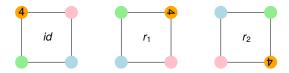


- The dihedral group (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:
 - id, r_1 (CW rotation by 90°),



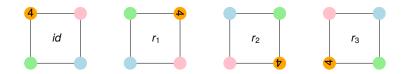


- The dihedral group (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:
 - id, r_1 (CW rotation by 90°), r_2 (CW rotation by 180°),



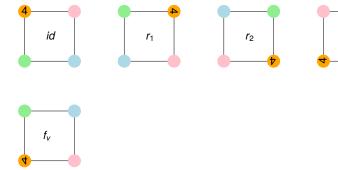


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 - id, r_1 (CW rotation by 90°), r_2 (CW rotation by 180°), r_3 (CW rotation by 270°);





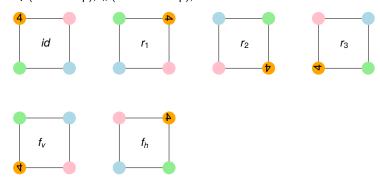
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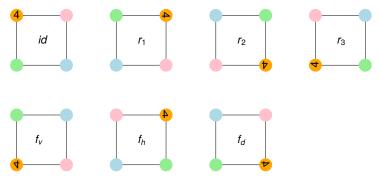
 r_3

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 - f_V (vertical flip), f_h (horizontal flip),



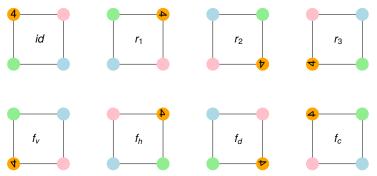


- The dihedral group (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:
 - id, r_1 (CW rotation by 90°), r_2 (CW rotation by 180°), r_3 (CW rotation by 270°);
 - f_v (vertical flip), f_h (horizontal flip), f_d (diagonal flip),



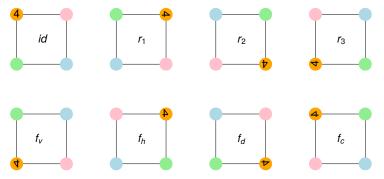


- The dihedral group (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:
 - id, r_1 (CW rotation by 90°), r_2 (CW rotation by 180°), r_3 (CW rotation by 270°);
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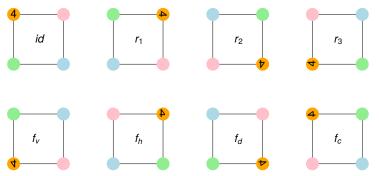
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• Does D₄ have eight elements? Or did we miss any element?



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- Does D_4 have eight elements? Or did we miss any element?
- No, we didn't!



- \bullet We denote the composition of functions by $\circ.$
- Multiplication table of D₄:

0	id	<i>r</i> ₁	<i>r</i> ₂	<i>r</i> ₃	f_{V}	f _h	f _d	f _c
id	id	<i>r</i> ₁	<i>r</i> ₂	<i>r</i> ₃	f_{V}	f _h	f _d	f _c
<i>r</i> ₁	<i>r</i> ₁	<i>r</i> ₂	<i>r</i> ₃	id	f _c	f _d	f_{V}	f _h
<i>r</i> ₂	<i>r</i> ₂	<i>r</i> ₃	id	<i>r</i> ₁	f_h	f_{V}	f _c	f _d
<i>r</i> ₃	<i>r</i> ₃	id	<i>r</i> ₁	<i>r</i> ₂	f _d	f _c	f _h	f_{V}
f_{v}	f_{V}	f_d	f _h	f_c	id	r_2	<i>r</i> ₁	<i>r</i> ₃
f _h	f_h	f_c	f_{V}	f_d	r_2	id	<i>r</i> ₃	<i>r</i> ₁
f_d	f_d	f_h	f _c	f_{ν}	r_3	<i>r</i> ₁	id	<i>r</i> ₂
f _c	f _c	f_{V}	f_d	f_h	<i>r</i> ₁	r_3	r_2	id

• E.g., $f_d \circ f_v$, which means flip vertically and then flip diagonally, corresponds to a (clockwise) rotation by 270^o , i.e., to r_3 .



- We denote the composition of functions by o.
- Multiplication table of D₄:

0	id	<i>r</i> ₁	<i>r</i> ₂	<i>r</i> ₃	f_{V}	f _h	f _d	f _c
id	id	<i>r</i> ₁	<i>r</i> ₂	<i>r</i> ₃	f_{V}	f _h	f _d	f_c
<i>r</i> ₁	<i>r</i> ₁	r ₂	<i>r</i> ₃	id	f _c	f _d	f_{V}	f _h
<i>r</i> ₂	<i>r</i> ₂	<i>r</i> ₃	id	<i>r</i> ₁	f_h	f_{V}	f _c	f _d
<i>r</i> ₃	<i>r</i> ₃	id	<i>r</i> ₁	<i>r</i> ₂	f _d	f _c	f _h	f_{V}
f_{v}	f_{V}	f_d	f _h	f_c	id	r_2	<i>r</i> ₁	<i>r</i> ₃
f _h	f_h	f _c	f_{V}	f_d	r_2	id	<i>r</i> ₃	<i>r</i> ₁
f_d	f _d	f_h	f _c	f_{ν}	r_3	<i>r</i> ₁	id	<i>r</i> ₂
f_c	f _c	f_{v}	f_d	f_h	<i>r</i> ₁	r_3	r_2	id

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<i>r</i> ₁	<i>r</i> ₁	r ₂	<i>r</i> ₃	id	f _c	f _d	f_{V}	f _h
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<i>r</i> ₃	<i>r</i> ₃	id	<i>r</i> ₁	r_2	f _d	f _c	f _h	f_{V}
f_{v}	f_{V}	f_d	f _h	f_c	id	r_2	<i>r</i> ₁	<i>r</i> ₃
$ f_h $	f_h	f _c	f_{V}	f_d	<i>r</i> ₂	id	<i>r</i> ₃	<i>r</i> ₁
$ f_d $	f_d	f_h	f _c	f_{ν}	<i>r</i> ₃	<i>r</i> ₁	id	r_2
f _c	f_c	f_{v}	f_d	f_h	<i>r</i> ₁	<i>r</i> ₃	r_2	id

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- Note: $f_d \circ f_V \neq f_V \circ f_d$. That is, D_4 is not commutative.
- Note that each one of the transformations appears exactly once in each row and each column of the table: Latin square.

Real-World Application: Geometric Crystal Classes

- D₄ is one of the so-called crystallographic point groups, which describe sets of symmetry operations relative to a fixed point. Aka geometric crystal class.
- Each operation leaves the structure of the crystal unchanged. That is, the same types of atoms appear in similar positions as before the transformation induced by the operation.



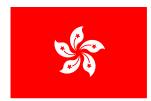
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The Bauhinia flower has C_5 symmetry, and each star has D_5 symmetry.



This (color-inverted) snowflake has D_6 symmetry.

Definition 53 (Ring, Dt.: Ring mit Eins)

A set R which possesses an "addition" $+: R \times R \to R$ and a "multiplication" $\cdot: R \times R \to R$ defines a *(unit) ring* if the following conditions hold:



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- Sample ring: The set of all continuous real-valued functions defined over an interval $[\alpha,\beta]\subset\mathbb{R}$, with addition and multiplication of functions as operations forms a ring.

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A set F which possesses an "addition" $+: F \times F \to F$ and a "multiplication" $\cdot: F \times F \to F$ defines a *field* if the following conditions hold:



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- **4** Commutativity: $\forall a, b \in F$ $a \cdot b = b \cdot a$.
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- **3** For all $a \in F$ there exists an additive inverse $b \in F$, satisfying a + b = 0.
- For all $a \in F \setminus \{0\}$ there exists a multiplicative inverse $b \in F$, satisfying $a \cdot b = 1$.



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- **9** For all $a \in F \setminus \{0\}$ there exists a multiplicative inverse $b \in F$, satisfying $a \cdot b = 1$.
- $0 \neq 1$.



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 - Again, the elements of F need not be numbers.



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- For all $a \in F \setminus \{0\}$ there exists a multiplicative inverse $b \in F$, satisfying $a \cdot b = 1$.
- **1** $0 \neq 1$
 - Again, the elements of F need not be numbers.
 - Note: The multiplication sign is often dropped if the meaning is clear within specific context: It is common to write ab rather than $a \cdot b$.

 In the sequel, we denote the additive neutral element of a field (F, +, ·) by 0 and its multiplicative neutral element by 1. Furthermore, we denote the inverse elements of b ∈ F by −b and b⁻¹.



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Let $(F, +, \cdot)$ be a field. We define the binary operation "subtraction" $-: F \times F \to F$:

$$\forall a, b \in F \quad a - b := a + (-b)$$



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Let $(F, +, \cdot)$ be a field. We define the binary operation "division" $\div : F \times (F \setminus \{0\}) \to F$:

$$\forall a \in F, b \in F \setminus \{0\} \quad a \div b := a \cdot b^{-1}$$



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Lemma 57

Let $(F, +, \cdot)$ be a field.

$$\forall a \in F \quad a - a = 0$$
 and $\forall a \in F \setminus \{0\} \quad a \div a = 1$.

Theorem 58

Let $(F, +, \cdot)$ be a field. Then

$$-0 = 0$$
 and $1^{-1} = 1$ and $\forall a \in F \ 0 \cdot a = 0$.



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Proof: We have 0 = 0 + (-0) = -0.



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$$0 = 0 + (-0) = -0$$
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Proof: We have 0 = 0 + (-0) = -0. Similarly, $1 = 1 \cdot 1^{-1} = 1^{-1}$. Let $a \in F$ be arbitrary but fixed. Then

$$0=0\cdot a+\big(-(0\cdot a)\big)$$



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$$0 = 0 \cdot a + (-(0 \cdot a)) = (0 + 0) \cdot a - 0 \cdot a$$



Theorem 58

Let $(F, +, \cdot)$ be a field. Then

$$-0 = 0$$
 and $1^{-1} = 1$ and $\forall a \in F \quad 0 \cdot a = 0$.

Proof: We have 0 = 0 + (-0) = -0. Similarly, $1 = 1 \cdot 1^{-1} = 1^{-1}$. Let $a \in F$ be arbitrary but fixed. Then

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Theorem 59

Let $(F, +, \cdot)$ be a field. Then, for all $a, b \in F$,

$$(-1) \cdot a = -a$$
 and $-(-a) = a$ and

$$(-a) \cdot b = -(a \cdot b) = a \cdot (-b)$$
 and $(-a) \cdot (-b) = a \cdot b$.



Theorem 60

Let $(F, +, \cdot)$ be a field. Then, for all $a, b \in F$,

$$a \cdot b = 0 \implies (a = 0 \text{ or } b = 0).$$



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$$0=a^{-1}\cdot 0$$



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Let $(F, +, \cdot)$ be a field. Then, for all $a, b \in F \setminus \{0\}$,

$$(a^{-1})^{-1} = a$$
 and $(a \cdot b)^{-1} = a^{-1} \cdot b^{-1}$.



Definition 62 (Fraction, Dt.: Bruch)

For $a \in F, b \in F \setminus \{0\}$, the fraction $\frac{a}{b}$ is defined as

$$\frac{a}{b} := a \div b.$$

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Let $(F, +, \cdot)$ be a field. Then, for all $a, b, x, y \in F$ for which no denominator equals 0,

$$\frac{a}{b} \pm \frac{x}{y} = \frac{a \cdot y \pm b \cdot x}{b \cdot y}$$
 and $\frac{a}{b} \cdot \frac{x}{y} = \frac{a \cdot x}{b \cdot y}$ and $\frac{a}{b} \div \frac{x}{y} = \frac{a \cdot y}{b \cdot x}$.

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- So, for two structures A, B and two binary operations ⋆_A (on A) and ⋆_B (on B), if
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- E.g., $(\mathbb{R},+)$ and (\mathbb{R}^+,\cdot) form groups. A group homomorphism from $(\mathbb{R},+)$ to (\mathbb{R}^+,\cdot) is given by the exponential function $x\mapsto e^x$. (Recall that $e^{x+y}=e^x\cdot e^y$.)



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- A ring homomorphism from A to B is compatible with ring addition and multiplication, and maps the multiplicative neutral element of A to the multiplicative neutral element of B.



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- A ring homomorphism from A to B is compatible with ring addition and multiplication, and maps the multiplicative neutral element of A to the multiplicative neutral element of B.
- An isomorphism is a bijective homomorphism.



Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
 - Orders
 - Peano's Axioms for Introducing the Natural Numbers
 - The Principle of Mathematical Induction
 - Cardinality
- Integers
- Rational Numbers
- Real Numbers
- More Proof Techniques



How Shall We Define Natural Numbers or Real Numbers?

- Three options:
 - Ignore all formal details and presuppose an "intuitive" understanding of reals, integers, . . .
 - ② Introduce the natural numbers, \mathbb{N} , and then construct a hierarchy of number systems: $\mathbb{N} \subset \mathbb{Z} \subset \mathbb{Q} \subset \mathbb{R}$.
 - $\textbf{ § Set up the reals, } \mathbb{R}, \text{ axiomatically and then define proper subsets for } \mathbb{N}, \mathbb{Z}, \mathbb{Q}.$



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 - **3** Set up the reals, \mathbb{R} , axiomatically and then define proper subsets for $\mathbb{N}, \mathbb{Z}, \mathbb{Q}$.
- What is the best approach for a course on (applied) discrete mathematics? Much scholarly debate — no consensus!
- ullet We will start with introducing the natural numbers. However, since the gory details result in a lengthy discussion which provides little additional insight in \mathbb{N} and this is no course on number theory we base our introduction of \mathbb{N} on a simplified treatment of the so-called Peano axioms; see a book on number theory for a more formalized introduction of \mathbb{N} .



- Intuitively, the natural numbers $\mathbb N$ are given by $\{1,2,3,4,5,\ldots\}$ or by $\{0,1,2,3,4,5,\ldots\}.$
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Convention

In this course we adopt the following convention:

$$\mathbb{N}:=\{1,2,3,4,5,\ldots\} \ \ \text{and} \ \ \mathbb{N}_0:=\{0,1,2,3,4,5,\ldots\}.$$



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- Caution: Read a text carefully to learn what an author means by "natural number".
 In particular, watch for clues such as terms like "positive natural numbers" (which indicates that zero is included) or statements like "n is a natural number, so it must be greater than zero" (which indicates that zero is not included).
- If one treats 0 as an element of $\mathbb N$ then $\{1,2,3,4,5,\ldots\}$ is often denoted by $\mathbb N^*$.



Definition 65 (Partial order, Dt.: Halbordnung)

A *partial order* on a set S is a binary relation \leq , i.e., a subset of $S \times S$, such that the following three properties hold for all $a, b, c \in S$:



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- Irreflexivity: $\neg (a < a)$.
- 2 Transitivity: $(a < b \land b < c) \Rightarrow a < c$.
 - A strict partial order is always asymmetric: If a < b then $\neg (b < a)$.
 - $(a < b \land b < a) \stackrel{trans.}{\Rightarrow} a < a$, in contradiction to the irreflexivity: $\neg(a < a) \stackrel{\land}{\rightarrow}$

Theorem 67

There is a one-to-one correspondence between non-strict and strict partial orders. Let S be a set and $a, b \in S$.

- If \leq is a non-strict partial order on S then the corresponding strict partial order "<" on S is the *reflexive reduction* given by
 - a < b : \Leftrightarrow $a \le b$ and $a \ne b$.
- ② If, on the other hand, < is a strict partial order on S then the corresponding non-strict partial order " \le " on S is the *reflexive closure* given by

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 As a notational convention, we omit the indication of an equality sign if we refer to a strict order, e.g., we write < rather than ≤ or ⊂ rather than ⊆.



 \bullet E.g., $(\mathbb{Z}, \trianglerighteq)$ with (the non-strict order) \trianglerighteq as defined below forms a poset:

if a and b are even: $a \ge b$: $\Leftrightarrow a \ge b$

if a and b are odd: $a \ge b :\Leftrightarrow a \le b$

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Note that we do not know $a \trianglerighteq b$ if one of a, b is even and the other one is odd. That is, if (S, \leq) is a poset then not all pairs of elements of S need to be comparable!

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• The subset relation, \subset , on the powerset $\mathcal{P}(X)$ of a set X is a strict partial order.



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Definition 68 (Dual order, Dt.: duale Ordnung)

Let (S, \leq) resp. (S, <) be a (strict) poset. The *dual order* (or *reverse order*) on S, \geq resp. >, is defined as follows for $a, b \in S$:

$$a \ge b$$
 : \Leftrightarrow $b \le a$

$$a > b$$
 : \Leftrightarrow $b < a$.



Definition 69 (Minimal element, Dt.: minimales Element)

Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *minimal element* of T if no $b \in T \setminus \{a\}$ exists such that $b \leq a$.



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Definition 70 (Least element, Dt.: kleinstes Element, Minimum)

Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *least element* (or *minimum*) of T if $\forall b \in T \setminus \{a\}$ $a \leq b$.



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Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *minimal element* of T if no $b \in T \setminus \{a\}$ exists such that $b \leq a$.

Definition 70 (Least element, Dt.: kleinstes Element, Minimum)

Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *least element* (or *minimum*) of T if $\forall b \in T \setminus \{a\}$ $a \leq b$.

Definition 71 (Maximal element, Dt.: maximales Element)

Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *maximal element* of T if no $b \in T \setminus \{a\}$ exists such that $a \leq b$.

Definition 72 (Greatest element, Dt.: größtes Element, Maximum)

Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *greatest element* (or *maximum*) of T if $\forall b \in T \setminus \{a\}$ $b \leq a$.

• Note: If a minimum or maximum exists then the anti-symmetry ensures that it is unique. Minimal or maximal elements need not be unique, though.

Definition 73 (Total order, Dt.: totale Ordnung)

A binary relation \leq on a set S forms a *total order* (or *linear order*) on S if the following three statements hold for all $a, b, c \in S$:



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- Note that (1) in Def. 73 implies reflexivity: $a \le a$ for all $a \in S$.
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Definition 74 (Well-order, Dt.: Wohlordnung)

A total order \leq on a set S forms a *well-order* if every non-empty subset of S has a least element.



 The following definition of N is based on a simplified version of Peano's Axioms, as proposed by Giuseppe Peano (1858–1932) in 1889.



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The number n + 1 is called the *successor* of n, sometimes denoted by succ(n).



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- ullet One can show that the standard algebraic rules are compatible with the conditions imposed on $\mathbb N$, and that algebra and order interact smoothly within $\mathbb N$.
- One can also show that (up to a renaming of elements) there is only one set that fulfills all conditions of Def. 75. Hence, $\mathbb N$ is uniquely defined.

Definition 76 (Inductive)

A set $K \subseteq \mathbb{N}$ is *inductive* if

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Theorem 78 (Weak Principle of Induction (W.P.I.))

Consider a predicate P over \mathbb{N} .

lf

and if

$$\forall k \in \mathbb{N} \ (P(k) \Rightarrow P(k+1))$$

then

$$\forall n \in \mathbb{N} \ P(n)$$
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Thus, Thm. 77 is applicable and we conclude $K = \mathbb{N}$. That is, the predicate P holds for all natural numbers.

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- **1** Inductive step ("IS"): We prove P(k + 1) based on the knowledge that P(k) is true.



• We claim that $\sum_{i=1}^{n} i = \frac{n \cdot (n+1)}{2}$ holds for all $n \in \mathbb{N}$.



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• Induction basis (IB): We establish the truth of P(1):

$$\sum_{i=1}^{1} i = 1 = \frac{1 \cdot 2}{2}.$$



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$$\sum_{i=1}^k i = \frac{k \cdot (k+1)}{2}.$$



Proof (cont'd):

• Inductive step (IS): We have to prove P(k + 1) based on the induction hypothesis. That is, we have to prove

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$$\sum_{i=1}^{k+1} i = \left(\sum_{i=1}^{k} i\right) + (k+1)$$



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Theorem 79 (Strong Principle of Induction (S.P.I.))

Consider a predicate P over \mathbb{N} .

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and if

$$\forall k \in \mathbb{N} \ \left[\left(P(1) \land P(2) \land \ldots \land P(k) \right) \ \Rightarrow \ P(k+1) \right]$$

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Since

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all theorems that can be proved by W.P.I. can also be proved by S.P.I.



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But W.P.I. and S.P.I. are equivalent, at least from a theoretical point of view.



Theorem 80 (S.P.I. with Larger Base)

Consider a predicate P over \mathbb{N} , and let $m \in \mathbb{N}$.

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Proof: We define a new predicate P' over \mathbb{N} with

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and apply the standard S.P.I.

 We could also carry out induction for smaller base values. That is, induction. works for claims over \mathbb{N}_0 . (And even for negative base values!)

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• Thus, proving the base is mandatory!



• Several base cases alone do not suffice! For $n \in \mathbb{N}_0$, let

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- No! For instance, f(41) is not prime.
- Thus, proving the inductive step is truly mandatory!



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 - IS: Consider a set S of n + 1 cats, and let A and B be two subsets of S such that

$$|A| = |B| = n$$
 and $A \cup B = S$.

Using the induction hypothesis and the transitivity of the equivalence, we conclude that all cats of the set S have the same color!



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As nature shows, this "proof" is seriously flawed . . .



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$$2 \cdot (n+1) = 2 \cdot (k_1 + k_2) = 2 \cdot k_1 + 2 \cdot k_2 \stackrel{\text{i.H.}}{=} 0 + 0 = 0,$$

thus finishing the inductive "proof" . . .



- Suppose that some limited and non-uniform resource has to be distributed fairly among n receivers, for some $n \in \mathbb{N}$ with n > 1.
- E.g., a cake (with fruits, whipped cream, chocolate crumbs, icing, etc.) might have to be distributed fairly among *n* kids. Aka: "Cake Cutting Problem".



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Definition 81 (Fair distribution protocol)

A protocol for the distribution of a resource among n receivers is considered fair if each receiver gets at least 1/n-th of the resource (by his/her preferences), no matter what the preferences of the other receivers are and what the other receivers get.



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A protocol for the distribution of a resource among n receivers is considered *fair* if each receiver gets at least 1/n-th of the resource (by his/her preferences), no matter what the preferences of the other receivers are and what the other receivers get.

• How can we come up with a fair distribution protocol? Is there a general algorithm for fair cake cutting in the presence of *n* kids??



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Alice cuts the cake into two equal pieces (equal by her preferences).



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- If n = 2: Cut-and-choose distribution protocol.
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- If n > 2: Recursive application of the cut-and-choose distribution protocol.
 - The first n-1 kids cut the cake into n-1 pieces by applying the cut-and-choose distribution protocol recursively to n-2, n-3 etc. kids, thus each obtaining (hopefully) at least a fair 1/n-1 portion of the cake.



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 - 2 The n-th kid asks all other n-1 kids to cut his/her portion of the cake into n pieces such that the cutting is fair according to his/her preferences. (That is, according to each kid's preferences, each of the n pieces of his/her portion is equally desirable, for all of the first n-1 kids.)



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 - **1** The n-th kid walks around and collects one piece the most desirable piece according to his/her preferences! from all the other n-1 kids.



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Theorem 82

The recursive cut-and-choose distribution protocol is fair.



Proof of Thm. 82 by induction: Assume that the total cake is worth 1 for each kid.

I.B.: n := 2 Alice cut the cake into two pieces that are equally desirable (according to her preferences) and, thus, both worth 1/2. Hence, she will get one half of the cake (by her preferences), no matter how Bob behaves.



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Bob sees two pieces, one worth w_1 and the other one worth $1 - w_1$ (by his preferences). Trivially, either $w_1 \ge 1/2$ or $1 - w_1 \ge 1/2$.



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Hence, Bob can choose at least one half of the cake (according to his preferences), and both kids have no reason to complain about an unfair cutting.



- I.B.: n := 2 Alice cut the cake into two pieces that are equally desirable (according to her preferences) and, thus, both worth ¹/₂. Hence, she will get one half of the cake (by her preferences), no matter how Bob behaves.
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 Hence, Bob can choose at least one half of the cake (according to his preferences), and both kids have no reason to complain about an unfair cutting.
- **I.H.:** Assume that the recursive cut-and-choose cake cutting has been considered fair by the first k-1 kids, for $k \ge 3$ arbitrary but fixed. Hence, each of the first k-1 kids got a portion that is a least worth (according to the kid's preferences) $\frac{1}{k-1}$.



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- **I.B.:** n:=2 Alice cut the cake into two pieces that are equally desirable (according to her preferences) and, thus, both worth 1/2. Hence, she will get one half of the cake (by her preferences), no matter how Bob behaves. Bob sees two pieces, one worth w_1 and the other one worth $1-w_1$ (by his preferences). Trivially, either $w_1 \ge 1/2$ or $1-w_1 \ge 1/2$.
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$$\frac{w_1}{k} + \frac{w_2}{k} + \ldots + \frac{w_{k-1}}{k} = \frac{1}{k}(w_1 + w_2 + \ldots + w_{k-1}) = \frac{1}{k}.$$



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- However, this notion of cardinality becomes tricky for "infinite" sets.
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Definition 83 (Cardinality; Dt.: Mächtigkeit, Kardinalität)

The set *A* has *n* elements, aka *cardinality n*, for some $n \in \mathbb{N}$, if there exists a bijection from $\{1, 2, ..., n-1, n\}$ to *A*. The cardinality of *A* is denoted by |A|.



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A subset of a countably infinite set is a finite or a countably infinite set itself.



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Theorem 86 (Cantor&Schröder&Bernstein)

Consider two sets A and B. If there exist injective functions $f: A \to B$ and $g: B \to A$ between the sets A and B, then there exists a bijective function between A and B.



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Corollary 87

Consider three sets A, B and C. If $A \subseteq B \subseteq C$ and |A| = |C| then |A| = |B| = |C|.

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Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers
 - Construction of the Integers
 - Integral Powers
 - Divisibility and Prime Numbers
 - Quotient and Remainder
 - Congruences
 - Greatest Common Divisor
 - Chinese Remainder Theorem



Integers: Z

- Intuitive way to define the integers: $\mathbb{Z} := \mathbb{N}_0 \cup \{-n : n \in \mathbb{N}\}.$
- Thus, $\mathbb{Z} = \{0, \pm 1, \pm 2, \pm 3, \pm 4, \pm 5, \ldots\}.$
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- But what are the properties of the elements -n??
- And how could we define a + b and $a \cdot b$ for $a, b \in \mathbb{Z}$??
- In order to put $\mathbb Z$ on a more solid basis, we "extend" $\mathbb N$ to obtain $\mathbb Z$.



• Let $\cong_{\mathcal{Z}}$ be a relation over \mathbb{N}_0 such that

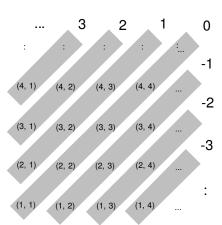
$$(a,b) \cong_{\mathbb{Z}} (c,d) :\Leftrightarrow a+d=c+b.$$



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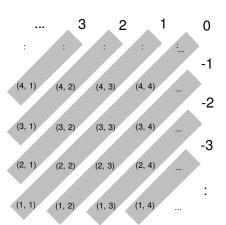
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• Easy to show: \cong_Z is an equivalence relation over \mathbb{N}_0 , with the equivalence classes shown below.



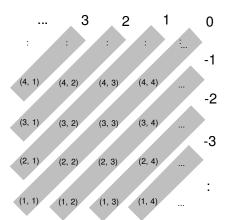


• We interprete $[(a,b)]_{\cong_{\mathbb{Z}}}$ as a-b.





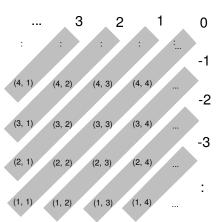
- We interprete $[(a,b)]_{\cong_{\mathbb{Z}}}$ as a-b.
- For $n \in \mathbb{N}$, the equivalence classes $[(n,0)]_{\cong_{\mathbb{Z}}}$ form the natural numbers, while $[(0,n)]_{\cong_{\mathbb{Z}}}$ form the negative integers.
- Zero is given by $[(0,0)]_{\cong_{\mathbb{Z}}}$.





Definition 88 (Integers)

The *integers* \mathbb{Z} are defined as $\mathbb{Z} := \{[(a,b)]_{\cong_Z} : a,b \in \mathbb{N}_0\}.$

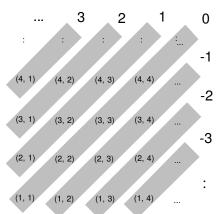




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• Furthermore, $\mathbb{Z}^+ := \mathbb{N}$ and $\mathbb{Z}_0^+ := \mathbb{N}_0$.





• It remains to define addition, multiplication and order on \mathbb{Z} . For $a, b, c, d \in \mathbb{N}_0$ we define an addition $+_Z$, a multiplication \cdot_Z and an order \leqslant_Z as follows:

$$\begin{array}{lcl} [(a,b)]_{\cong_{\mathbb{Z}}} +_{\mathbb{Z}} [(c,d)]_{\cong_{\mathbb{Z}}} & := & [(a+c,b+d)]_{\cong_{\mathbb{Z}}} \\ [(a,b)]_{\cong_{\mathbb{Z}}} \cdot_{\mathbb{Z}} [(c,d)]_{\cong_{\mathbb{Z}}} & := & [(a\cdot c+b\cdot d,a\cdot d+b\cdot c)]_{\cong_{\mathbb{Z}}} \\ [(a,b)]_{\cong_{\mathbb{Z}}} \leqslant_{\mathbb{Z}} [(c,d)]_{\cong_{\mathbb{Z}}} & :\Leftrightarrow & a+d\leqslant b+c \end{array}$$



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- It is easy to show that
 - addition, multiplication and order are well-defined,
 - the standard rules of arithmetic hold, with $[(0,0)]_{\cong_Z}$ as zero element ("zero"),
 - \leq_Z defines a total order on $\mathbb{N}_0 \times \mathbb{N}_0$.



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An integer is *positive* if it is greater than zero and *negative* if it is less than zero; zero is neither positive nor negative.



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Theorem 90

 \mathbb{Z} is a countably infinite set. That is, $|\mathbb{N}| = |\mathbb{Z}|$.



Definition 91 (Integral power, Dt.: ganzzahlige Potenz)

Consider $x \in F$ for a field $(F, +, \cdot)$, with additive neutral element e. For $n \in \mathbb{N}_0$, we define integral powers of x as follows:

$$x^n := \left\{ egin{array}{lll} 1 & & ext{if} & n=0 ext{ and } x
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Lemma 92

Let $(F, +, \cdot)$ be a field. Then, for all $x, y \in F$ and all $m, n \in \mathbb{Z}$,

$$\mathbf{x}^m \cdot \mathbf{x}^n = \mathbf{x}^{m+n}$$

$$x^n \cdot y^n = (x \cdot y)^n$$

Definition 93 (Divisor, Dt.: Teiler, Faktor)

Let $a, b \in \mathbb{Z}$ with $a \neq 0$. Then a divides b, denoted by $a \mid b$, if there exists $c \in \mathbb{Z}$ such that $b = c \cdot a$.

$$a \mid b :\Leftrightarrow \exists c \in \mathbb{Z} \ b = c \cdot a.$$



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In this case, we also say that b is a *multiple* of a, or a is a *divisor* or *factor* of b, or b is *divisible* by a. Otherwise we have $a \nmid b$. We have a *genuine divisor* if $a \mid b$ and $a \neq \pm 1$ and $a \neq \pm b$.



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Definition 94 (Even/odd, Dt.: gerade/ungerade)

A number $b \in \mathbb{Z}$ is said to be *even* if and only if $2 \mid b$; otherwise, b is *odd*.



Lemma 95



Lemma 95

Lemma 96

For all $a, b, c \in \mathbb{Z}$ and all $k \in \mathbb{Z} \setminus \{0\}$,

$$(a = b + c \land k \mid b) \Rightarrow (k \mid a \Leftrightarrow k \mid c).$$



Lemma 97

A number $a \in \mathbb{N}$ is divisible by

if its last digit is even, i.e., 0, 2, 4, 6 or 8;



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- if its last digit is even, i.e., 0, 2, 4, 6 or 8;
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- if it is divisible by two and three;



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- if its last two digits form a number that is divisible by four;
- if its last digit is 0 or 5;
- if it is divisible by two and three;
- if the hundreds digit is even and the number formed by the last two digits is divisible by eight, or if the hundreds digit is odd and the number formed by the last two digits plus four is divisible by eight;



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- if the hundreds digit is even and the number formed by the last two digits is divisible by eight, or if the hundreds digit is odd and the number formed by the last two digits plus four is divisible by eight;
- if the sum of its digits is divisible by nine;



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- if its last two digits form a number that is divisible by four;
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- if the hundreds digit is even and the number formed by the last two digits is divisible by eight, or if the hundreds digit is odd and the number formed by the last two digits plus four is divisible by eight;
- if the sum of its digits is divisible by nine;
- if its last digit is 0;



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- if the hundreds digit is even and the number formed by the last two digits is divisible by eight, or if the hundreds digit is odd and the number formed by the last two digits plus four is divisible by eight;
- if the sum of its digits is divisible by nine;
- if its last digit is 0;
- if the alternating sum of its digits is divisible by eleven;



Lemma 97

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- if it is divisible by three and four.
 - There also exist divisibility rules for seven but all of them are a bit ackward

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Proof of Lem. 97: We prove only the divisibility by three. Let $n \in \mathbb{N}$ and $a_0, a_1, \ldots, a_n \in \{0, 1, \ldots, 9\}$ such that

$$a=\sum_{i=0}^n a_i\cdot 10^i.$$



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We get

$$a = \sum_{i=0}^{n} a_i \cdot 10^i = \sum_{i=0}^{n} a_i \cdot (10^i - 1) + \sum_{i=0}^{n} a_i.$$



Proof of Lem. 97: We prove only the divisibility by three. Let $n \in \mathbb{N}$ and $a_0, a_1, \ldots, a_n \in \{0, 1, \ldots, 9\}$ such that

$$a=\sum_{i=0}^n a_i\cdot 10^i.$$

We get

$$a = \sum_{i=0}^{n} a_i \cdot 10^i = \sum_{i=0}^{n} a_i \cdot (10^i - 1) + \sum_{i=0}^{n} a_i.$$

Since

$$3 \mid \left(\sum_{i=0}^n a_i \cdot (10^i - 1)\right),\,$$

Lemma 96 implies that the number a is divisible by three if and only if

$$3 \mid \left(\sum_{i=0}^n a_i\right)$$
.



Prime Numbers

Definition 98 (Prime, Dt.: Primzahl)

A natural number $p \in \mathbb{N}$ is a *prime number*, or is *prime*, if $p \ge 2$ and if $p \in \mathbb{N}$ is divisible only by 1 and p itself. All other natural numbers $p \ge 2$ are called *composite*.



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Definition 99 (Prime factor, Dt.: Primfaktor)

A natural number $p \in \mathbb{N}$ is a *prime factor* of $n \in \mathbb{N}$ if p is prime and $p \mid n$. If p is a prime factor of n then its *multiplicity* (Dt.: Vielfachheit) is the largest exponent k for which $p^k \mid n$.



Lemma 100

Let $k \in \mathbb{N}$ and $a_1, a_2, \dots, a_k \in \mathbb{Z}$ and $p \in \mathbb{P}$. Then

$$p \mid a_1 \cdot a_2 \cdot \ldots \cdot a_k \quad \Leftrightarrow \quad (\exists (1 \leqslant j \leqslant k) \quad p \mid a_j).$$



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Corollary 101

If two products of primes are identical then the primes are identical up to the order in which they appear in the products.



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Theorem 102 (Fundamental Theorem of Arithmetic)

Every natural number n > 1 is representable uniquely in the form

$$n = p_1^{m_1} \cdot p_2^{m_2} \cdot \ldots \cdot p_k^{m_k},$$

where $p_1 < \ldots < p_k$ are primes and $m_j \in \mathbb{N}$ are multiplicities for every $j = 1, \ldots, k$.



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Corollary 103

There are infinitely many prime numbers.

Definition 104 (Mersenne prime)

A *Mersenne number* is of the form 2^n-1 for $n\in\mathbb{N}$. A *Mersenne prime* is a Mersenne number which is prime.



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- INT_MAX (in C/C++) is the eight Mersenne prime: $2147483647 = 2^{31} 1$.
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- Several unsolved problems related to Mersenne numbers:
 - Since $2^{11} 1 = 2047 = 23 \cdot 89$, not all Mersenne numbers are primes!
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Lemma 105

If $2^n - 1$ is prime for some $n \in \mathbb{N}$ then n is prime.



Conjecture 106 (Goldbach 1742, "weak version" or "ternary conjecture")



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- The strong version of this conjecture implies the weak version: If $n \in \mathbb{N}$, with $n \ge 7$, is odd then n' := n 3 is even with n' > 3. Hence, if n' can be written as the sum of two primes, then n can be written as the sum of three primes.



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- Also by distributed computing, in 2013 Harald Helfgott and David Platt verified the weak Goldbach conjecture up to (roughly) 8 · 10³⁰.



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Conjecture 108 (Polignac, 1849)

For every natural number k there exist infinitely many numbers p such that p and p + 2k are consecutive primes.



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- In November 2013, James Maynard reduced this bound to 600.
- This bound seems to have been further reduced to 246 by the Polymath project.



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- E.g., (3,4,5) is an integer solution triple for $a^2 + b^2 = c^2$.



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For every natural number n > 2, the Diophantine equation $a^n + b^n = c^n$ has no solution $(a, b, c) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$.

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- Finally proved by Andrew Wiles in 1993; a gap in the proof was fixed by Wiles and his former student Richard Taylor; the full proof was published in 1995.



Quotient and Remainder

Lemma 110

Let $a\in\mathbb{N}$ and $b\in\mathbb{Z}$. Then there exist a unique *quotient* $q\in\mathbb{Z}$ and a unique *remainder* $r\in\mathbb{N}_0$ such that

$$b = a \cdot q + r$$
 and $0 \le r < a$.

We will use the operators div and mod for computing the quotient and remainder.
 That is, q and r of Lemma 110 are given by q := b div a and r := b mod a.



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- IEEE 754 defines a remainder function based on the round-to-nearest convention.

Warning

If one or both of *a* and *b* are allowed to be negative integers then the sign of the remainder may differ among different implementations!



• We know that $25 = (11001)_2$, i.e., $(11001)_2$ is the base-two representation of $25 = (25)_{10}$. (After all, $25 = 1 \cdot 2^4 + 1 \cdot 2^3 + 0 \cdot 2^2 + 0 \cdot 2^1 + 1 \cdot 2^0$.)



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and so on until some $q_i = 0$.

• E.g.,

$$25 = 12 \cdot 2 + 1
12 = 6 \cdot 2 + 0
6 = 3 \cdot 2 + 0
3 = 1 \cdot 2 + 1
1 = 0 \cdot 2 + 1$$

and therefore $25 = (11001)_2$.



• Introduced by Carl Friedrich Gauss (1777–1855) in 1801.

Definition 111 (Congruence, Dt.: Kongruenz)

Let $a, b \in \mathbb{Z}$ and $m \in \mathbb{N}$. We say that a is congruent to b modulo m, and write

$$a \equiv_m b$$
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if a - b is divisible by m. The term $a \equiv_m b$ is called a *congruence*.



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$$a \equiv_m b :\Leftrightarrow m \mid (a-b)$$
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Lemma 112

For all $a, b \in \mathbb{Z}$ and $m \in \mathbb{N}$, we have $a \equiv_m b$ if and only if a and b have the same remainder after dividing by m, i.e., if and only if $a \mod m = b \mod m$.



$$38 \equiv_{12} 2$$

$$-3 \equiv_{\mathbf{5}} 2$$

$$0 \equiv_3 3$$

$$8 \equiv_3 2$$

$$7 \equiv_3 1$$

$$7 \equiv_3 -8$$

$$even + even \equiv_2 even$$

$$even + odd \equiv_2 odd$$

$$\textit{odd} + \textit{odd} \equiv_2 \textit{even}$$

$$even \cdot even \equiv_2 even$$

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Lemma 113

For $m \in \mathbb{N}$, the relation \equiv_m is an equivalence relation on \mathbb{Z} , i.e., for all $a, b, c \in \mathbb{Z}$,

reflexivity $a \equiv_m a$,

symmetry if $a \equiv_m b$ then $b \equiv_m a$, and

transitivity if $a \equiv_m b$ and $b \equiv_m c$ then $a \equiv_m c$

hold.



Residues and \mathbb{Z}_m

Lemma 114

For $m \in \mathbb{N}$, the relation \equiv_m is a congruence relation on \mathbb{Z} , i.e., it respects addition, subtraction, and multiplication: Let $a,b,c,d\in\mathbb{Z}$ and $m\in\mathbb{N}$, and suppose that

$$a \equiv_m b$$
 and $c \equiv_m d$.

Then

$$a+c\equiv_m b+d$$
 and $a-c\equiv_m b-d$ and $a\cdot c\equiv_m b\cdot d$.



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Definition 115 (Residue, Dt.: Residuum, Restklasse)

Let $m \in \mathbb{N}$ with $m \geqslant 2$. The equivalence classes of \mathbb{Z} modulo m are called *residues* (or remainders) modulo m. For $a \in \mathbb{Z}$, its equivalence class modulo m is denoted by $[a]_m$. The set of residues modulo m is denoted by \mathbb{Z}_m or $\mathbb{Z}/m\mathbb{Z}$.



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Lemma 116

Let $m \in \mathbb{N}$ with $m \ge 2$. Then $\mathbb{Z}_m = \{ [a]_m : a \in \mathbb{N}_0 \land a < m \}$.



Definition 117 (Arithmetic on \mathbb{Z}_m **)**

Let $m \in \mathbb{N}$ with $m \ge 2$, and $[a]_m, [b]_m \in \mathbb{Z}_m$. On \mathbb{Z}_m we define an addition $+_m$ and a multiplication \cdot_m as follows.

$$[a]_m +_m [b]_m := [a+b]_m$$

 $[a]_m \cdot_m [b]_m := [a \cdot b]_m$



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 $[a]_m \cdot_m [b]_m := [a \cdot b]_m$

Lemma 118

Let $m \in \mathbb{N}$ with $m \ge 2$. Then addition $+_m$ and multiplication \cdot_m on \mathbb{Z}_m are well-defined. Furthermore, $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a commutative ring.



Definition 117 (Arithmetic on \mathbb{Z}_m)

Let $m \in \mathbb{N}$ with $m \ge 2$, and $[a]_m, [b]_m \in \mathbb{Z}_m$. On \mathbb{Z}_m we define an addition $+_m$ and a multiplication \cdot_m as follows.

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Lemma 118

Let $m \in \mathbb{N}$ with $m \geqslant 2$. Then addition $+_m$ and multiplication \cdot_m on \mathbb{Z}_m are well-defined. Furthermore, $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a commutative ring.

• Often the notation $[a]_m$ is simplified by omitting the modulus m, i.e., by writing [a], or even by simply writing a if it is clear that $a \in \mathbb{Z}_m$. Similarly for $+_m$ and \cdot_m .



Definition 117 (Arithmetic on \mathbb{Z}_m **)**

Let $m \in \mathbb{N}$ with $m \ge 2$, and $[a]_m, [b]_m \in \mathbb{Z}_m$. On \mathbb{Z}_m we define an addition $+_m$ and a multiplication \cdot_m as follows.

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Let $m \in \mathbb{N}$ with $m \geqslant 2$. Then addition $+_m$ and multiplication \cdot_m on \mathbb{Z}_m are well-defined. Furthermore, $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a commutative ring.

- Often the notation $[a]_m$ is simplified by omitting the modulus m, i.e., by writing [a], or even by simply writing a if it is clear that $a \in \mathbb{Z}_m$. Similarly for $+_m$ and \cdot_m .
- It is also common to write

a mod m

instead of

$$[a]_m$$
.



Theorem 119 (Fermat's Little Theorem)

If $p \in \mathbb{N}$ is prime then $a^p \equiv_p a$ for every $a \in \mathbb{N}$.



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• If *a* is not divisible by *p* then this yields $a^{p-1} \equiv_p 1$. In particular, this congruence holds for all $1 \le a \le p-1$.



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- If a is not divisible by p then this yields $a^{p-1} \equiv_p 1$. In particular, this congruence holds for all $1 \le a \le p-1$.
- Hence, if $a^{n-1} \not\equiv_n 1$ for $n \in \mathbb{N}$ (and $1 \leqslant a \leqslant n-1$) then n is composite, i.e., not a prime.



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- Fermat Primality Test for $n \in \mathbb{N}$:
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 - **3** Otherwise, repeat the test for some other value of $a \in \{2, 3, ..., n-2\}$.



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One can prove that the probability for incorrectly classifying n as prime goes to zero (in most cases) as the number of tests is increased.



• Since computers cannot flip a coin to obtain a random result, one resorts to algorithms that generate "random" numbers: pseudo-random number generators.



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- Linear congruential generators (LCG, [Lehmer 1954]) have been well studied, are easy to implement and used frequently.
- They generate a sequence of non-negative integers less than some specified modulus $m \in \mathbb{N}$ according to the following recursive definition:

$$x_{n+1} := (a \cdot x_n + c) \mod m$$
,

where

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m \in \mathbb{N} with m > 1 ..... modulus, a \in \mathbb{N} with a < m ..... multiplier, c \in \mathbb{N}_0 with c < m .... increment, x_0 \in \mathbb{N}_0 with x_0 < m .... seed.
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 E.g., m := 15, a := 1, c := 4 and x₀ := 2 yields the following sequence of numbers:

2 6 10 14 3 7 11 0 4 8 12 1 5 9 13 2 6 ...



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• GCC/glibc: $m := 2^{31} - 1$, a := 1103515245, c := 12345.



- Since computers cannot flip a coin to obtain a random result, one resorts to algorithms that generate "random" numbers: pseudo-random number generators.
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• E.g., m := 15, a := 1, c := 4 and $x_0 := 2$ yields the following sequence of numbers:

```
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• GCC/glibc: $m := 2^{31} - 1$, a := 1103515245, c := 12345. More advanced pseudo-random number generators exist, e.g., Mersenne twister.



Greatest Common Divisor

Lemma 120

Let $a, b \in \mathbb{N}$. Then there exists a unique $n \in \mathbb{N}$ such that

- \bigcirc $n \mid a$ and $n \mid b$, and
- ② for all $m \in \mathbb{N}$, if $m \mid a$ and $m \mid b$ then $m \leqslant n$.



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Definition 121 (Greatest common divisor, Dt.: größter gemeinsamer Teiler (ggT))

Let $a, b \in \mathbb{N}$. The unique number $n \in \mathbb{N}$ that exists according to Lem. 120 is called *greatest common divisor* of a and b, and is denoted by gcd(a, b).



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Definition 123 (Pairwise relatively prime)

A set S of natural numbers is called *pairwise relatively prime* (or *pairwise coprime* or *mutually coprime*) if all pairs of numbers a and b in S, with $a \neq b$, are relatively prime.

Lemma 124 (Bézout's Identity)

Let $a, b \in \mathbb{N}$. Then there exist $x, y \in \mathbb{Z}$ such that $gcd(a, b) = a \cdot x + b \cdot y$. Conversely, the smallest positive number $a \cdot x + b \cdot y$, for all $x, y \in \mathbb{Z}$, equals gcd(a, b).



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• That is, $gcd(a,b) = min (\mathbb{N} \cap \{a \cdot x + b \cdot y \colon x,y \in \mathbb{Z}\}).$



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- For $a, b, d \in \mathbb{Z}$ given, the identity $d = a \cdot x + b \cdot y$ over $\mathbb{Z} \times \mathbb{Z}$ is called a *linear Diophantine equation* in x and y.
- Lemma 124 was first stated by Étienne Bézout (1730–1783), and numbers $x, y \in \mathbb{Z}$ with $gcd(a, b) = a \cdot x + b \cdot y$ are called Bézout numbers.



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- Note: Bézout numbers are not unique! For instance, gcd(10, 15) = 5, and 10x + 15y = 5 has the solutions x := -1 and y := 1, and x := 2 and y := -1.



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- Note: Bézout numbers are not unique! For instance, gcd(10, 15) = 5, and 10x + 15y = 5 has the solutions x := -1 and y := 1, and x := 2 and y := -1.

Corollary 125

The numbers $a,b\in\mathbb{N}$ are relatively prime if and only if the linear Diophantine equation $a\cdot x+b\cdot y=1$ has a solution, i.e., if and only if there exist $x,y\in\mathbb{Z}$ such that $a\cdot x+b\cdot y=1$.



Euclidean Algorithm for GCD Computation

Theorem 126 (Euclidean Algorithm)

The following algorithm computes gcd(a,b) for $a,b \in \mathbb{N}_0$ with a > b.

```
function gcd(a,b)

precondition: a,b \in \mathbb{N}_0 with a > b.

postcondition: t = gcd(a,b)

while b > 0 do

t \leftarrow b

b \leftarrow a \mod b

a \leftarrow t

end while

t \leftarrow a
```



```
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 We want to compute the gcd of 78 and 99. Hence, b := 78 and a := 99 = 1 · 78 + 21.



```
function \gcd(a,b)
precondition: a,b \in \mathbb{N}_0 with a > b.
postcondition: t = \gcd(a,b)
while b > 0 do
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end while
t \leftarrow a
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We want to compute the gcd of 78 and 99. Hence, b := 78 and
 a := 99 = 1 · 78 + 21. We get after different passes through the loop:

after 1st pass:
$$t = 78$$
, $b = 21$, $a = 78 = 3 \cdot 21 + 15$



```
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while b > 0 do

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```
after 1st pass: t = 78, b = 21, a = 78 = 3 \cdot 21 + 15
after 2nd pass: t = 21, b = 15, a = 21 = 1 \cdot 15 + 6
```



```
function \gcd(a,b)

precondition: a,b \in \mathbb{N}_0 with a > b.

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after 1st pass: t = 78, b = 21, a = 78 = 3 \cdot 21 + 15
after 2nd pass: t = 21, b = 15, a = 21 = 1 \cdot 15 + 6
after 3rd pass: t = 15, b = 6, a = 15 = 2 \cdot 6 + 3
```



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after 1st pass: t = 78, b = 21, a = 78 = 3 \cdot 21 + 15
after 2nd pass: t = 21, b = 15, a = 21 = 1 \cdot 15 + 6
after 3rd pass: t = 15, b = 6, a = 15 = 2 \cdot 6 + 3
after 4th pass: t = 6, b = 3, a = 6 = 2 \cdot 3 + 0
```



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function gcd(a,b)

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postcondition: t = gcd(a,b)

while b > 0 do

t \leftarrow b

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We want to compute the gcd of 78 and 99. Hence, b := 78 and
 a := 99 = 1 · 78 + 21. We get after different passes through the loop:

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after 1st pass: t = 78, b = 21, a = 78 = 3 \cdot 21 + 15
after 2nd pass: t = 21, b = 15, a = 21 = 1 \cdot 15 + 6
after 3rd pass: t = 15, b = 6, a = 15 = 2 \cdot 6 + 3
after 4th pass: t = 6, b = 3, a = 6 = 2 \cdot 3 + 0
after 5th pass: t = 3, b = 0, a = 3
```

• Hence, $t = 3 = \gcd(78, 99)$.



Theorem 127

Let $m \in \mathbb{N}$ with $m \geqslant 2$. An element $[a]_m \in \mathbb{Z}_m$ has a multiplicative inverse if and only if a is relatively prime to m.



Theorem 127

Let $m \in \mathbb{N}$ with $m \ge 2$. An element $[a]_m \in \mathbb{Z}_m$ has a multiplicative inverse if and only if a is relatively prime to m.

Corollary 128

Let $m \in \mathbb{N}$ with $m \ge 2$. The ring $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a (finite) field if and only if m is prime.



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• If m is not prime then $(\mathbb{Z}_m, +_m, \cdot_m)$ may contain non-trivial zero divisors.



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• If m is not prime then $(\mathbb{Z}_m, +_m, \cdot_m)$ may contain non-trivial zero divisors.

Lemma 129

Let $m \in \mathbb{N}$ with $m \geqslant 2$ and $[a]_m \in \mathbb{Z}_m$ such that m and a are relatively prime. Let $x, y \in \mathbb{Z}$ such that $a \cdot x + m \cdot y = 1$. Then $[a]_m \cdot_m [x]_m = [1]_m$, i.e., $[x]_m$ is the multiplicative inverse element for $[a]_m$.



Euclidean Algorithm Revisited

Recursive formulation of the Euclidean Algorithm.

```
 \begin{split} & \textbf{function} \ \ \mathsf{gcd\_recursive}(a,b) \\ & \textbf{precondition:} \ \ a,b \in \mathbb{N} \ \ \mathsf{with} \ \ a > b. \\ & \textbf{if} \ \ (a \bmod b) = 0 \ \ \textbf{then} \\ & \textbf{return} \ \ b \\ & \textbf{else} \\ & \textbf{return} \ \ \mathsf{gcd\_recursive}(b,a \bmod b) \\ & \textbf{end if} \\ \end{split}
```



Extended Euclidean Algorithm for GCD Computation

Theorem 130 (Extended Euclidean Algorithm)

```
The following algorithm computes x,y\in\mathbb{Z} and d\in\mathbb{N} such that \gcd(a,b)=d=a\cdot x+b\cdot y for a,b\in\mathbb{N}_0 with a>b.

function \gcd\_\text{extended}(a,b) precondition: a,b\in\mathbb{N}_0 with a>b.

postcondition: (d,x,y)\in\mathbb{N}\times\mathbb{Z}\times\mathbb{Z} such that \gcd(a,b)=d=a\cdot x+b\cdot y if (a\bmod b)=0 then return (b,0,1) else (d,x,y)\leftarrow\gcd\_\text{extended}(b,a\bmod b) return (d,y,x-y\cdot (a\ \text{div}\ b)) end if
```



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```
function \gcd_{a,b}

postcondition: (d,x,y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that \gcd(a,b) = d = a \cdot x + b \cdot y

if (a \bmod b) = 0 then

return (b,0,1)

else

(d,x,y) \leftarrow \gcd_{a,b} = \gcd_{a,b}

(d,x,y) \leftarrow \gcd_{a,b} = \gcd_{a,b}

return (d,y,x-y\cdot (a \bmod b))

end if
```



```
function gcd\_extended(a,b)

postcondition: (d,x,y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a,b) = d = a \cdot x + b \cdot y

if (a \bmod b) = 0 then

return (b,0,1)

else

(d,x,y) \leftarrow gcd\_extended(b,a \bmod b)

return (d,y,x-y\cdot(a \ div \ b))

end if
```

• We want to compute $x, y \in \mathbb{Z}$ and $d \in \mathbb{N}$ such that gcd(99, 78) = d = 99x + 78y.



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```
function \gcd_a extended(a,b) postcondition: (d,x,y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that \gcd(a,b) = d = a \cdot x + b \cdot y if (a \mod b) = 0 then return (b,0,1) else (d,x,y) \leftarrow \gcd_a extended(b,a \mod b) return (d,y,x-y\cdot (a \operatorname{div} b)) end if
```

а	b	a div b	a mod b	d	X	У
99	78	1	21			
78	21	3	15			



```
function gcd_extended(a,b)

postcondition: (d,x,y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that \gcd(a,b) = d = a \cdot x + b \cdot y

if (a \mod b) = 0 then

return (b,0,1)

else

(d,x,y) \leftarrow \gcd_extended(b,a \mod b)

return (d,y,x-y\cdot(a \text{ div }b))

end if
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а	b	a div b	a mod b	d	X	У
99	78	1	21			
78	21	3	15			
21	15	1	6			



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78	21	3	15			
21	15	1	6			
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а	b	a div b	a mod b	d	X	У
99	78	1	21			
78	21	3	15	3	3	-11
21	15	1	6	3	-2	3
15	6	2	3	3	1	-2
6	3	_	0	3	0	1



```
function \gcd_-extended(a,b) postcondition: (d,x,y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that \gcd(a,b) = d = a \cdot x + b \cdot y if (a \bmod b) = 0 then return (b,0,1) else (d,x,y) \leftarrow \gcd_-extended(b,a \bmod b) return (d,y,x-y\cdot (a \ \text{div}\ b)) end if
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а	b	a div b	a mod b	d	X	У
99	78	1	21	3	-11	14
78	21	3	15	3	3	-11
21	15	1	6	3	-2	3
15	6	2	3	3	1	-2
6	3	_	0	3	0	1



• We want to compute $x, y \in \mathbb{Z}$ and $d \in \mathbb{N}$ such that gcd(99, 78) = d = 99x + 78y.

а	b	a div b	a mod b	d	X	У
99	78	1	21	3	-11	14
78	21	3	15	3	3	-11
21	15	1	6	3	-2	3
15	6	2	3	3	1	-2
6	3	_	0	3	0	1

• Hence, $gcd(99, 78) = -11 \cdot 99 + 14 \cdot 78 = -1089 + 1092 = 3$.



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 - Then the soldiers should form groups of 7 and report back the number of soldiers that could not join a group consisting of 7 soldiers.
 - Then the soldiers should form groups of 11 and report back the number of soldiers that did not end up in a group consisting of 11 soldiers.



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 - **⑤** ...
- Based on this information he was able to figure out the number n of soldiers in his army.



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 - Then the soldiers should form groups of 11 and report back the number of soldiers that did not end up in a group consisting of 11 soldiers.
 - **⑤** ...
- Based on this information he was able to figure out the number n of soldiers in his army.
- Indeed, a mathematical solution was provided by the Chinese mathematician Sun Tzu sometime in the third to fifth century, and republished by Qin Jiushao in 1247!







 $n \mod 3 = 1$







$$n \mod 3 = 1$$





$$n \mod 5 = 2$$







 $n \mod 3 = 1$



9 9

 $n \mod 5 = 2$



9 (9

 $n \mod 7 = 2$



Theorem 131 (Chinese Remainder Theorem, Dt.: Chinesischer Restsatz)

If, for some $k \in \mathbb{N}$, the numbers $m_1, m_2, \cdots, m_k \in \mathbb{N}$ are pairwise relatively prime, then the following system of simultaneous congruences has an integer solution b for all $a_1, a_2, \ldots, a_k \in \mathbb{Z}$ given:

$$\begin{vmatrix}
b \equiv_{m_1} a_1 \\
b \equiv_{m_2} a_2 \\
\vdots \\
b \equiv_{m_k} a_k
\end{vmatrix} (*)$$

Furthermore, all solutions of (*) are congruent modulo $m := \prod_{i=1}^k m_i$. That is, the solution is unique if constrained to $\{1, 2, \dots, m\}$.



Proof: We show the existence of an integer solution. Consider $i \in \mathbb{N}$ with $1 \le i \le k$. Since m_1, m_2, \dots, m_k are pairwise relatively prime, $\gcd(\frac{m}{m_i}, m_i) = 1$. Using the extended Euclidean algorithm (Thm. 130), we can find integers x_i and y_i such that

$$x_i \cdot m_i + y_i \cdot \frac{m}{m_i} = 1. \tag{*}$$



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$$x_i \cdot m_i + y_i \cdot \frac{m}{m_i} = 1. \tag{*}$$

Let $b_i := y_i \cdot \frac{m}{m_i}$. Equation (\star) guarantees that the remainder of b_i when divided by m_i is 1. On the other hand, for $j \neq i$ every m_i divides b_i evenly. Thus,

$$b_i \equiv_{m_i} 1$$
 and $b_i \equiv_{m_j} 0$ for all j with $j \neq i$ and $1 \leqslant j \leqslant k$.



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Since congruences respect multiplication, we get

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Thus, one solution of the simultaneous congruences is given by

$$b:=\sum_{i=1}^k a_i\cdot b_i.$$



- The Emperor collected the following information:
 - When the soldiers formed groups of 3, one soldier was left out.
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- That is, since $a_1 = 1$, $a_2 = 2$, $a_3 = 2$ and $m_1 = 3$, $m_2 = 5$, $m_3 = 7$ and $m = 3 \cdot 5 \cdot 7 = 105$:

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• Hence, we are to find $x_1, x_2, x_3, y_1, y_2, y_3 \in \mathbb{Z}$ such that

$$3x_1 + 35y_1 = 1$$
 $5x_2 + 21y_2 = 1$ $7x_3 + 15y_3 = 1$.



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$$3x_1 + 35y_1 = 1$$
 $5x_2 + 21y_2 = 1$ $7x_3 + 15y_3 = 1$.

• We have $x_1 := 12$, $y_1 := -1$, $x_2 := -4$, $y_2 := 1$, $x_3 := -2$, $y_3 := 1$ and, thus,

$$n = (35 \cdot (-1) \cdot 1 + 21 \cdot 1 \cdot 2 + 15 \cdot 1 \cdot 2) \mod 105 = 37 \mod 105 = 37.$$



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- Secret sharing refers to the distribution of information related to a secret (e.g., a number) among a group of receivers such that the secret can only be reconstructed if all or, at least, a large percentage of the receivers cooperate.
- Ideally, the information received by one individual receiver shall be of no (or very little) help for the receiver to obtain the secret without the help of the others.



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- A secret sharing scheme is called a (t, n) threshold scheme, or t-out-of-n scheme, if at least t of the n receivers have to cooperate. (Of course, $t \le n$.)
- Typically, *t* is large relative to *n* but not identical to *n*.



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- Typically, *t* is large relative to *n* but not identical to *n*.
- Several different variants of schemes for secret sharing are used in practice.
- At least two published schemes rely on the Chinese Remainder Theorem 131.
- We sketch the very basic idea of a scheme based on the Chinese Remainder Theorem 131. (In our simple scheme we have t := n.)



• Suppose that the number 1234 is the secret *b* to be shared by five receivers.



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- We choose

$$m_1 := 2,$$
 $m_2 := 3,$ $m_3 := 5,$ $m_4 := 7,$ $m_5 := 11.$

Note that

$$m := \prod_{i=1}^{5} m_i = 2 \cdot 3 \cdot 5 \cdot 7 \cdot 11 = 2310 > 1234.$$



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• Now consider $a_i := 1234 \mod m_i$, for $1 \le i \le 5$. This gives us the numbers

$$a_1 = 0,$$
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- Note that each individual receiver has gained little information about the secret b.

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- The numbers m_i and a_i are passed to the *i*-th receiver.
- Note that each individual receiver has gained little information about the secret b.
- Rather, in our simple approach, all five receivers need to cooperate in order to recover b: They have to solve the following set of five congruences:

$$b \equiv_2 0$$
 $b \equiv_3 1$ $b \equiv_5 4$ $b \equiv_7 2$ $b \equiv_{11} 2$



• The five receivers have to solve the following set of five congruences:

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 Since a₁ = 0, we need to solve only four congruences and get the following four Diophantine equations.

$$3x_2 + 770y_2 = 1$$
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Solving these equations yields the following solutions:

$$x_2 := 257, y_2 := -1$$
 $x_3 := 185, y_3 := -2$ $x_4 := -47, y_4 := 1$ $x_5 := -19, y_5 := 1$



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Hence, the secret sought is recovered as

$$b = (-1) \cdot 770 \cdot 1 + (-2) \cdot 462 \cdot 4 + 1 \cdot 330 \cdot 2 + 1 \cdot 210 \cdot 2 = -3386 \equiv_{2310} 1234.$$



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- One way to carry out integer arithmetic with large integers is to apply modulo arithmetic and the Chinese Remainder Theorem 131:
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 - **3** Represent an integer n < m by its k remainders n_1, n_2, \ldots, n_k upon division by m_1, m_2, \ldots, m_k .



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 - **131** Recover the actual result by applying the Chinese Remainder Theorem 131.



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 - Recover the actual result by applying the Chinese Remainder Theorem 131.
- This approach works as long as all intermediate results are less than *m*.
- Advantages:
 - One can use (mostly) standard arithmetic to handle integers larger than those normally handled.
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- Advantages:
 - One can use (mostly) standard arithmetic to handle integers larger than those normally handled.
 - One can run the computations for the different remainders in parallel, thus speeding up the computation.
- Standard choices for the modules are numbers of the form $2^{i} 1$:
 - One can prove $gcd(2^{i} 1, 2^{i} 1) = 2^{gcd(i,j)} 1$, which makes it easy to ensure that the modules are relatively prime.

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- Suppose that we want to limit our arithmetic operations to numbers less than 12.
- We choose the five modules

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and remember that $m := m_1 \cdot m_2 \cdot m_3 \cdot m_4 \cdot m_5 = 2310$.

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$$\begin{array}{l} (0,1,4,2,2) + (0,1,0,6,10) = (0 \text{ mod } 2,2 \text{ mod } 3,4 \text{ mod } 5,8 \text{ mod } 7,12 \text{ mod } 11) \\ &= (0,2,4,1,1). \end{array}$$



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 Thus, b := 1234 + 1000 is uniquely determined as the solution of the following set of five congruences:

$$b \equiv_2 0$$
 $b \equiv_3 2$ $b \equiv_5 4$ $b \equiv_7 1$ $b \equiv_{11} 1$



- 4
 - **Numbers and Basics of Number Theory**
 - Algebraic Structures
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 - Rational Numbers
 - Construction of the Rational Numbers
 - Properties
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Definition 132 (Rational equivalence)

On $\mathbb{Z} \times \mathbb{N}$ we define the binary relation \cong_Q as follows:

$$(p_1,q_1)\cong_{Q}(p_2,q_2)\quad :\Leftrightarrow\quad p_1\cdot q_2=p_2\cdot q_1.$$



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Definition 134 (Rational numbers)

The rational numbers $\mathbb Q$ are defined as

$$\mathbb{Q}:=\{[(p,q)]_{\cong_Q}:\ p\in\mathbb{Z},q\in\mathbb{N}\}.$$



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The canonical representative of $[(p,q)]_{\cong_Q}$ is denoted by $\frac{p'}{q'}$, where p':=p div $\gcd(|p|,q)$ and q':=q div $\gcd(|p|,q)$.



• It is easy to define an addition $+_Q$, multiplication \cdot_Q and order \leqslant_Q on $\mathbb Q$ that turns $(\mathbb Q,+,\cdot)$ into a totally ordered field. E.g.,

$$[(p_1,q_1)]_{\cong_Q} +_Q [(p_2,q_2)]_{\cong_Q} := [(p_1 \cdot q_2 + p_2 \cdot q_1, q_1 \cdot q_2)]_{\cong_Q}$$



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• Of course, it is standard to simplify the notation and write

$$\frac{p}{q}$$
 instead of $[(p,q)]_{\cong_Q}$.

But keep in mind that fractions are equivalence classes. Thus,

$$(1,3) \cong_Q (3,9) \cong_Q (3000,9000)$$
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 In the sequel we resort to standard knowledge and deal with rational numbers as we learned in school. (However, this could be formalized based on Def. 134!)



Theorem 135

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Proof: Suppose that there exists $x \in \mathbb{Q}$ such that $x^2 = 2$. Hence, there exist $p \in \mathbb{Z}$ and $q \in \mathbb{N}$ such that

$$\gcd(|p|,q)=1$$
 and $2=\left(\frac{p}{q}\right)^2=\frac{p^2}{q^2}$.



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The second equation is equivalent to $2q^2 = p^2$, implying that $p^2 \equiv_2 0$, and, thus, also $p \equiv_2 0$.



Properties: \bigcirc Is Not Complete

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• Hence, $\sqrt{2} \notin \mathbb{Q}$.

Lemma 136

There exists a rational number between any two distinct rational numbers.



Properties: \mathbb{Q} Is Countably Infinite

Theorem 137

 $\ensuremath{\mathbb{Q}}$ is a countably infinite set.



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Properties: Q Is Countably Infinite

Theorem 137

Q is a countably infinite set.

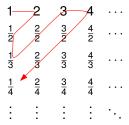
Proof by Cantor: Construct a bijection between $\mathbb N$ and $\mathbb Z\times\mathbb N$ (as a "superset" of $\mathbb Q$).

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This gives the sequence $1, 2, \frac{1}{2}, \frac{1}{3}, \frac{2}{2}, 3, \dots$

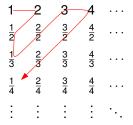


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This gives the sequence $1, 2, \frac{1}{2}, \frac{1}{3}, \frac{2}{2}, 3, \dots$ Now start with zero and include every number's negative number, thus obtaining a systematic enumeration of $\mathbb{Z} \times \mathbb{N}$:

$$0 \quad 1 \quad -1 \quad 2 \quad -2 \quad \frac{1}{2} \quad -\frac{1}{2} \quad \frac{1}{3} \quad -\frac{1}{3} \quad \frac{2}{2} \quad -\frac{2}{2} \quad 3 \quad -3 \quad \dots$$

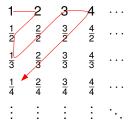


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Numbering this sequence yields a bijection from $\mathbb N$ onto $\mathbb Z\times\mathbb N$, and Cor. 87 implies the claim.

4

Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers
- Rational Numbers
- Real Numbers
 - Decimal Notation
 - Properties and Cardinality
- More Proof Techniques



Real Numbers: \mathbb{R}

- Intuitively, the reals comprise both rational and irrational numbers like $\sqrt{2}$ or π .
- A formal introduction of the reals, \mathbb{R} , based on \mathbb{Q} e.g., based on Dedekind cuts or based on equivalence classes of Cauchy sequences is beyond the scope of this lecture.



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- Convenient notations for intervals of real numbers:

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\forall a, b \in \mathbb{R} \quad [a, b] := \{x \in \mathbb{R} : a \leqslant x \leqslant b\};
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- Floor and ceiling function (Dt.: Ab- und Aufrundungsfunktion): For $x \in \mathbb{R}$,

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[x] := \max\{k \in \mathbb{Z} : k \leqslant x\},[x] := \min\{k \in \mathbb{Z} : k \geqslant x\}.
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- Gauß introduced the square-bracket notation [x] ("Gaussklammer") in 1808. The names "floor" and "ceiling" and the corresponding notations were introduced by Iverson in 1962 in his book on APL.
- We have [x] = |x| for all $x \in \mathbb{R}$.



Definition 138 (Decimal representation, Dt.: Dezimalzahl)

A real number $x \in \mathbb{R}_0^+$ is in *decimal representation* (or a *decimal number*) if it is represented as a sum of (negative) powers of ten:

$$x = x_0 + \sum_{i=1}^{\infty} \frac{x_i}{10^i}$$
, with an integer part $x_0 \in \mathbb{N}_0$ and with $0 \le x_i \le 9$ for all $i \in \mathbb{N}$.



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The decimal representation is *finite* if, for some $n_0 \in \mathbb{N}_0$, we have $x_i = 0$ for all $i \ge n_0$.



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• It is straightforward to extend Def. 138 to negative reals.



Definition 139 (Recurring decimal, Dt.: periodische Dezimalzahl)

A decimal representation of a real number is a *recurring decimal* (or *repeating decimal*) if it becomes periodic at some point: a finite subsequence of the digits after the decimal separator is repeated indefinitely.



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Recurring decimals, e.g.,

$$\frac{1}{3}=0.333\cdots$$

or

$$\frac{1}{7} = 0.142857142857142857 \cdots$$

are written as $0.\overline{3}$ or $0.\dot{3}$, and $0.\overline{142857}$. (The horizontal line is known as *vinculum*.)



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- Note: The decimal representation is not unique: we have 1.0 = 0.9 = 0.9999...where the ellipsis "..." represents an infinite sequence of the digit 9.
- In fact, every non-zero, finitely represented decimal number has an alternate. representation with trailing 9s, such as 123.4567 as 123.45669.

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Lemma 140

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Decimal Notation: Is It a Rational Number?

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Definition 141 (Irrational)

A number $x \in \mathbb{R} \setminus \mathbb{Q}$ is called *irrational*.

Decimal Notation

Definition 142 (Decimal separator)

The decimal separator is a symbol which is used to mark the boundary between the integer part and the fractional part of a number in decimal representation.



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Warning

A least two symbols are in wide-spread use for the decimal separator!

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- Dots or commas are frequently used to group three digits into groups within the integer part. However, this practice is discouraged by ISO!



• By definition, (\mathbb{N}, \leqslant) is well-ordered. And we have already hinted at well-orderings for \mathbb{Z} and \mathbb{Q} .



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Well-Order "Theorem"

Every set can be well-ordered.

 In 1883, Georg Cantor stated that the Well-Order Theorem is a "fundamental law of thought". This statement started a mathematical flame war!



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Well-Order "Theorem"

Every set can be well-ordered.

- In 1883, Georg Cantor stated that the Well-Order Theorem is a "fundamental law of thought". This statement started a mathematical flame war!
- In any case, this "theorem" can only be taken as an axiom, since it has been
 proved that it does not follow from any of the other commonly accepted axioms of
 set theory.
- In first-order logic, the Well-Order Theorem is equivalent to the Axiom of Choice (Dt.: Auswahlaxiom) and to Zorn's Lemma, in the sense that either one of them together with the Zermelo-Fraenkel Axioms allows to deduce the other ones

Theorem 143

The real numbers are uncountable, i.e., there exists no bijection from $\mathbb N$ onto $\mathbb R$.



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$$a_1 = -.d_1 - - - ...$$

 $a_2 = -.. - d_2 - - ...$
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If $d_k = 1$ then $r_k := 2$ else $r_k := 1$. Now regard r_k as the k-th digit of a number $r \in \mathbb{R}$: we have $r = 0.r_1r_2r_3r_4...$ Since at least the k-th digit of r differs from the k-th digit of a_k , we conclude that $r \neq a_n$ for all $n \in \mathbb{N}$.



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• Hence, $|\mathbb{N}| < |\mathbb{R}|$.



Theorem 144

For every $x \in \mathbb{R}$, every arbitrarily small neighborhood of x contains a rational number.

Sketch of proof: Let $x \in \mathbb{R}$ and $\varepsilon \in \mathbb{R}^+$ be arbitrary but fixed. W.I.o.g, 0 < x < 1. Let $k \in \mathbb{N}$ such that $10^{-k} < \varepsilon$.



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We define the rational number P/q as follows:

$$p := \left| x \cdot 10^k \right| \qquad q := 10^k$$



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This gives

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• E.g., $\pi \approx 3.1415 = \frac{31415}{10000} \text{ with } |\pi - \frac{31415}{10000}| \leqslant \frac{1}{10000}.$



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This gives

$$\left|x-\frac{p}{q}\right|\leqslant \frac{1}{10^k}<\varepsilon.$$

- E.g., $\pi \approx 3.1415 = \frac{31415}{10000}$ with $|\pi \frac{31415}{10000}| \leqslant \frac{1}{10000}$.
- Thus, we can approximate a real number by a rational number P/q.
- If the denominator q is a power of 10 then we can guarantee the error to be at most 1/q. Otherwise, if we allow an arbitrary integer q as denominator, we can guarantee the error to be at most $1/q^2$.

Theorem 145

No (non-empty) set A has the same cardinality as its power set $\mathcal{P}(A)$.



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- \bullet This implies that $|\mathbb{N}|<|\mathcal{P}(\mathbb{N})|<|\mathcal{P}(\mathcal{P}(\mathbb{N}))|<\dots$
- The cardinality of $\mathbb N$ is denoted by \aleph_0 .

Theorem 146

$$|\mathbb{R}|=|\mathcal{P}(\mathbb{N})|=2^{\aleph_0}>\aleph_0=|\mathbb{N}|=|\mathbb{Z}|=|\mathbb{Q}|.$$



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- The continuum hypothesis started out as a conjecture, until it was shown to be consistent with the usual axioms of the reals (by Gödel in 1940), and independent of those axioms (by Cohen in 1963).



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- The continuum hypothesis started out as a conjecture, until it was shown to be consistent with the usual axioms of the reals (by Gödel in 1940), and independent of those axioms (by Cohen in 1963).
- Under this hypothesis, the cardinality of \mathbb{R} equals \aleph_1 , and we have $2^{\aleph_0} = \aleph_1$.
- Furthermore, $|\mathcal{P}(\mathbb{R})| =: \aleph_2$, etc.



4

Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers
- Rational Numbers
- Real Numbers
- More Proof Techniques
 - Pigeonhole Principle
 - Well-founded Induction
 - Structural Induction



The Pigeonhole Principle

- In 1834, Johann Dirichlet noted that if there are five objects in four drawers then there is a drawer with two or more objects.
- Pigeonhole Principle: If *n* letters are posted to *m* pigeonholes, then
 - at least one pigeonhole receives more than one letter if n > m.
 - at least one pigeonhole remains empty if n < m.
 - each pigeonhole might receive exactly one letter if n = m.



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Theorem 147 (Pigeonhole Principle, Dt.: Schubfachschluss)

Consider two finite sets A and B. If A has more elements then B then every mapping from A to B will cause at least one element of B to be the target of two or more elements of A.



Lemma 148

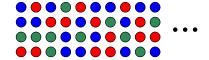
Consider a rectangular grid of points which consists of four rows and 100 columns.

```
0000000000
0000000000
0000000000
```



Lemma 148

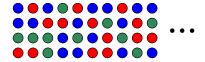
Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue.





Lemma 148

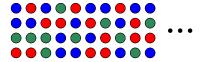
Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue. Prove that there always exist four points of the same color that form the corners of a rectangle (with sides parallel to the grid), no matter how the coloring is done.





Lemma 148

Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue. Prove that there always exist four points of the same color that form the corners of a rectangle (with sides parallel to the grid), no matter how the coloring is done.



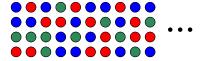
Proof: A column pattern is the top-to-bottom sequence of colors assigned to the four points of a column of the grid.





Lemma 148

Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue. Prove that there always exist four points of the same color that form the corners of a rectangle (with sides parallel to the grid), no matter how the coloring is done.

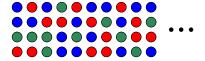


Proof: A column pattern is the top-to-bottom sequence of colors assigned to the four points of a column of the grid. There are exactly $3^4 = 81$ different column patterns.



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Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue. Prove that there always exist four points of the same color that form the corners of a rectangle (with sides parallel to the grid), no matter how the coloring is done.

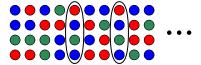


Proof: A column pattern is the top-to-bottom sequence of colors assigned to the four points of a column of the grid. There are exactly $3^4 = 81$ different column patterns. Since there are more than 81 columns, we are guaranteed to have at least two columns with the same column pattern.



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Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue. Prove that there always exist four points of the same color that form the corners of a rectangle (with sides parallel to the grid), no matter how the coloring is done.

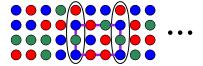


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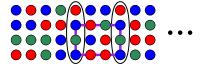


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Lemma 148

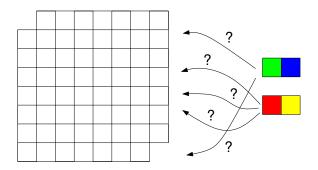
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Proof: A column pattern is the top-to-bottom sequence of colors assigned to the four points of a column of the grid. There are exactly 3⁴ = 81 different column patterns. Since there are more than 81 columns, we are guaranteed to have at least two columns with the same column pattern. Consider two such columns. Since there are four rows but only three colors, we conclude that two of the rows have the same color, thus giving us the four corners of the rectangle sought. □

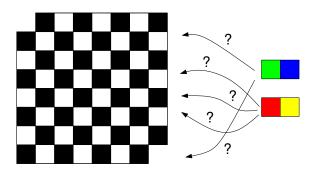
• Note: Just 19 columns suffice to guarantee the existance of such a rectangle

 Question: Can our modified chessboard be covered completely by 31 domino blocks of arbitrary color combinations?





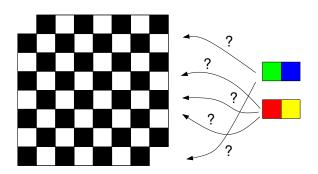
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 We observe that every permissible domino placement covers exactly one black square and one white square of the chessboard.



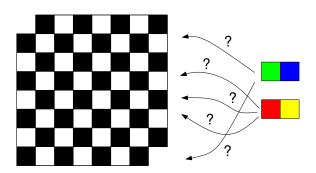
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 Question: Can our modified chessboard be covered completely by 31 domino blocks of arbitrary color combinations?



- We observe that every permissible domino placement covers exactly one black square and one white square of the chessboard.
- Thus, all domino placements would establish a one-to-one mapping between black and white squares. However, there are 32 black squares and only 30 white squares! We conclude that our chessboard cannot be covered completely domino blocks.

 Could one design an algorithm for lossless data compression that is guaranteed not to increase the file size of some input file while achieving a genuine compression for at least one other file?



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 Suppose further that there exists a compression algorithm that transforms every file into a distinct file which is no longer than the original file, and that at least one file will be compressed into something that is shorter than itself.



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- Since n < m, every file of length n keeps its size during compression. There are 2^n many such files. Together with f we would have $2^n + 1$ files which all compress into one of the 2^n files of length n.



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- By the pigeonhole principle there must exist some file f' of length n which is the
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- Hence, every compression algorithm will increase the size of at least some file or keep the sizes of all files unchanged.

Definition 149 (Well-founded order, Dt.: wohlfundierte Ordnung)

A strict partial order < on M is called *well-founded* if every $X \subseteq M$, with $X \neq \emptyset$, has at least one minimal element relative to <. A poset (M, <) is called a well-founded poset if < is well-founded.



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- Some authors call a well-founded order also a Noetherian order, named after Emmy Noether (1882-1935).
- Not to be confused with a well-order (Dt.: Wohlordnung).



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Lemma 150

The poset (M, <) is well-founded if and only if no infinite strictly decreasing sequence in M exists, i.e., if an $a : \mathbb{N} \to M$ with $a_{i+1} < a_i$ for all $i \in \mathbb{N}$ does not exist.



Lexicographical Order

Definition 151

Let $(M_1, <_1)$ and $(M_2, <_2)$ be two posets. The *lexicographical ordering* $(<_1, <_2)_{lex}$ on $M_1 \times M_2$ is defined as

$$(a_1,b_1) \ (\prec_1, \prec_2)_{\text{lex}} \ (a_2,b_2) \quad :\Leftrightarrow \quad \big((a_1 \prec_1 a_2) \lor ((a_1 = a_2) \land (b_1 \prec_2 b_2))\big),$$

where $(a_1, b_1), (a_2, b_2) \in M_1 \times M_2$.



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Lemma 152

Let (M_1, \prec_1) and (M_2, \prec_2) be two posets. Then $M_1 \times M_2$ together with the lexicographical order $(\prec_1, \prec_2)_{lex}$ is a poset.

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• Similarly for a non-strict partial order \leq .

Lemma 153

The posets $(M_1, <_1)$ and $(M_2, <_2)$ are well-founded if and only if $(M_1 \times M_2, (<_1, <_2)_{lex})$ is well-founded.



Induction Revisited

• Consider a predicate P over $\mathbb N$ and recall the Strong Induction Principle (Thm 79): If P(1) and if

$$\forall k \in \mathbb{N} \ \left[\left(\forall (m \in \mathbb{N}, m \leqslant k) \ P(m) \right) \ \Rightarrow \ P(k+1) \right]$$

then

 $\forall n \in \mathbb{N} \ P(n).$



Induction Revisited

And yet another version with "implicit" base:
 If

$$\forall k \in \mathbb{N} \left[\left(\forall (m \in \mathbb{N}, m < k) \ P(m) \right) \Rightarrow P(k) \right]$$

then

$$\forall n \in \mathbb{N} \ P(n).$$



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Induction Revisited

And yet another version with "implicit" base:

$$\forall k \in \mathbb{N} \left[\left(\forall (m \in \mathbb{N}, m < k) \ P(m) \right) \Rightarrow P(k) \right]$$

then

$$\forall n \in \mathbb{N} \ P(n).$$

• Note: The base case was not lost! Rather, it is included since we have to prove P(1) using the "helpful knowledge" that P(m) holds for all $m \in \mathbb{N}$ with m < 1.



Well-founded Induction

Theorem 154 (Principle of Well-founded Induction, Dt.: wohlfundierte Induktion)

Let (M, <) be well-founded and P be a predicate on M.

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• That is, as inductive step we have to prove that the predicate holds for *k* if it holds for all predecessors *m* of *k* relative to <.



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 That is, as inductive step we have to prove that the predicate holds for k if it holds for all predecessors m of k relative to <.

Proof: Let $X := \{m \in M : P(m) \text{ is false}\}$, and suppose $X \neq \emptyset$. Since (M, <) is well-founded, X has a minimal element n. Thus, $\forall m \in M$ with m < n the predicate P(m) holds. The inductive step

$$(\forall (m \in M, m < n) \ P(m)) \ \Rightarrow \ P(n)$$

yields that P(n) holds, in contradiction to $n \in X$.



 We give a proof of the existance claim made by the Fundamental Theorem of Arithmetic (Thm. 102): Every natural number n > 1 is either a prime number or has a prime factorization.



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 m_1 is genuine divisor of k and m_2 is genuine divisor of k.

Hence, both m_1 and m_2 are predecessors of k.



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Hence, both m_1 and m_2 are predecessors of k. By the inductive hypothesis, we know that m_1 is either prime or has a prime factorization; same for m_2 . Thus, also k has a prime factorization, which establishes the inductive step.



Partial Order on Recursive Structures

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- E.g., for $a \in \Sigma$ and $\sigma, \sigma' \in \Sigma^*$, if $\sigma = a\sigma'$ then we could regard σ' to be "smaller" than σ .
- More generally, $\sigma' <_{\Sigma} \sigma$ if and only if σ can be obtained from σ' and other words over Σ by applying constructors finitely often. (Hence, in this case σ' is a *sub-string* of σ .)
- Easy to prove: $<_{\Sigma}$ is a well-founded partial order on Σ^* .



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P(s) for every instance $s \in S$ under the assumption $P(s_1), P(s_2), \ldots, P(s_k)$, for some suitable $k \in \mathbb{N}$, if s can be obtained in one recursive construction step from $s_1, s_2, \ldots, s_k \in S$,



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Structural induction can be seen as a special case of a well-founded induction.



Lemma 156

Let Σ be a finite set. For every $\sigma \in \Sigma^*$ we have $\sigma \bullet \epsilon = \epsilon \bullet \sigma = \sigma$.

Proof: Def. 37 immediately gives $\epsilon \bullet \sigma = \sigma$ for all $\sigma \in \Sigma^*$.



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The empty word ϵ is the only minimal element stated in the base case of the definition of Σ^* , and we have

$$\epsilon \bullet \epsilon \stackrel{\textit{Def.}}{=} \stackrel{\textit{37}}{\epsilon} \epsilon.$$



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Theorem 157 (Functional completeness of NAND)

The Nand junctor, \(\frac{1}{2}\), is (truth-functionally) complete.

- Thus, every formula of propositional logic has a logically equivalent formula that uses only Nand junctors.
- Hence, any digital circuit can be realized by using only one type of gate: NAND gates. (This is also true for the Non inverter.)



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Lemma 158

Let p, q denote two Boolean variables. The following logical equivalences hold:

$$\neg p \equiv (p \uparrow p) \qquad (p \land q) \equiv ((p \uparrow q) \uparrow (p \uparrow q)) \qquad (p \lor q) \equiv ((p \uparrow p) \uparrow (q \uparrow q))$$
$$(p \Rightarrow q) \equiv (\neg p \lor q) \qquad (p \Leftrightarrow q) \equiv ((p \Rightarrow q) \land (q \Rightarrow p))$$
$$\top \equiv (p \uparrow (p \uparrow p)) \qquad \bot \equiv (\top \uparrow \top)$$



Proof of Thm. 157: Recall Def. 2: Propositional formulas (over some fixed set of n propositional variables p_1, p_2, \ldots, p_n) follow a rigid recursive construction scheme.



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• The minimal elements of the base case are given by the variables p₁, p₂,...,p_n and the constants ⊥ and ⊤. Lem. 158 tells us that ⊥ and ⊤ can be expressed using Nand junctors.



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- ② Consider an arbitrary but fixed propositional formula ϕ_0 that contains at least one junctor.



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- ② Consider an arbitrary but fixed propositional formula ϕ_0 that contains at least one junctor. By the construction scheme of propositional formulas, the formula ϕ_0 is of the form $(\neg \phi_1)$ or $(\phi_1 \# \phi_2)$, for suitable propositional formulas ϕ_1, ϕ_2 and where # is one of the junctors $\land, \lor, \Leftrightarrow, \Rightarrow$.



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Assume as inductive hypothesis that ϕ_1,ϕ_2 can be expressed using only Nand junctors (or are simply variables).



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By using the scheme outlined in Lem. 158, also ϕ_0 can be expressed using only NAND junctors.



- **5** Principles of Elementary Counting and Combinatorics
 - Sum and Product Rule
 - Inclusion-Exclusion Principle
 - Binomial Coefficient
 - Permutations
 - Ordered Selection (Variation)
 - Unordered Selection (Combination)



5

Principles of Elementary Counting and Combinatorics

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Theorem 159 (Sum rule, Dt.: Additionsprinzip)

Let A, B be two finite sets with $A \cap B = \emptyset$. Then

$$|A \cup B| = |A| + |B|.$$



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Corollary 160

For $n \in \mathbb{N}$, let A_1, A_2, \dots, A_n be n finite sets that are pairwise disjoint. Then

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Theorem 161 (Product rule, Dt.: Multiplikationsprinzip)

Let A, B be two finite sets. Then

$$|A \times B| = |A| \cdot |B|.$$



Proof of Theorem 161:

We observe that

$$A\times B=\bigcup_{b\in B}(A\times\{b\}),\quad \text{ with } (A\times\{b_1\})\cap (A\times\{b_2\})=\emptyset \ \text{ if } \ b_1\neq b_2.$$



Proof of Theorem 161:

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• There exists a bijective mapping between A and $A \times \{b\}$ for every $b \in B$. Thus, $|A| = |A \times \{b\}|$, and the theorem is a consequence of Corollary 160.



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Corollary 162

For $n \in \mathbb{N}$, let A_1, A_2, \dots, A_n be n finite sets. Then

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Corollary 163

For a propositional formula that contains n variables, 2^n evaluations are necessary in order to test all possible combinations of truth assignments to its variables.



Characteristic Function and Cardinality of Power Set

Definition 164 (Characteristic function, Dt.: Indikatorfunktion)

Let A be a finite set, and $B \subseteq A$. The *characteristic function* $1_B : A \to \{0, 1\}$ indicates membership of an element of A in B:

$$1_B(a) := \left\{ \begin{array}{ll} 1 & \text{if} \quad a \in B, \\ 0 & \text{if} \quad a \notin B. \end{array} \right.$$



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Lemma 165

A finite set A has $2^{|A|}$ many different subsets. That is, $|\mathcal{P}(A)| = 2^{|A|}$.



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Proof: We observe that every subset of A, including \emptyset and A itself, has a one-to-one correspondance to a characteristic function. Thus, every subset of A corresponds to a sequence of n 0's and 1's, where n := |A|. We conclude that the power set $\mathcal{P}(A)$ has 2^n members.



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Lemma 166

Let A be a finite set, and $B \subseteq A$. Then $|B| = \sum_{a \in A} 1_B(a)$.



• How many 3-element strings *s* can be formed over the standard Latin alphabet — 26 lower-case letters — such that every string contains at least one *x*?



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 26 lower-case letters such that every string contains at least one x?
- Obviously such a 3-element string *s* is in exactly one of the following sets:

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A_1 := \{s : \text{ first } x \text{ in first place of } s\},

A_2 := \{s : \text{ first } x \text{ in second place of } s\},

A_3 := \{s : \text{ first } x \text{ in third place of } s\}.
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Applying the Product Rule 161 yields

$$|A_1| = |\{x\} \times \{a, b, \dots, z\} \times \{a, b, \dots, z\}| = 1 \cdot 26 \cdot 26,$$



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 26 lower-case letters such that every string contains at least one x?
- Obviously such a 3-element string s is in exactly one of the following sets:

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A_1 := \{s : \text{ first } x \text{ in first place of } s\},

A_2 := \{s : \text{ first } x \text{ in second place of } s\},

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Applying the Product Rule 161 yields

$$|A_1| = |\{x\} \times \{a, b, \dots, z\} \times \{a, b, \dots, z\}| = 1 \cdot 26 \cdot 26,$$

 $|A_2| = |(\{a, b, \dots, z\} \setminus \{x\}) \times \{x\} \times \{a, b, \dots, z\}| = 25 \cdot 1 \cdot 26,$



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• Since A_1, A_2, A_3 are pairwise disjoint, the Sum Rule 159 implies

$$|A_1 \cup A_2 \cup A_3| = |A_1| + |A_2| + |A_3| = 26 \cdot 26 + 25 \cdot 26 + 25 \cdot 25 = 1951.$$



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- How many different passwords do exist under these restrictions?



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- By the Product Rule 161, the total number of six-character strings (over the 26 letters and the 10 digits) is 36⁶, with 26⁶ of them containing no digit at all. Hence,

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Hence, by the Sum Rule 159,

$$N = N_6 + N_7 + N_8 = 2684483063360.$$



5

Principles of Elementary Counting and Combinatorics

- Sum and Product Rule
- Inclusion-Exclusion Principle
- Binomial Coefficient
- Permutations
- Ordered Selection (Variation)
- Unordered Selection (Combination)



Theorem 167 (Inclusion-exclusion principle, Dt.: Siebprinzip, Poincaré-Formel)

Let A_1, A_2, \ldots, A_n be finite sets. Then

$$|\bigcup_{i=1}^{n} A_i| = \sum_{\substack{l \neq \emptyset \\ l \subseteq \{1, \dots, n\}}} (-1)^{|l|+1} |\bigcap_{i \in I} A_i|.$$



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- For |*I*| = 1:
 - $\sum_{1 \le i \le n} (-1)^{1+1} |A_i| = \sum_{i=1}^n |A_i|.$
- For |*I*| = 2:

$$\sum_{1\leqslant i< j\leqslant n} (-1)^{2+1} |A_i\cap A_j| = -\sum_{1\leqslant i< j\leqslant n} |A_i\cap A_j|.$$



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In particular:

$$|A_1 \cup A_2| = |A_1| + |A_2| - |A_1 \cap A_2|$$

$$|A_1 \cup A_2 \cup A_3| = |A_1| + |A_2| + |A_3| - |A_1 \cap A_2| - |A_1 \cap A_3| - |A_2 \cap A_3| + |A_1 \cap A_2 \cap A_3|$$



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 How many bit strings of length eight either start with 1 as first bit or end in 00 as the two last bits? (This is a non-exclusive or!)



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 and $|A_2| = 2^6 = 64$ and $|A_1 \cap A_2| = 2^5 = 32$.



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 and $|A_2| = 2^6 = 64$ and $|A_1 \cap A_2| = 2^5 = 32$.

Hence, by the Inclusion-Exclusion Principle (Thm. 167),

$$|A_1 \cup A_2| = |A_1| + |A_2| - |A_1 \cap A_2| = 128 + 64 - 32 = 160.$$



5

Principles of Elementary Counting and Combinatorics

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Definition 168 (Binomial coefficient, Dt.: Binomialkoeffizient)

Let $n \in \mathbb{N}_0$ and $k \in \mathbb{Z}$. The binomial coefficient $\binom{n}{k}$ of n and k is defined as follows:

$$\begin{pmatrix} n \\ k \end{pmatrix} := \begin{cases} 0 & \text{if } k < 0, \\ \frac{n!}{k! \cdot (n-k)!} & \text{if } 0 \le k \le n, \\ 0 & \text{if } k > n. \end{cases}$$



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- Recall k! := 1 for k := 0.
- The binomial coefficient $\binom{n}{k}$ is pronounced as "n choose k"; Dt.: "n über k".



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Lemma 169

Let $n \in \mathbb{N}_0$ and $k \in \mathbb{Z}$.

$$\binom{n}{0} = \binom{n}{n} = 1 \qquad \qquad \binom{n}{1} = \binom{n}{n-1} = n \qquad \qquad \binom{n}{k} = \binom{n}{n-k}$$

• The following table contains the non-zero values of $\binom{n}{k}$ for $0 \le n, k \le 6$.

				k			
n	0	1	2	3	4	5	6
0	1						
1	1	1					
2	1	2	1				
3	1	3	3	1			
4	1	4	6	4	1		
5	1	5	10	10	5	1	
6	1	6	15	20	15	6	1

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5	1	5	10	10	5	1	
6	1	6	15	20	15	6	1

- Trivial to observe:
 - Each row begins and ends with 1.
 - Initially each row contains increasing numbers till its middle but then the numbers start to decrease.
 - Each row's first half is exactly the mirror image of its second half.



Binomial Coefficients: Pascal's Triangle

A simple rearrangement of the previous table yields what is known as Pascal's Triangle in the Western world (Blaise Pascal, 1623–1662). But it was already studied in India in the 10th century, and discussed by Omar Khayyam (1048–1131)!

						1						
					1		1					
				1		2		1				
			1		3		3		1			
		1		4		6		4		1		
	1		5		10		10		5		1	
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	1		5		10		10		5		1	
1		6		15		20		15		6		1

 All entries in this triangle, except for the left-most and right-most entries per row, are the sum of the two entries above them in the previous row.



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Theorem 170 (Khayyam, Yang Hui, Tartaglia, Pascal)

For $n \in \mathbb{N}$ and $k \in \mathbb{Z}$,

$$\binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k}.$$

Binomial Theorem

• We know: $(a + b)^2 = a^2 + 2ab + b^2$ and $(a + b)^3 = a^3 + 3a^2b + 3ab^2 + b^3$.



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Theorem 171 (Binomial Theorem, Dt.: Binomischer Lehrsatz)

For all $n \in \mathbb{N}_0$ and $a, b \in \mathbb{R}$,

$$(a+b)^n = \binom{n}{0}a^n + \binom{n}{1}a^{n-1}b + \cdots + \binom{n}{n}b^n$$

or, equivalently,

$$(a+b)^n = \sum_{i=0}^n \binom{n}{i} a^{n-i} b^i.$$



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or, equivalently,

$$(a+b)^n = \sum_{i=0}^n \binom{n}{i} a^{n-i} b^i.$$

Corollary 172

For all $n \in \mathbb{N}$ and all $x \in \mathbb{R}$:

$$\sum_{i=0}^{n} \binom{n}{i} x^{i} = (1+x)^{n}$$

$$\sum_{i=0}^{n} \binom{n}{i} = 2^{n}$$

$$\sum_{i=0}^{n} (-1)^{i} \binom{n}{i} = 0$$

5

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Permutations

Definition 173 (Permutation)

Let A be a finite set. A *permutation* of A is a bijective function from A to A.



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- Many encryption schemes used in cryptography can be seen as permutations.



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- Many encryption schemes used in cryptography can be seen as permutations.
- Standard notation for a permutation π of $\{1, 2, ..., n\}$:

$$\left(\begin{array}{cccc} 1 & 2 & 3 & \dots & n \\ \pi(1) & \pi(2) & \pi(3) & \dots & \pi(n) \end{array}\right)$$



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$$\begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 3 & 1 & 4 \end{pmatrix}$$



Definition 173 (Permutation)

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- Many encryption schemes used in cryptography can be seen as permutations.
- Standard notation for a permutation π of $\{1, 2, ..., n\}$: • E.g., for *n* := 4:

$$\left(\begin{array}{cccc} 1 & 2 & 3 & \dots & n \\ \pi(1) & \pi(2) & \pi(3) & \dots & \pi(n) \end{array}\right)$$

$$\left(\begin{array}{ccccc}
1 & 2 & 3 & 4 \\
2 & 3 & 1 & 4
\end{array}\right)$$

Definition 174 (Product of permutations)

Let A be a finite set together with two permutations α, β . Then the *product* (or *composition*) $\alpha \circ \beta$ is the function

$$\alpha \circ \beta : A \to A$$
 with $(\alpha \circ \beta)(a) := \alpha(\beta(a))$ for all $a \in A$.



- The product of two permutations is itself a bijective function, i.e., a permutation.
- Note: It is common to drop \circ in $\alpha \circ \beta$ and simply write $\alpha\beta$.



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- The product of two permutations is not commutative.

$$\alpha := \begin{pmatrix} 1 & 2 & 3 & 4 \\ 1 & 3 & 4 & 2 \end{pmatrix} \qquad \beta := \begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 1 & 3 & 4 \end{pmatrix}$$
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Lemma 175

For all $n \in \mathbb{N}$ and all finite sets A with n = |A|, the set of all permutations, S_n , over A together with \circ as operation forms a group, the so-called *symmetric group*.



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Lemma 176

For all $n \in \mathbb{N}$ and all finite sets A with n = |A|, the group (S_n, \circ) is a finite group with exactly n! members.

Definition 177 (Cycle, Dt.: Zyklus)

Let A be a finite set of cardinality n. A permutation π of A is a *cycle of length* $k \le n$ if there exists a set $B \subseteq A$ with |B| = k such that, with $B := \{b_1, b_2, \dots, b_k\}$,

$$\pi(b_1) = b_2, \ \pi(b_2) = b_3, \ \dots, \ \pi(b_{k-1}) = b_k, \ \pi(b_k) = b_1,$$

and $\pi(a) = a$ for all $a \in A \backslash B$. In this case this k-cycle is written as

$$(b_1 \ b_2 \ \dots \ b_k)$$
 or as $b_1 \mapsto b_2 \mapsto \dots \mapsto b_k \mapsto b_1$.

A cycle is *non-trivial* if $k \ge 2$.



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Definition 178 (Transposition)

A transposition is a cycle of length two, aka 2-cycle.



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Lemma 179

Every permutation (of two or more elements) can be written as

- (1) a product of cycles,
- (2) a product of transpositions.

Lemma 180

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The *signature* of a permutation is +1, and the permutation is *even*, if it consists of an even number of transpositions. Otherwise, the signature is -1 and the permutation is *odd*.



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200

Definition 183 (Inversion, Dt.: Inversion, Fehlstand)

A permutation $\pi \in S_n$ has an *inversion* (i,j) if $\pi(i) > \pi(j)$ for $1 \le i < j \le n$.

5

Principles of Elementary Counting and Combinatorics

- Sum and Product Rule
- Inclusion-Exclusion Principle
- Binomial Coefficient
- Permutations
- Ordered Selection (Variation)
- Unordered Selection (Combination)



Definition 184 (Ordered selection without repetition, Dt.: Variation ohne Zurücklegen, Variation ohne Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \le n$, and A be a finite set of cardinality n. An ordered selection without repetition of k elements from A is a k-tuple

$$(a_1, a_2, \dots, a_k)$$
 with $a_i \in A$ for $i = 1, 2, \dots, k$ and $a_i \neq a_j$ for $1 \leq i < j \leq k$.



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Lemma 185

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \leq n$, and A be a finite set of cardinality n. There exist

$$V_k^n := \frac{n!}{(n-k)!}$$

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- Convention: $V_k^n := 0$ for k > n.
- V_k^n is the number of injective functions from I_k to A.
- Sometimes, V(n,k) is written instead of V_k^n . English-language textbooks often speak of a k-permutation rather than of an ordered selection without repetition.

Definition 186 (Ordered selection with repetition, Dt.: Variation mit Zurücklegen, Variation mit Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. An ordered selection with repetition of k elements from A is a k-tuple

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- Note: ${}^{r}V_{k}^{n} = |A^{k}|$.
- Sometimes, $V_r(n,k)$ is written instead of ${}^rV_k^n$.



- 5
 - Principles of Elementary Counting and Combinatorics
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Definition 188 (Unordered selection without repetition, Dt.: Kombination ohne Zurücklegen, Kombination ohne Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \le n$, and A be a finite set of cardinality n. An unordered selection without repetition of k elements from A is a set B such that

$$B \subseteq A$$
 with $|B| = k$.



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Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \le n$, and A be a finite set of cardinality n. There exist

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Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \le n$, and A be a finite set of cardinality n. An unordered selection without repetition of k elements from A is a set B such that

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many different unordered selections without repetition of k elements from A.

- Convention: $C_k^n := 0$ for k > n. Sometimes, C(n,k) is written instead of C_k^n .
- Lemma 189 yields an alternate proof of $|\mathcal{P}(A)| = 2^n$. It also implies that there exist $\binom{n}{\nu}$ different binary sequences where exactly k elements are 1.

Definition 190 (Unordered selection with repetition, Dt.: Kombination mit Zurücklegen, Kombination mit Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. An unordered selection with repetition of k elements from A is a k-element multiset, i.e., a set $B \subseteq A$ together with a multiplicity function, mult: $A \to \mathbb{N}_0$, such that

$$\operatorname{mult}(a) = 0 \text{ for all } a \in A \backslash B \text{ and } \operatorname{mult}(b) > 0 \text{ for all } b \in B \text{ and } \sum_{b \in B} \operatorname{mult}(b) = k.$$



Definition 190 (Unordered selection with repetition, Dt.: Kombination mit Zurücklegen, Kombination mit Wiederholung)

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Lemma 191

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. There exist

$${}^{r}C_{k}^{n}:=\binom{n+k-1}{k}$$

many different unordered selections with repetition of k elements from A.



Definition 190 (Unordered selection with repetition, Dt.: Kombination mit Zurücklegen, Kombination mit Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. An unordered selection with repetition of k elements from A is a k-element multiset, i.e., a set $B \subseteq A$ together with a multiplicity function, mult: $A \to \mathbb{N}_0$, such that

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many different unordered selections with repetition of *k* elements from *A*.

• Sometimes, $C_r(n,k)$ is written instead of ${}^rC_k^n$.



Proofs of Lemmas 185-191

Proof of Lemma 185: We have n options for a_1 , leaving n-1 options for a_2 , etc.

Thus, we have
$$n \cdot (n-1) \cdot \ldots \cdot (n-k+1) = \frac{n!}{(n-k)!}$$
 options.

Proof of Lemma 187: We have n options for every selection. Thus, we have n^k options in total.

Proof of Lemma 189: We know that $V_k^n = \frac{n!}{(n-k)!}$. There are k! many different ordered selections that correspond to the same unordered selection. Thus,

$$C_k^n = V_k^n/k! = \frac{n!}{(n-k)!k!} = \binom{n}{k}.$$

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Proof of Lemma 185: We have n options for a_1 , leaving n-1 options for a_2 , etc.

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Proof of Lemma 191: Let a_1, \ldots, a_n be the *n* elements of *A*, and $k \in \mathbb{N}_0$. We encode such an unordered selection with repetition of k elements from A as a sequence of length n + k - 1 of k crosses \times which are separated by n - 1 vertical bars |, where i crosses between the j-th vertical bar and the (j + 1)-st vertical bar, for $1 \le j \le n - 2$, indicate that element a_{i+1} was chosen with multiplicity i.



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Proof of Lemma 185: We have *n* options for a_1 , leaving n-1 options for a_2 , etc.

Thus, we have $n \cdot (n-1) \cdot \ldots \cdot (n-k+1) = \frac{n!}{(n-k)!}$ options.

 $= \frac{1}{(n-k)!} \text{ options.}$

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Proof of Lemma 191: Let a_1, \ldots, a_n be the n elements of A, and $k \in \mathbb{N}_0$. We encode such an unordered selection with repetition of k elements from A as a sequence of length n+k-1 of k crosses \times which are separated by n-1 vertical bars |, where i crosses between the j-th vertical bar and the (j+1)-st vertical bar, for $1 \le j \le n-2$, indicate that element a_{j+1} was chosen with multiplicity i. Similarly for the multiplicities of a_1 and a_n for crosses before the first and after the last vertical bar. We note that we have exactly

$$C_k^{n+k-1} = \binom{n+k-1}{k}$$

ways to choose the positions of the k crosses within this sequence.



Real-World Application: Elementary Probability

 What is the probability to win in the Austrian "6-aus-45" lottery after choosing one set of six numbers?



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 equally-likely outcomes as the number of favorable outcomes divided by the total
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- As usual, we define the probability of an event among (finitely many)
 equally-likely outcomes as the number of favorable outcomes divided by the total
 number of possible outcomes.
- Assuming that the lottery is fair and, thus, that all combinations are equally likely to win, we get

$$\frac{1}{C_6^{45}} = \frac{1}{\binom{45}{6}} = \frac{1}{8145060} \approx 1.22774 \cdot 10^{-7}$$

as probability for having all six numbers right.



 A standard deck of cards contains 52 cards grouped into four suits (Dt.: Farben) - diamonds (Dt.: Schelle, Karo), clubs (Dt.: Eichel, Kreuz), hearts (Dt.: Herz), and spades (Dt.: Laub, Pik) — with 13 cards in each suit (ace, 2, 3, 4, 5, 6, 7, 8, 9, 10, jack, queen, king).



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- What is the probability that all hearts appear in consecutive (but arbitrary) order after a decent shuffling of the deck?
- There are 52! different permutations of the 52 cards.
- There are 40! different permutations of the block of 13 hearts and the other 39 cards, and 13! many permutations of the 13 hearts within that block.



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- There are 52! different permutations of the 52 cards.
- There are 40! different permutations of the block of 13 hearts and the other 39 cards, and 13! many permutations of the 13 hearts within that block.
- Hence, the probability that all hearts are consecutive is given by

$$\frac{40! \cdot 13!}{52!} \approx 6.29908 \cdot 10^{-11}.$$



- 6
 - **Complexity Analysis and Recurrence Relations**
 - Growth Rates
 - Bachmann-Landau (Asymptotic) Notation
 - Recurrence Relations
 - Master Theorem



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- Determine the dominating term in the complexity function it gives the order of magnitude of the asymptotic behavior.

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$$\log n$$
, $\log^2 n$, \sqrt{n} , n , $n \log n$, $n \log^2 n$, $n^{\frac{7}{6}}$, n^2 , n^3 , ..., 2^n , 3^n , $2^{(2^n)}$, ...



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• Recall that $\log_{\alpha} n = \frac{1}{\log_2 \alpha} \log_2 n$.



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Thus, we have

$$c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)$$

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g grows at most as fast as $c_2 \cdot f$ f is an asymptotic upper bound on g

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g grows at most as fast as $c_2 \cdot f$ f is an asymptotic upper bound on g we'll say that $g \in O(f)$

g has same growth rate as f we'll say that $g \in \Theta(f)$



6

Complexity Analysis and Recurrence Relations

- Growth Rates
- Bachmann-Landau (Asymptotic) Notation
 - Bachmann-Landau Symbols
 - Limit of a Sequence
 - Basic Facts
 - Conditional Asymptotic Notation and Smoothness Rule
- Recurrence Relations
- Master Theorem



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Definition 192 (Big-O, Dt.: Groß-O)

Let $f: \mathbb{N} \to \mathbb{R}^+$. Then the set O(f) is defined as

$$O(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_2 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \forall n \geqslant n_0 \qquad g(n) \leqslant c_2 \cdot f(n) \}.$$



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$$O(f) \ := \ \left\{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \quad \exists c_2 \in \mathbb{R}^+ \quad \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad g(n) \leqslant c_2 \cdot f(n) \right\}.$$

• Some authors prefer to use the symbol \mathcal{O} instead of \mathcal{O} .



$$c_1\cdot f(n)\leq \dot{g}(n)\leq c_2\cdot f(n)$$
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Definition 192 (Big-O, Dt.: Groß-O)

Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set O(f) is defined as

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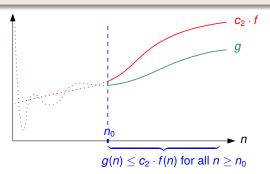
- Some authors prefer to use the symbol \mathcal{O} instead of O.
- Note: O(f) is a set of functions! Definitions of the form $O(f(n)) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_2 \in \mathbb{R}^+ \mid \exists n_0 \in \mathbb{N} \mid \forall n \geqslant n_0 \quad g(n) \leqslant c_2 \cdot f(n) \}$ are a (wide-spread) formal nonsense.

Graphical Illustration of O(f)

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Let $f: \mathbb{N} \to \mathbb{R}^+$. Then the set O(f) is defined as

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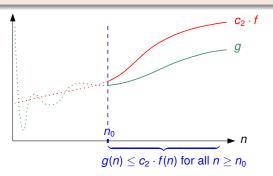


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Equivalent definition used by some authors:

$$O(f) \ := \ \left\{g \colon \mathbb{N} o \mathbb{R}^+ \mid \ \exists c_2 \in \mathbb{R}^+ \ \exists n_0 \in \mathbb{N} \ orall n_0 \geqslant n_0 \ rac{g(n)}{f(n)} \leqslant c_2
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Note that this notation hides all lower-order terms and multiplicative constants.
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- Consider the following two nested for-loops:

```
for i = 1 to n do
  for j = i to n do
    Compute(i, j)
  end for
end for
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- We get

$$g(n) = n + (n-1) + ... + 2 + 1$$

= $\frac{n(n+1)}{2} = \frac{1}{2}n^2 + \frac{1}{2}n$.



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= $\frac{n(n+1)}{2} = \frac{1}{2}n^2 + \frac{1}{2}n$.

- Consider $f: \mathbb{N} \to \mathbb{R}^+$ with $f(n) := n^2$.
- Let's compare the growth rates of f and g when we double n:

n	g(n)	<i>f</i> (<i>n</i>)
5	15	25
10	55	100
20	210	400
40	820	1600
80	3240	6400



- Note that this notation hides all lower-order terms and multiplicative constants.
 Why don't we care?
- Since it doesn't matter for large values of *n*.
- Consider the following two nested for-loops:

for
$$i = 1$$
 to n do
for $j = i$ to n do
Compute (i, j)
end for

- How often is Compute() being called?
 Let g: N→ R⁺ be the function that models the number of calls in dependence on n.
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$$g(n) = n + (n-1) + ... + 2 + 1$$

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- Consider $f: \mathbb{N} \to \mathbb{R}^+$ with $f(n) := n^2$.
- Let's compare the growth rates of *f* and *g* when we double *n*:

n	g (n)	<i>f</i> (<i>n</i>)
5	15	25
10	55	100
20	210	400
40	820	1600
80	3240	6400

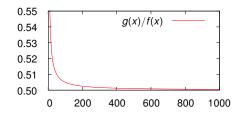
 Doubling n causes both f(n) and g(n) to (roughly) quadruple!



• We plot the growth ratio $\frac{g(n)}{f(n)}$ for $f,g:\mathbb{N}\to\mathbb{R}^+$ with $f(n):=n^2$ and $g(n):=\frac{1}{2}n^2+\frac{1}{2}n$.

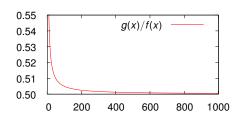


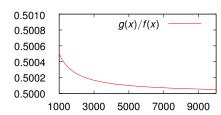
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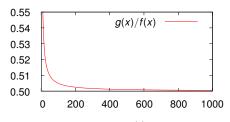
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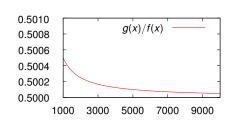






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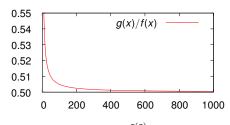


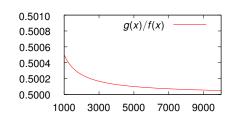


• The plots suggest $\frac{g(n)}{f(n)} \le 1$ for all $n \ge 200$, that is, $g(n) \le f(n)$, which would imply $g \in O(f)$.



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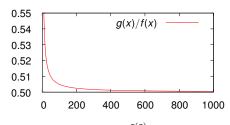


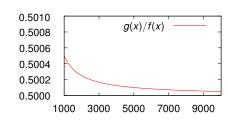


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- More precisely, they suggest $\frac{g(n)}{f(n)} \leq \frac{1}{2} + \varepsilon$ for any positive ε and all sufficiently large values of n.



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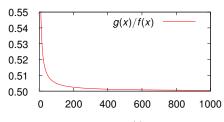


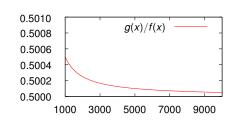


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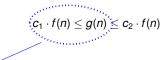




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- The plots also suggest $\frac{g(n)}{f(n)} \ge \frac{1}{2}$, which would imply $g \in \Omega(f)$.
- Hence $g(n) \approx \frac{1}{2}f(n)$, which would imply $g \in \Theta(f)$.



Asymptotic Notation: Big-Omega



 $\left\{\begin{array}{l} \text{ for all } n \geq n_0 \text{ and} \\ \text{ fixed } c_1, c_2 \in \mathbb{R}^+. \end{array}\right.$

g grows at least as fast as $c_1 \cdot f$ f is an asymptotic lower bound on g we'll say that $g \in \Omega(f)$



Asymptotic Notation: Big-Omega

$$c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)$$

 $\left\{ \begin{array}{l} \text{for all } n \geq n_0 \text{ and} \\ \text{fixed } c_1, c_2 \in \mathbb{R}^+. \end{array} \right.$

g grows at least as fast as $c_1 \cdot f$

f is an asymptotic lower bound on g

we'll say that $g \in \Omega(f)$

Definition 193 (Big-Omega, Dt.: Groß-Omega)

Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set $\Omega(f)$ is defined as

$$\Omega(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_1 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad c_1 \cdot f(n) \leqslant g(n) \}.$$



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g grows at least as fast as $c_1 \cdot f$

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Equivalently,

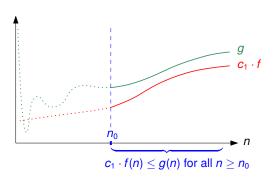
$$\Omega(f) \ := \ \left\{g \colon \mathbb{N} o \mathbb{R}^+ \mid \ \exists c_1 \in \mathbb{R}^+ \quad \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad c_1 \leqslant \dfrac{g(n)}{f(n)} \right\}$$

Graphical Illustration of $\Omega(f)$

Definition 193 (Big-Omega, Dt.: Groß-Omega)

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Asymptotic Notation: Big-Theta

$$c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)$$

for all $n \geq n_0$ and fixed $c_1, c_2 \in \mathbb{R}^+$.

g has same growth rate as f we'll say that $g \in \Theta(f)$



Asymptotic Notation: Big-Theta

$$\underbrace{c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)}_{\textit{g} \text{ has same growth rate as } f} \left\{ \begin{array}{l} \text{for all } n \geq n_0 \text{ and} \\ \text{fixed } c_1, c_2 \in \mathbb{R}^+. \end{array} \right.$$

Definition 194 (Big-Theta, Dt.: Groß-Theta)

Let $f: \mathbb{N} \to \mathbb{R}^+$. Then the set $\Theta(f)$ is defined as

$$\Theta(f) := \begin{cases} g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_1, c_2 \in \mathbb{R}^+ & \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \\ c_1 \cdot f(n) \leqslant g(n) \leqslant c_2 \cdot f(n) \end{cases}.$$





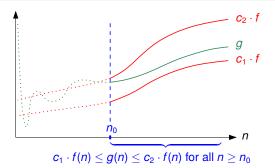
Graphical Illustration of $\Theta(f)$

Definition 194 (Big-Theta, Dt.: Groß-Theta)

Let $f: \mathbb{N} \to \mathbb{R}^+$. Then the set $\Theta(f)$ is defined as

$$\Theta(f) \ := \ \big\{g \colon \mathbb{N} \to \mathbb{R}^+ \, | \quad \exists c_1, \ c_2 \in \mathbb{R}^+ \quad \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0$$

$$c_1 \cdot f(n) \leqslant g(n) \leqslant c_2 \cdot f(n) \}.$$



which is equivalent to $c_1 \leq \frac{g(n)}{f(n)} \leq c_2$ for all $n \geq n_0$



• We prove $g \in \Theta(f)$ for $f(n) := n^2$ and $g(n) := \frac{1}{2}n^2 + \frac{1}{2}n$.



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Proof:

• We get, for all $n \in \mathbb{N}$,

$$g(n) = \frac{1}{2}n^2 + \frac{1}{2}n \leqslant \frac{1}{2}n^2 + \frac{1}{2}n^2 = n^2 = f(n),$$



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• Thus, $g \in O(f)$ with $c_2 := 1$ and $n_0 := 1$.



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- Thus, $g \in O(f)$ with $c_2 := 1$ and $n_0 := 1$.
- Now we prove $g \in \Omega(f)$ and get, again for all $n \in \mathbb{N}$,

$$g(n) = \frac{1}{2}n^2 + \frac{1}{2}n \geqslant \frac{1}{2}n^2 = \frac{1}{2}f(n),$$



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Proof:

• We get, for all $n \in \mathbb{N}$,

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- Thus, $g \in O(f)$ with $c_2 := 1$ and $n_0 := 1$.
- Now we prove $g \in \Omega(f)$ and get, again for all $n \in \mathbb{N}$,

$$g(n)=\frac{1}{2}n^2+\frac{1}{2}n\geqslant\frac{1}{2}n^2=\frac{1}{2}f(n),\quad \text{that is }\frac{1}{2}f(n)\leqslant g(n).$$

• Thus, $g \in \Omega(f)$ with $c_1 := \frac{1}{2}$ and $n_0 := 1$. Def. 194 or Lemma 201 yield $g \in \Theta(f)$.



• We prove $g \in \Theta(f)$ for $f(n) := n^2$ and $g(n) := \frac{1}{2}n^2 + \frac{1}{2}n$.

Proof:

• We get, for all $n \in \mathbb{N}$,

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- Thus, $g \in O(f)$ with $c_2 := 1$ and $n_0 := 1$.
- Now we prove $g \in \Omega(f)$ and get, again for all $n \in \mathbb{N}$,

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• Thus, $g \in \Omega(f)$ with $c_1 := \frac{1}{2}$ and $n_0 := 1$. Def. 194 or Lemma 201 yield $a \in \Theta(f)$.

Don't be overly zealous!

There is no need to try to obtain the "best-possible" values for n_0 and c_1, c_2 !

Definition 195 (Small-Oh, Dt.: Klein-O)

Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set o(f) is defined as

$$o(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \forall c \in \mathbb{R}^+ \mid \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad g(n) \leqslant c \cdot f(n) \}.$$

Definition 195 (Small-Oh, Dt.: Klein-O)

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Mind the difference

$$O(f) := \begin{cases} g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ & \exists n_0 \in \mathbb{N} & \forall n \geqslant n_0 \\ o(f) := \begin{cases} g \colon \mathbb{N} \to \mathbb{R}^+ \mid \forall c \in \mathbb{R}^+ & \exists n_0 \in \mathbb{N} & \forall n \geqslant n_0 \\ \end{cases} g(n) \leqslant c \cdot f(n) \end{cases}$$



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Mind the difference

$$O(f) := \begin{cases} g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ & \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad g(n) \leqslant c \cdot f(n) \end{cases}$$

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• Similarly, $\omega(f)$ can be defined relative to $\Omega(f)$.



Definition 195 (Small-Oh, Dt.: Klein-O)

Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set o(f) is defined as

$$o\left(f\right) \;:=\; \left\{g\colon \mathbb{N}\to\mathbb{R}^+ \;\middle| \quad \forall c\in\mathbb{R}^+ \quad \exists n_0\in\mathbb{N} \quad \forall n\geqslant n_0 \qquad g(n)\leqslant c\cdot f(n)\right\}.$$

Mind the difference

$$O(f) := \begin{cases} g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ & \exists n_0 \in \mathbb{N} & \forall n \geqslant n_0 \\ o(f) := \begin{cases} g \colon \mathbb{N} \to \mathbb{R}^+ \mid \forall c \in \mathbb{R}^+ & \exists n_0 \in \mathbb{N} & \forall n \geqslant n_0 \\ \end{cases} g(n) \leqslant c \cdot f(n) \end{cases}$$

- Similarly, $\omega(f)$ can be defined relative to $\Omega(f)$.
- It is trivial to extend Definitions 192–195 such that \mathbb{N}_0 rather than \mathbb{N} is taken as the domain.
- We can also replace the codomain \mathbb{R}^+ by \mathbb{R}_0^+ (or even \mathbb{R}) provided that all functions are eventually positive.
- The same comments apply to the subsequent slides.



Definition 196 (Sequence, Dt.: Folge)

A (real) sequence is a function from $\mathbb N$ (or $\mathbb N_0$) to $\mathbb R$. For $x \colon \mathbb N \to \mathbb R$ it is common to write the sequence as $(x_n)_{n \in \mathbb N}$ or $\langle x_n \rangle_{n \in \mathbb N}$, or simply (x_n) or $\langle x_n \rangle$.



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Definition 197 (Limit, Dt. Grenzwert)

The value $\bar{x} \in \mathbb{R}$ is the limit of the (real) sequence (x_n) , denoted by $\lim_{n\to\infty} x_n = \bar{x}$, if

$$\forall \varepsilon \in \mathbb{R}^+ \ \exists n_0 \in \mathbb{N} \ \forall n \geqslant n_0 \ |x_n - \bar{x}| < \varepsilon.$$



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Lemma 198

If $z_n = x_n + y_n$ for three sequences $(x_n), (y_n), (z_n)$ and if $\lim_{n \to \infty} x_n$ and $\lim_{n \to \infty} y_n$ exist, then $\lim_{n \to \infty} z_n$ exists and we have $\lim_{n \to \infty} z_n = \lim_{n \to \infty} x_n + \lim_{n \to \infty} y_n$.



Theorem 199 (Squeeze theorem, Dt.: Einschnürungssatz)

Consider three real sequences $(x_n), (y_n), (z_n)$ and suppose that $x_n \le y_n \le z_n$ for all $n \ge n_0$ for some $n_0 \in \mathbb{N}$. If the limits of (x_n) and (z_n) exist such that

$$\lim_{n\to\infty} x_n = \lim_{n\to\infty} z_n,$$

then the limit of (y_n) exists with

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• For $z_n := \frac{8}{n}$ it is easy to see that $\lim_{n\to\infty} z_n = 0$.



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- Now consider the following sequences:

$$x_n := 0$$
 $y_n := \frac{\log n + 7\sqrt{n} - 10}{n^2}$ $z_n := \frac{8}{n}$



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- Now consider the following sequences:

$$x_n := 0$$
 $y_n := \frac{\log n + 7\sqrt{n} - 10}{n^2}$ $z_n := \frac{8}{n}$

• We have for all $n \in \mathbb{N} \setminus \{1, 2, 3\}$

$$x_n \leqslant y_n \leqslant z_n$$
 and $\lim_{n \to \infty} x_n = 0 = \lim_{n \to \infty} z_n$.

Thus, $\lim_{n\to\infty} y_n = 0$.



 The following theorem (by Guillaume de l'Hôpital, 1661–1704) allows to handle limits that involve indeterminate terms of the form

$$\frac{0}{0}$$
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then

$$\lim_{x\to c}\frac{f(x)}{g(x)}=\lim_{x\to c}\frac{f'(x)}{g'(x)}.$$



Lemma 201



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$$O(c \cdot f_1) = O(f_1)$$

$$\Theta(f_1) = O(f_1) \cap \Omega(f_1)$$



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Lemma 201

Let $f_1, f_2, g_1, g_2 : \mathbb{N} \to \mathbb{R}^+$, and $c \in \mathbb{R}^+$. Then the following relations hold:

$$O(c \cdot f_1) = O(f_1)$$



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Lemma 202

Let $f,g \colon \mathbb{N} \to \mathbb{R}^+$ and $c \in \mathbb{R}^+$. Then:

$$\lim_{n\to\infty}\frac{g(n)}{f(n)}=c\quad \Rightarrow\quad g\in\Theta(f),$$



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• For example, let $f, g, h \colon \mathbb{N} \to \mathbb{R}^+$ with $f(n) := n^2 - 7n$, $g(n) := 3n^2 + 5n\sqrt{n}$ and $h(n) := n^2$.



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$$\lim_{n\to\infty}\frac{f(n)}{h(n)}=\lim_{n\to\infty}\frac{n^2-7n}{n^2}=\lim_{n\to\infty}\left(1-\frac{7}{n}\right)=1$$



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$$\lim_{n\to\infty}\frac{g(n)}{h(n)}=\lim_{n\to\infty}\frac{3n^2+5n\sqrt{n}}{n^2}=\lim_{n\to\infty}\left(3+\frac{5}{\sqrt{n}}\right)=3$$



Asymptotic Notation: Wide-spread Notational Abuse

It is convenient to be a bit sloppy and write, e.g.,

$$g(n) = O(n^2)$$
 or $g \in O(n^2)$

rather than to resort to the λ -quantifier and write $g \in O(\lambda n.n^2)$, or

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Similarly,

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means

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means

$$|g-h| \in O(f)$$
 with $f: \mathbb{N} \to \mathbb{R}^+, n \mapsto n^3$.

Furthermore,

$$g(n) = n^{O(1)}$$

indicates that

$$g \in O(f)$$
 with $f: \mathbb{N} \to \mathbb{R}^+, n \mapsto n^c$

for some constant $c \in \mathbb{R}^+$.



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- So, keep in mind that an is-element-of or subset relation is meant even if an
 equality sign is used!
- Unfortunately, several textbooks are fuzzy about this important distinction . . .



Definition 203 (Conditional Asymptotic Notation)

Consider a function $f: \mathbb{N} \to \mathbb{R}^+$

$$O(f \quad) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 :$$

$$g(n) \leqslant c \cdot f(n)$$
.



Definition 203 (Conditional Asymptotic Notation)

Consider a function $f: \mathbb{N} \to \mathbb{R}^+$ and a predicate $P: \mathbb{N} \to \{F, T\}$.

$$O(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 : \}$$

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.

• E.g., let
$$P(n) :\Leftrightarrow n \equiv_2 0$$
, or $P(n) :\Leftrightarrow (\exists k \in \mathbb{N}_0 \ n = 2^k)$.



Definition 203 (Conditional Asymptotic Notation)

Consider a function $f: \mathbb{N} \to \mathbb{R}^+$ and a predicate $P: \mathbb{N} \to \{F, T\}$.

$$O(f \mid P) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \forall n \geqslant n_0 : P(n) \Rightarrow g(n) \leqslant c \cdot f(n) \}.$$

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$$\Omega(f \mid P) := \left\{ g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 : \\ P(n) \Rightarrow g(n) \geqslant c \cdot f(n) \right\}.$$

$$\Theta(f \mid P) := \left\{ g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_1, c_2 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \forall n \geqslant n_0 : \\ P(n) \Rightarrow c_1 \cdot f(n) \leqslant g(n) \leqslant c_2 \cdot f(n) \right\}.$$

$$o(f \mid P) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \forall c \in \mathbb{R}^+ \mid \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 : P(n) \Rightarrow g(n) \leqslant c \cdot f(n) \}.$$

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200

Smoothness

Definition 204 (Eventually non-decreasing, Dt.: schlussendlich nicht abnehmend)

A function $f: \mathbb{N} \to \mathbb{R}^+$ is eventually non-decreasing exactly if

$$\exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad f(n) \leqslant f(n+1).$$



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A function $f: \mathbb{N} \to \mathbb{R}^+$ is *eventually non-decreasing* exactly if

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Definition 205 (b-smooth, Dt.: b-glatt)

A function $f: \mathbb{N} \to \mathbb{R}^+$ is *b-smooth* for some integer $b \ge 2$ exactly if f is eventually non-decreasing and if

$$\exists c \in \mathbb{R}^+ \quad \exists n_0 \in \mathbb{N} \quad \forall n \geqslant n_0 \qquad f(b \cdot n) \leqslant c \cdot f(n).$$



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Definition 206 (smooth, Dt.: glatt)

A function $f: \mathbb{N} \to \mathbb{R}^+$ is *smooth* if it is *b*-smooth for all integers $b \ge 2$.



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Lemma 207

If $f: \mathbb{N} \to \mathbb{R}^+$ is b'-smooth for some integer $b' \geqslant 2$ then it is smooth.

Smoothness Rule

Theorem 208 (Smoothness Rule)

Let $f, g: \mathbb{N} \to \mathbb{R}^+$, and consider an integer $b \ge 2$.



Smoothness Rule

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Let $f, g: \mathbb{N} \to \mathbb{R}^+$, and consider an integer $b \ge 2$. If

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then $g \in O(f)$.

- Similarly for $\Omega(f)$ and $\Theta(f)$.
- Again, it is trivial to extend the definitions and lemmas such that \mathbb{N}_0 rather than \mathbb{N} is taken as the base set. Similarly, we can replace \mathbb{R}^+ by \mathbb{R}_0^+ or even by \mathbb{R} provided that all functions are eventually positive.
- The same comments apply to the subsequent slides.



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• Note that $\left\lceil \frac{n}{2} \right\rceil = 2^{k-1}$ if $n = 2^k$.



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 - prove that f, with $f(n) := n^2$, is smooth,
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• For $a,b\in\mathbb{R}_0^+$ we define $g\colon\mathbb{N}\to\mathbb{R}_0^+$ as

- Note that $\lceil \frac{n}{2} \rceil = 2^{k-1}$ if $n = 2^k$.
- We would like to show that $g \in \Theta(n^2)$: It suffices to
 - prove that f, with $f(n) := n^2$, is smooth,
 - prove that $g \in \Theta(f \mid \text{``is power of 2''})$,
 - prove that g is eventually non-decreasing.



• For $a, b \in \mathbb{R}_0^+$ we define $g: \mathbb{N} \to \mathbb{R}_0^+$ as

- Note that $\left\lceil \frac{n}{2} \right\rceil = 2^{k-1}$ if $n = 2^k$.
- We would like to show that $g \in \Theta(n^2)$: It suffices to
 - prove that f, with $f(n) := n^2$, is smooth,
 - prove that $g \in \Theta(f \mid \text{``is power of 2''})$,
 - prove that q is eventually non-decreasing.
- Standard application in computer science: Solving the recurrence relation

$$T(n) = T\left(\left\lceil\frac{n}{2}\right\rceil\right) + T\left(\left\lfloor\frac{n}{2}\right\rfloor\right) + b \cdot n,$$

e.g., as derived when analyzing the complexity of merge sort.



- 6
 - **Complexity Analysis and Recurrence Relations**
 - Growth Rates
 - Bachmann-Landau (Asymptotic) Notation
 - Recurrence Relations
 - Heuristics for Solving Recurrences
 - Solving Linear Recurrence Relations
 - Master Theorem



• Sample sequence $t: \mathbb{N}_0 \to \mathbb{R}$: (1, 2, 4, 8, 16, 32, 64, 128, 256, ...)



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Definition 209 (Recurrence relation, Dt.: Rekurrenzgleichung)

A *recurrence relation* for a sequence t is an equation that relates elements of t. It is of order k, for some $k \in \mathbb{N}$, if t_n can be expressed in terms of n and $t_{n-1}, t_{n-2}, \ldots, t_{n-k}$, i.e., if t_n is of the form $t_n = f(t_{n-1}, t_{n-2}, \ldots, t_{n-k}, n)$ for $f : \mathbb{R}^k \times \mathbb{N} \to \mathbb{R}$ (or for $f : \mathbb{R}^k \times \mathbb{N}_0 \to \mathbb{R}$).



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Recurrence relation (of order 1) for the sample sequence given above:

$$t_n := \left\{ \begin{array}{ll} 1 & \text{if } n = 0, \\ 2 \cdot t_{n-1} & \text{if } n > 0. \end{array} \right.$$



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• Easy to see: $t_n = 2^n$ for all $n \in \mathbb{N}_0$.



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Note

We will freely mix the notations t_k and t(k) for denoting the k-th element of a sequence $(t_n)_{n\in\mathbb{N}}$ or $(t_n)_{n\in\mathbb{N}_0}$.



 According to legend, life on Earth will end once the Brahmin priests managed to move the last disk in their 64-disk Tower-of-Hanoi problem . . .



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 - \bullet the top n-1 disks (recursively) from pole I to the auxiliary pole III,
 - the largest (bottom-most) disk from pole I to pole II,
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 - \bullet the top n-1 disks (recursively) from pole III to pole II.
- If T(n) denotes the number of moves for the n-disk ToH problem, the priests need two times T(n-1) moves for the recursive Steps (1) and (3), and one move for getting the largest disk from pole I to II in Step (2).
- Of course, T(1) = 1.



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- Hence, we get the recurrence relation

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for the number *T* of moves for solving the Tower-of-Hanoi problem recursively.



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for the number T of moves for solving the Tower-of-Hanoi problem recursively.

- A solution of this recurrence relation tells us when life on Earth might end . . .
- So, is it already time for an apocalyptic mood?
- We start with heuristics for solving recurrence relations.



- Constructive Induction:
 - First "guess" a solution.
 - Use "constructive" induction to verify that the solution guessed is correct.



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 - Restate the recurrence relation for $t_n, t_{n-1}, t_{n-2}, \ldots$
 - Manipulate and rearrange the individual equations such that summing over all equations yields a closed-form expression for t_n .



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Note

All heuristics require induction to prove that the result obtained is indeed correct!



- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
- Guess: $t \in O(f)$ for $f(n) := n^2$.



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



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$$\stackrel{a:=2}{=} 2(n^2 + 2n + 1)$$

$$= 2(n+1)^2.$$

• Now use standard induction to show that $t_n \leq 2n^2$ is indeed correct for all $n \in \mathbb{N}_0$.



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- Restating the recurrence yields the following set of equations:

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$$\vdots$$

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Heuristics for Solving Recurrences: Cascading

- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
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$$t_{n} = t_{0} + 1 + 2 + \dots + (n-2) + (n-1) + n$$



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$$t_{n} = t_{0} + 1 + 2 + \dots + (n-2) + (n-1) + n$$

$$= 0 + 1 + 2 + \dots + (n-2) + (n-1) + n$$

This indicates that

$$t_n = \sum_{i=0}^{n} i = \frac{n(n+1)}{2} \in \Theta(n^2),$$

which is proved by induction.



- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
- Iterating the recurrence yields

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- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
- Iterating the recurrence yields

$$t_{n} = t_{n-1} + n$$

$$= (t_{n-2} + (n-1)) + n = t_{n-2} + ((n-1)) + n)$$

$$= (t_{n-3} + (n-2)) + ((n-1) + n) = t_{n-3} + ((n-2) + (n-1) + n)$$



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$$= (t_{n-4} + (n-3)) + ((n-2) + (n-1) + n)$$



- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
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$$= (t_{n-3} + (n-2)) + ((n-1) + n) = t_{n-3} + ((n-2) + (n-1) + n)$$

$$= (t_{n-4} + (n-3)) + ((n-2) + (n-1) + n)$$

$$\vdots$$

$$= t_{0} + 1 + 2 + \dots + (n-1) + n$$

$$= 0 + 1 + 2 + \dots + (n-1) + n.$$



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Again, this indicates that

$$t_n = \sum_{i=0}^n i = \frac{n(n+1)}{2} \in \Theta(n^2),$$

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• We have the Tower-of-Hanoi recurrence relation

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$$= 2^{2}(2^{1}T(n-3) + 2^{0}) + 2^{1} + 2^{0} = 2^{3}T(n-3) + 2^{2} + 2^{1} + 2^{0}$$



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$$\vdots$$

$$= 2^{n-1}T(n-(n-1)) + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$$

$$= 2^{n-1} + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$$



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$$\vdots$$

$$= 2^{n-1}T(n - (n-1)) + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$$

$$= 2^{n-1} + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$$

$$= 2^{n} - 1$$



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 with $T(1) := 1$.

Iteration yields the following identities:

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$$= 2(2^{1}T(n-2) + 2^{0}) + 2^{0} = 2^{2}T(n-2) + 2^{1} + 2^{0}$$

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$$= 2^{n-1} + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$$

$$= 2^{n} - 1$$

• Hence, if the priests manage to move one disk per second then we would have to expect the end of Earth $2^{64}-1$ seconds after they started, i.e., roughly within $5\cdot 10^{11}$ years ...

UNIVERSITÄT SALZBLIRG Computational Gacretry and Applications Lab

Definition 210 (Homogeneous recurrence, Dt.: homogene Rekurrenz)

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$$t_n := 3 \cdot n^2 \cdot t_{n-1} \cdot t_{n-2}$$
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Definition 211 (Linear homogeneous recurrence)

A homogeneous recurrence relation of order k is *linear* if $t_n = \sum_{i=1}^k a_i(n) \cdot t_{n-i}$, where $a_i : \mathbb{N} \to \mathbb{R}$ for $i = 1, 2, \dots, k$.



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• E.g.,
$$t_n := n^2 \cdot t_{n-1} + 3 \cdot t_{n-2}$$
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• E.g., $t_n := 2 \cdot t_{n-1} + 3 \cdot t_{n-2}$.



Lemma 213

Consider the recurrence relation $a_0t_n+a_1t_{n-1}+\cdots+a_kt_{n-k}=0$, with $a_i\in\mathbb{R}$. If (f_n) and (g_n) satisfy the recurrence relation then $(\alpha f_n+\beta g_n)$ satisfies the recurrence relation for all $\alpha,\beta\in\mathbb{R}$.



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Proof: Suppose that

$$a_0 f_n + a_1 f_{n-1} + \dots + a_k f_{n-k} = \sum_{i=0}^k a_i f_{n-i} = 0$$
 and $\sum_{i=0}^k a_i g_{n-i} = 0$

for all $n \ge k$.



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for all $n \ge k$. Let $\alpha, \beta \in \mathbb{R}$ arbitrary but fixed and consider $(\alpha f_n + \beta g_n)$. We get

$$\sum_{i=0}^k a_i(\alpha f_{n-i} + \beta g_{n-i}) = \alpha \sum_{i=0}^k a_i f_{n-i} + \beta \sum_{i=0}^k a_i g_{n-i} = 0.$$



- So, consider $a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0$
- Guess $t_n = x^n$ for some unknown $x \in \mathbb{R}$.



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- Then $a_0x^n + a_1x^{n-1} + \cdots + a_kx^{n-k} = 0$.



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- Further $x^{n-k}(a_0x^k + a_1x^{k-1} + \cdots + a_k) = 0$.



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- Further $x^{n-k}(a_0x^k + a_1x^{k-1} + \cdots + a_k) = 0$.
- If we ignore the trivial solution x := 0 then we get

$$a_0x^k + a_1x^{k-1} + \cdots + a_k = 0$$

as the so-called *characteristic equation* of the recurrence relation

$$a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0.$$

• Hence, any root r of this equation serves as a partial solution of the recurrence relation, with $t_n := r^n$.



• Suppose that the characteristic equation has k distinct roots r_1, \ldots, r_k such that all roots are real numbers. I.e., the characteristic equation is given as

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Then, the general solution of the recurrence relation is of the form

$$t_n = \sum_{i=1}^k c_i \cdot r_i^n,$$

for some constants $c_1, c_2, \ldots, c_k \in \mathbb{R}$.



Solving Linear Homogeneous Recurrence Relations With Constant Coefficients

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• The constants c_i are determined based on the initial condition(s).



Solving Linear Homogeneous Recurrence Relations With Constant Coefficients: Fibonacci Sequence

Consider the Fibonacci sequence (over N₀)

$$F_n := \begin{cases} 0 & \text{if } n = 0, \\ 1 & \text{if } n = 1, \\ F_{n-1} + F_{n-2} & \text{if } n \geqslant 2. \end{cases}$$



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• Hence, $F_n - F_{n-1} - F_{n-2} = 0$



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$$x^2 - x - 1 = 0$$

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• This characteristic equation has the roots

$$r_1 := \frac{1 + \sqrt{5}}{2}$$
 and $r_2 := \frac{1 - \sqrt{5}}{2}$.

• Note: r_1 is known as the *golden ratio*, ϕ , with $\phi \approx 1.618$.



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- Note: r_1 is known as the *golden ratio*, ϕ , with $\phi \approx 1.618$.
- This yields

$$F_n = c_1 \cdot \left(\frac{1+\sqrt{5}}{2}\right)^n + c_2 \cdot \left(\frac{1-\sqrt{5}}{2}\right)^n.$$



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$$F_n = c_1 \cdot \left(\frac{1+\sqrt{5}}{2}\right)^n + c_2 \cdot \left(\frac{1-\sqrt{5}}{2}\right)^n.$$

• The constants c_1, c_2 are determined by resorting to the initial conditions.

$$n := 0:$$
 $F_0 = 0 = c_1 + c_2$
 $n := 1:$ $F_1 = 1 = c_1 \cdot \frac{1 + \sqrt{5}}{2} + c_2 \cdot \frac{1 - \sqrt{5}}{2}$



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- By solving this linear system we obtain $c_1 = -c_2 = \frac{1}{\sqrt{5}}$.
- Hence,

$$F_n = \frac{1}{\sqrt{5}} \cdot \left(\frac{1+\sqrt{5}}{2}\right)^n - \frac{1}{\sqrt{5}} \cdot \left(\frac{1-\sqrt{5}}{2}\right)^n.$$



Solving Linear Homogeneous Recurrence Relations With Constant Coefficients

• *Multiple roots*: Suppose that the characteristic equation has s distinct roots r_1, \ldots, r_s of multiplicities m_1, \ldots, m_s such that all roots are real numbers. I.e., the characteristic equation is given as

$$\prod_{i=1}^s (x-r_i)^{m_i}=0.$$



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Then we have

$$t_n = \sum_{i=1}^s \sum_{j=0}^{m_i-1} c_{ij} \cdot n^j \cdot r_i^n,$$

for constants $c_{ij} \in \mathbb{R}$.



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• E.g., for the characteristic equation $(x-1) \cdot (x-2)^2 = 0$ we have s = 2, $r_1 = 1$, $r_2 = 2$, $m_1 = 1$, $m_2 = 2$,



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• E.g., for the characteristic equation $(x-1) \cdot (x-2)^2 = 0$ we have s = 2, $r_1 = 1$, $r_2 = 2$, $m_1 = 1$, $m_2 = 2$, and get

$$t_n = c_{10} \cdot n^0 \cdot 1^n + c_{20} \cdot n^0 \cdot 2^n + c_{21} \cdot n^1 \cdot 2^n = c_{10} + c_{20} \cdot 2^n + c_{21} \cdot n \cdot 2^n$$



Solving Linear Inhomogeneous Recurrence Relations With Constant Coefficients

Assume we have an inhomogeneous recurrence relation of the following form:

$$a_0 \cdot t_n + a_1 \cdot t_{n-1} + \cdots + a_k \cdot t_{n-k} = b_1^n \cdot p_1(n) + b_2^n \cdot p_2(n) + \cdots + b_t^n \cdot p_t(n),$$

where $t \in \mathbb{N}_0$ and b_i is constant and p_i is a polynomial in n of degree $d_i \in \mathbb{N}_0$ for each $1 \le i \le t$.



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where $t \in \mathbb{N}_0$ and b_i is constant and p_i is a polynomial in n of degree $d_i \in \mathbb{N}_0$ for each $1 \le i \le t$.

• Then the characteristic polynomial is

$$(a_0 \cdot x^k + a_1 \cdot x^{k-1} + \cdots + a_k) \cdot \prod_{i=1}^t (x - b_i)^{d_i+1} = 0.$$

Now proceed as in the homogeneous case.



Solving Linear Inhomogeneous Recurrence Relations With Constant Coefficients

Theorem 214

Consider the linear inhomogeneous recurrence relation

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$$(a_0x^k + a_1x^{k-1} + \cdots + a_k) \cdot \prod_{i=1}^t (x - b_i)^{d_i+1} = 0$$

has s distinct roots r_1, \ldots, r_s of multiplicities m_1, \ldots, m_s such that all roots are real numbers.

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for constants $c_{ii} \in \mathbb{R}$.

Consider

$$t_n := \left\{ \begin{array}{ll} 0 & \text{if } n = 0, \\ 2t_{n-1} + n + 2^n & \text{otherwise}. \end{array} \right.$$



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Hence, relative to Thm. 214, we get

$$k = 1$$
 $a_0 = 1$ $a_1 = -2$ $t = 2$

$$b_1 = 1$$
 $p_1(n) = n$ $d_1 = 1$ $b_2 = 2$ $p_2(n) = 1$ $d_2 = 0$.



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This results in

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as the characteristic equation, and we get, with $r_1 := 1, r_2 := 2, m_1 = m_2 := 2,$

$$t_n = c_{10} \cdot n^0 \cdot 1^n + c_{11} \cdot n^1 \cdot 1^n + c_{20} \cdot n^0 \cdot 2^n + c_{21} \cdot n^1 \cdot 2^n$$

= $c_{10} + c_{11} \cdot n + c_{20} \cdot 2^n + c_{21} \cdot n \cdot 2^n$.



So, we know that

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• The constants $c_{10}, c_{11}, c_{20}, c_{21}$ are determined by resorting to the initial conditions:

$$n := 0: 0 = c_{10} + c_{11} \cdot 0 + c_{20} \cdot 2^{0} + c_{21} \cdot 0 \cdot 2^{0} = c_{10} + c_{20}$$

$$n := 1: 3 = c_{10} + c_{11} + 2 \cdot c_{20} + 2 \cdot c_{21}$$

$$n := 2: 12 = c_{10} + 2 \cdot c_{11} + 4 \cdot c_{20} + 8 \cdot c_{21}$$

$$n := 3: 35 = c_{10} + 3 \cdot c_{11} + 8 \cdot c_{20} + 24 \cdot c_{21}$$



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• Solving this system of four linear equations for $c_{10}, c_{11}, c_{20}, c_{21}$ yields

$$c_{10} = -2,$$
 $c_{11} = -1,$ $c_{20} = 2,$ $c_{21} = 1.$



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We conclude that

$$t_n = -2 - n + 2 \cdot 2^n + n \cdot 2^n$$
, i.e., $t_n = -2 - n + 2^{n+1} + n \cdot 2^n$.



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 - **Complexity Analysis and Recurrence Relations**
 - Growth Rates
 - Bachmann-Landau (Asymptotic) Notation
 - Recurrence Relations
 - Master Theorem



Master Theorem

Theorem 215 (Master theorem, Dt.: Hauptsatz der Laufzeitfunktionen)

Consider constants $c \in \mathbb{R}^+$, $k, n_0 \in \mathbb{N}$ and $a, b \in \mathbb{N}$ with $b \ge 2$, and let $T : \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = a \cdot T\left(\frac{n}{b}\right) + c \cdot n^{k}$$

for all $n \in \mathbb{N}$ with $n \ge n_0$, where we interpret $T(\frac{n}{b})$ as (a combination of) $T(\lceil \frac{n}{b} \rceil)$ or $T(\lfloor \frac{n}{b} \rfloor)$.



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Then we have

$$T \in \left\{ \begin{array}{ll} \Theta(n^k) & \text{if } a < b^k, \\ \Theta(n^k \log n) & \text{if } a = b^k, \\ \Theta(n^{\log_b a}) & \text{if } a > b^k. \end{array} \right.$$



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• E.g., we get $T \in \Theta(n \log n)$ for T defined as follows:

$$T(n) = T\left(\left\lceil \frac{n}{2} \right\rceil\right) + T\left(\left\lceil \frac{n}{2} \right\rceil\right) + c \cdot n.$$



Master Theorem (Asymptotic Version)

Theorem 216

Consider constants $k, n_0 \in \mathbb{N}$ and $a, b \in \mathbb{N}$ with $b \geqslant 2$, and a function $f \colon \mathbb{N} \to \mathbb{R}_0^+$ with $f \in \Theta(n^k)$. Let $T \colon \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = a \cdot T\left(\frac{n}{b}\right) + f(n)$$

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Master Theorem (Refined Asymptotic Version)

Theorem 217

Consider constants $n_0 \in \mathbb{N}$ and $a \in \mathbb{N}$, $b \in \mathbb{R}$ with b > 1, and a function $f : \mathbb{N} \to \mathbb{R}_0^+$. Let $T : \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = a \cdot T\left(\frac{n}{b}\right) + f(n)$$

for all $n \in \mathbb{N}$ with $n \ge n_0$, where we interpret $T(\frac{n}{b})$ as (a combination of) $T(\lceil \frac{n}{b} \rceil)$ or $T(\lfloor \frac{n}{b} \rfloor)$.



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Consider constants $n_0 \in \mathbb{N}$ and $a \in \mathbb{N}$, $b \in \mathbb{R}$ with b > 1, and a function $f : \mathbb{N} \to \mathbb{R}_0^+$. Let $T: \mathbb{N} \to \mathbb{R}_n^+$ be an eventually non-decreasing function such that

$$T(n) = a \cdot T\left(\frac{n}{b}\right) + f(n)$$

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Then we have

$$T \in \left\{ \begin{array}{ll} \Theta(f) & \text{if} \left\{ \begin{array}{ll} f \in \Omega(n^{(\log_b a) + \varepsilon}) \text{ for some } \varepsilon \in \mathbb{R}^+, \\ \text{ and if the following regularity condition holds} \\ \text{ for some } 0 < s < 1 \text{ and all sufficiently large n:} \\ a \cdot f(n/b) \leqslant s \cdot f(n), \\ \Theta\left(n^{\log_b a} \log n\right) & \text{ if } f \in \Theta(n^{\log_b a}), \\ \Theta(n^{\log_b a}) & \text{ if } f \in O(n^{(\log_b a) - \varepsilon}) \text{ for some } \varepsilon \in \mathbb{R}^+. \end{array} \right.$$

This is a simplified version of the Akra-Bazzi Theorem [Akra&Bazzi 1998].



Real-World Application: Analysis of Fast Integer Multiplication

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$$(a_{2n-1}a_{2n-2}\cdots a_1a_0)_2$$
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be the 2*n*-bit binary representations of *a* and *b*. Hence, $a = \sum_{i=0}^{2n-1} a_i 2^i$ and $b = \sum_{i=0}^{2n-1} b_i 2^i$.



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We have

$$a \sim 2^n A_1 + A_0$$
 and $b \sim 2^n B_1 + B_0$

with

$$A_1 := (a_{2n-1}a_{2n-2}\cdots a_{n+1}, a_n)_2, \quad A_0 := (a_{n-1}a_{n-2}\cdots a_1a_0)_2,$$

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Definition 218 (Graph, Dt.: (schlichter endlicher ungerichteter) Graph)

For $n \in \mathbb{N}$ and $m \in \mathbb{N}_0$, a (simple finite undirected) graph $\mathcal{G} := (V, E)$ with n vertices (aka nodes) and m edges consists of a vertex set $V := \{v_1, v_2, \ldots, v_n\}$ and an edge set $E := \{e_1, e_2, \ldots, e_m\}$, where $V \cap E = \emptyset$ and each edge is an unordered pair of distinct vertices:



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a

e

a

e

(C)

(C)

(b)

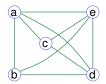
d

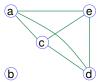
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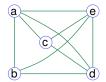


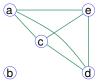


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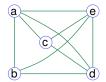


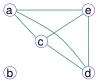
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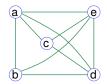


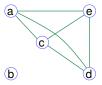
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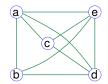


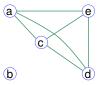


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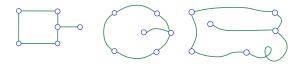
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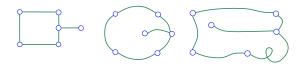
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- If we allow edges of the form *uu* then we get a *loop* (Dt.: Schlinge, Schleife) and the graph is no longer simple (Dt.: schlicht, einfach).
- If we allow multiple edges between two vertices then we get a multigraph.

- Graphical representation of a graph:
 - Denote the vertices by markers of the same form (circles, dots, squares, ...).
 - For every pair of vertex markers, draw a curve between them if the graph contains an edge between the corresponding vertices.



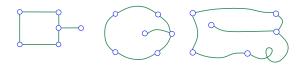


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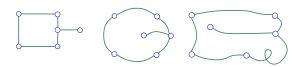


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- Use arrows to denote directed edges.





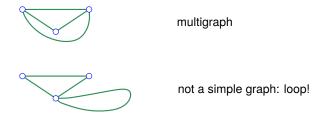
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multigraph

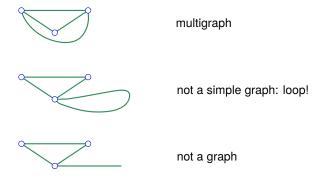


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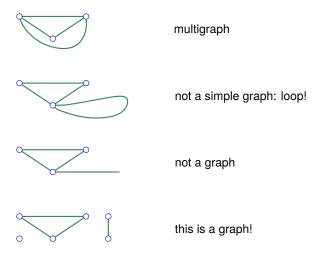


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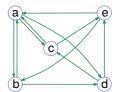
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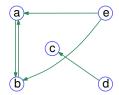




Definition 219 (Directed graph, Dt.: (schlichter endlicher) gerichteter Graph)

For $n \in \mathbb{N}$ and $m \in \mathbb{N}_0$, a (simple finite) directed graph, or digraph, $\mathcal{G} := (V, E)$ with n vertices (aka nodes) and m edges consists of a vertex set $V := \{v_1, v_2, \ldots, v_n\}$ and an edge set $E := \{e_1, e_2, \ldots, e_m\}$, where $V \cap E = \emptyset$ and each edge is an ordered pair of distinct vertices:



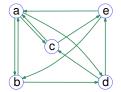


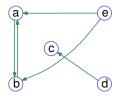


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For $n \in \mathbb{N}$ and $m \in \mathbb{N}_0$, a (simple finite) directed graph, or digraph, $\mathcal{G} := (V, E)$ with n vertices (aka nodes) and m edges consists of a vertex set $V := \{v_1, v_2, \ldots, v_n\}$ and an edge set $E := \{e_1, e_2, \ldots, e_m\}$, where $V \cap E = \emptyset$ and each edge is an ordered pair of distinct vertices:

$$E \subseteq \{(u, v) : u, v \in V \text{ and } u \neq v\}.$$





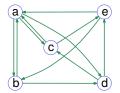
For a digraph, uv indicates the edge (u, v), i.e., an edge where u is the tail and v is the head.

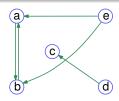


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- For a digraph, uv indicates the edge (u, v), i.e., an edge where u is the tail and v is the head.
- In this lecture we will always specify a directed graph explicitly; that is, the term "graph" without the qualifier "directed" shall mean "undirected graph".

Basic Definitions: How to Deal with $V = \emptyset$

- There is no consensus on whether or not to allow V = Ø in the definition of a graph. (Of course, if V = Ø then E = Ø.)
- And, indeed, there are pros and cons of allowing $V = \emptyset$.



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- Furthermore, if V = Ø is allowed then there is little consensus on how to call such a graph:
 - Common terms are *order-zero graph*, K_0 , and *null graph*.
 - Some authors also use the term *empty graph* to indicate $V = \emptyset$ while other authors prefer to reserve this term for a graph with $E = \emptyset$ but $V \neq \emptyset$.



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Convention

We will always assume that every (directed) graph has at least one node.



Basic Definitions — Warning!

No common terminology

The terminology in graph theory lacks a rigorous standardization, both in the German and in the English literature.



Basic Definitions — Warning!

No common terminology

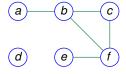
The terminology in graph theory lacks a rigorous standardization, both in the German and in the English literature.

- In several cases the meanings of different terms coincide for simple undirected graphs, which seems to serve as a justification for authors to freely mix and match terms.
- Thus, always make sure to check how some author defines standard terms of graph theory . . .



Undirected Graphs as Directed Graphs

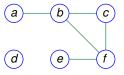
• It is straightforward to represent an undirected graph as a directed graph.

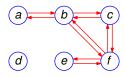




Undirected Graphs as Directed Graphs

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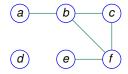


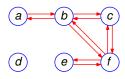




Undirected Graphs as Directed Graphs

- It is straightforward to represent an undirected graph as a directed graph.
- Hence, undirected graphs can be seen as a special case of directed graphs, and most algorithms that work for directed graphs are applicable to undirected graphs, too.







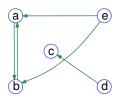
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- E.g., the relation R on the set $\{a, b, c, d, e\}$, with

$$R := \{(a,b), (b,a), (d,c), (e,a), (e,b)\},\$$

corresponds to the following directed graph:

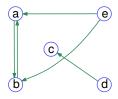




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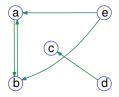
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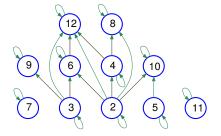


- Hence, statements about relations can be translated to statements about digraphs, and vice versa.
- Note, though, that the digraph corresponding to a relation
 - need not be simple but might contain loops,
 - need not have a finite vertex set.
- Simplified representation of the digraph of an order relation: Hasse diagram

• Consider the poset (S, R), where $S := \{n \in \mathbb{N} : 1 < n \le 12\}$ and R denotes the partial order of divisibility on S. (That is, for $a, b \in S$, we have a R b iff $a \mid b$.)

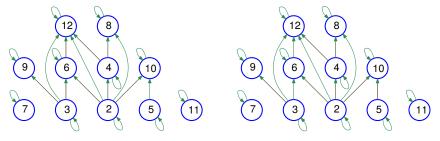


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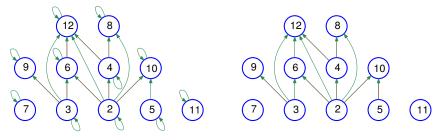
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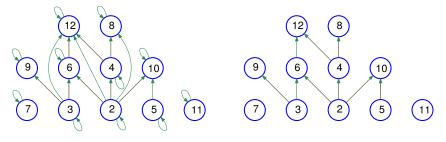


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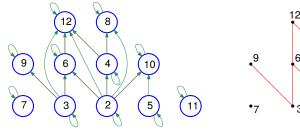
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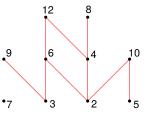


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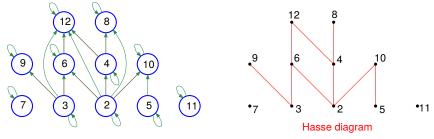




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Definition 220 (Hasse diagram)

The graph obtained after carrying out Steps (1)–(4) is the *Hasse diagram* of the poset.

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• Typically, some statements of a computer program could be executed in parallel.

- (1) a := 1
- (2) b := 2
- (3) c := 3
- (4) d := a + 2
- (5) e := 2a + b
- (6) f := d + c
- (7) g := c + e
- (8) h := d + e + f



- Typically, some statements of a computer program could be executed in parallel.
- A precedence graph is a directed graph that models dependences. E.g., the dependence of statements of a computer program on other statements:
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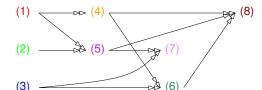
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$$b := 2$$

(3)
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$$(4) d := a + 2$$

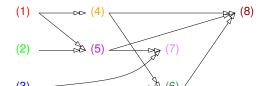
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- Precedence graphs are used in all sorts of scheduling tasks: E.g., job scheduling, concurrency control and instruction scheduling, resolving linker dependencies, data serialization, automated parallelization of sequential code.
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Definition 221 (Adjacent, Dt.: benachbart)

Two vertices $u, v \in V$ of a graph $\mathcal{G} := (V, E)$ are adjacent if $uv \in E$; the edge uv is incident to the vertices u and v.



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The *degree* (aka *valence*) of a vertex u of a graph $\mathcal{G} := (V, E)$ is the number of edges incident to u. It is denoted by deg(u).



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The *degree of a graph* is the maximum of the degrees of its vertices.



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Definition 223 (Subgraph, Dt.: Teilgraph)

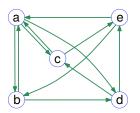
A graph $\mathcal{G}' := (V', E')$ is a *subgraph* of a (directed) graph $\mathcal{G} := (V, E)$ if $V' \subseteq V$ and $E' \subseteq E$ such that all edges of E' are formed by vertices of V'.



Definition 224 (Adjacency matrix, Dt.: Adjazenzmatrix)

The *adjacency matrix* of a (directed) graph $\mathcal{G} := (V, E)$ is an $n \times n$ matrix \mathbf{M} , where n := |V| and

$$m_{ij} := \left\{ egin{array}{ll} 1 & ext{if } v_i v_j \in E, \ 0 & ext{otherwise.} \end{array}
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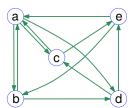
X	а	b	С	d	е
а	0	1	1	1	0
b	1	0	0	1	0
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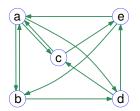
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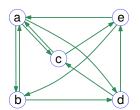
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- Note: Storing **M** (as an $n \times n$ array) requires $\Theta(n^2)$ memory!
- Adjacency lists (and their variants) help to preserve memory if $|E| \ll |V|^2$.

Basic Definitions: Regularity

Definition 225 (Regular graph, Dt.: regulärer Graph)

A graph \mathcal{G} is *regular* if every vertex of \mathcal{G} has the same degree. A regular graph with vertices of degree k is called a k-regular graph or regular graph of degree k.

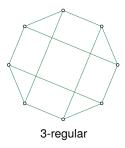


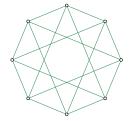
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4-regular

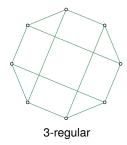


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- A 3-regular graph is known as a cubic graph, and a 4-regular graph is known as a quartic graph.
- For directed regular graphs it is common to demand that the in-degree and the out-degree of each vertex is identical.



4-regular



Lemma 226 (Degree sum formula)

The sum over all degrees of vertices of a graph $\mathcal{G}:=(V,E)$ equals twice the number of its edges, i.e., $\sum_{\nu\in V} \deg(\nu)=2|E|$.



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Sketch of proof: Adding one edge increases the sum of the degrees by two.



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Corollary 227 (Euler's Handshaking Lemma, Dt.: Handschlag-Lemma)

In every graph the number of vertices of odd degree is even.



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- Simple application of Euler's Handshaking Lemma:
 - Suppose that a party is attended by 15 guests. Is it possible that every guest at the party knows all others except for precisely one guest?
 - No: Consider a graph with 15 nodes (guests) where two nodes are linked by an edge if the corresponding guests do not know each other. Hence, we would get 15 nodes of degree one, in contradiction to Cor. 227.



- Graph Theory
 - What is a (Directed) Graph?
 - Paths
 - Walks
 - Connectedness
 - Euler Tour and Hamilton Cycle
 - Trees
 - Special Graphs
 - Graph Coloring



Walks

Definition 228 (Walk, Dt.: Wanderung, Kantenfolge)

A walk of length k, with $k \in \mathbb{N}_0$, on $\mathcal{G} := (V, E)$ is an alternating sequence

$$V_0e_1V_1e_2V_2\ldots e_kV_k$$

of k + 1 vertices $v_0, v_1, \dots, v_k \in V$ and k edges $e_1, \dots, e_k \in E$ such that

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Often, a walk of length k is written simply as

$$V_0V_1V_2\ldots V_k$$
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A walk of length k, with $k \in \mathbb{N}_0$, on $\mathcal{G} := (V, E)$ is an alternating sequence

$$V_0e_1V_1e_2V_2\ldots e_kV_k$$

of k + 1 vertices $v_0, v_1, \dots, v_k \in V$ and k edges $e_1, \dots, e_k \in E$ such that

$$\forall (1 \leqslant i \leqslant k) \ e_i = v_{i-1}v_i.$$

Often, a walk of length k is written simply as

$$V_0V_1V_2\ldots V_k$$
.

• Conventionally, v_0 is called the *start vertex* (or *initial vertex*) of the walk, and v_k is called its *end vertex* (or *terminal vertex*). Note that $v_{i-1} \neq v_i$ for $i \in \{1, 2, ..., k\}$.



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Definition 228 (Walk, Dt.: Wanderung, Kantenfolge)

A walk of length k, with $k \in \mathbb{N}_0$, on $\mathcal{G} := (V, E)$ is an alternating sequence

$$V_0e_1V_1e_2V_2\ldots e_kV_k$$

of k + 1 vertices $v_0, v_1, \dots, v_k \in V$ and k edges $e_1, \dots, e_k \in E$ such that

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Often, a walk of length k is written simply as

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Definition 229 (Closed walk, Dt.: geschlossene Wanderung)

A walk is called *closed* if the start vertex and the end vertex are identical. A closed walk of length k is called *trivial* if $k \le 2$.



Definition 230 (Trail, Dt.: Weg)

A *trail* in a (directed) graph G is a walk in which all edges are distinct.



Definition 230 (Trail, Dt.: Weg)

A *trail* in a (directed) graph \mathcal{G} is a walk in which all edges are distinct.

Definition 231 (Path, Dt.: Pfad)

A path in a (directed) graph \mathcal{G} is a walk in which all vertices are distinct.



Definition 230 (Trail, Dt.: Weg)

A trail in a (directed) graph $\mathcal G$ is a walk in which all edges are distinct.

Definition 231 (Path, Dt.: Pfad)

A path in a (directed) graph $\mathcal G$ is a walk in which all vertices are distinct.

Definition 232 (Tour, Dt.: Tour)

A *tour* in a (directed) graph G is a closed trail.



Definition 230 (Trail, Dt.: Weg)

A *trail* in a (directed) graph G is a walk in which all edges are distinct.

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A path in a (directed) graph $\mathcal G$ is a walk in which all vertices are distinct.

Definition 232 (Tour, Dt.: Tour)

A *tour* in a (directed) graph G is a closed trail.

Definition 233 (Cycle, Dt.: Zyklus, Kreis)

A *cycle* in a (directed) graph $\mathcal G$ is a non-trivial closed walk in which all but the start and the end vertices are distinct.



Definition 230 (Trail, Dt.: Weg)

A *trail* in a (directed) graph G is a walk in which all edges are distinct.

Definition 231 (Path, Dt.: Pfad)

A path in a (directed) graph G is a walk in which all vertices are distinct.

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A *cycle* in a (directed) graph $\mathcal G$ is a non-trivial closed walk in which all but the start and the end vertices are distinct.

 Note: Distinct vertices implies distinct edges; i.e., every path is a trail and every cycle is a tour.



Definition 230 (Trail, Dt.: Weg)

A *trail* in a (directed) graph \mathcal{G} is a walk in which all edges are distinct.

Definition 231 (Path, Dt.: Pfad)

A path in a (directed) graph \mathcal{G} is a walk in which all vertices are distinct.

Definition 232 (Tour, Dt.: Tour)

A *tour* in a (directed) graph \mathcal{G} is a closed trail.

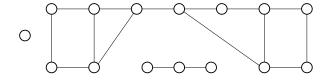
Definition 233 (Cycle, Dt.: Zyklus, Kreis)

A cycle in a (directed) graph \mathcal{G} is a non-trivial closed walk in which all but the start and the end vertices are distinct.

- Note: Distinct vertices implies distinct edges; i.e., every path is a trail and every cycle is a tour.
- Note that some authors prefer to use the terms "path", "simple path", "cycle" and "simple cycle" instead of "trail", "path", "tour" and "cycle" . . .

Definition 234 (Connected component, Dt.: Zusammenhangskomponente)

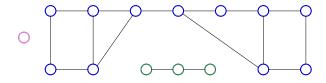
A connected component of a graph $\mathcal{G}:=(V,E)$ is a maximal subgraph $\mathcal{G}':=(V',E')$ of \mathcal{G} such that for every unordered pair $\{u,v\}$, with $u,v\in V'$ and $u\neq v$, there exists a path between u and v within \mathcal{G}' .





Definition 234 (Connected component, Dt.: Zusammenhangskomponente)

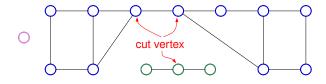
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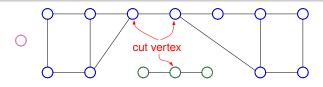
Definition 235 (Cut vertex, Dt.: Artikulationspunkt, Schnittknoten)

A *cut vertex* of a graph $\mathcal{G} := (V, E)$ is a vertex $v \in V$ such that the removal of v and of all edges incident to v would increase the number of connected components.



Definition 234 (Connected component, Dt.: Zusammenhangskomponente)

A *connected component* of a graph $\mathcal{G}:=(V,E)$ is a maximal subgraph $\mathcal{G}':=(V',E')$ of \mathcal{G} such that for every unordered pair $\{u,v\}$, with $u,v\in V'$ and $u\neq v$, there exists a path between u and v within \mathcal{G}' .



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Definition 236 (Connected, Dt.: zusammenhängend)

A graph is connected if it contains only one connected component.



Definition 237 (Weakly connected, Dt.: schwach zusammenhängend)

A directed graph is *weakly connected* if replacing all its directed edges by undirected edges results in a connected (undirected) graph.



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A directed graph is *weakly connected* if replacing all its directed edges by undirected edges results in a connected (undirected) graph.

Definition 238 (Strong component, Dt.: starke Zusammenhangskomponente)

A strong component (aka strongly connected component) of a directed graph $\mathcal{G}:=(V,E)$ is a maximal subgraph $\mathcal{G}'=(V',E')$ of \mathcal{G} such that for every ordered pair (u,v), with $u,v\in V'$ and $u\neq v$, there exists a path from u to v within \mathcal{G}' .



Definition 237 (Weakly connected, Dt.: schwach zusammenhängend)

A directed graph is *weakly connected* if replacing all its directed edges by undirected edges results in a connected (undirected) graph.

Definition 238 (Strong component, Dt.: starke Zusammenhangskomponente)

A strong component (aka strongly connected component) of a directed graph $\mathcal{G}:=(V,E)$ is a maximal subgraph $\mathcal{G}'=(V',E')$ of \mathcal{G} such that for every ordered pair (u,v), with $u,v\in V'$ and $u\neq v$, there exists a path from u to v within \mathcal{G}' .

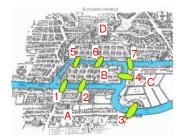
Definition 239 (Strongly connected, Dt.: stark zusammenhängend)

A directed graph $\mathcal{G}:=(V,E)$ is *strongly connected* if it consists of only one strong component, i.e., if for every ordered pair (u,v), with $u,v\in V$ and $u\neq v$, there exists a path from u to v.



Seven Bridges of Königsberg

 Early 18th century: Does there exist a trail (or even a tour) through the city of Königsberg that crosses every of its seven bridges exactly once? (Of course, every bridge had to be crossed fully, and no other means to get across the river Pregel were allowed.)



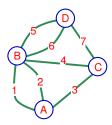
[Image credit for background image: Wikipedia.]



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[Image credit for background image: Wikipedia.]

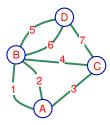
• In 1736, Leonhard Euler (1707–1783) treated this problem as a graph problem and proved, using a parity argument, that such a trail or tour does not exist.



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[Image credit for background image: Wikipedia.]

- In 1736, Leonhard Euler (1707–1783) treated this problem as a graph problem and proved, using a parity argument, that such a trail or tour does not exist.
- His solution is generally regarded as the first theorem of graph theory.



Definition 240 (Euler trail, Dt.: Eulerscher Weg)

An Euler trail is a trail that contains all edges of a graph exactly once.



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An *Euler trail* is a trail that contains all edges of a graph exactly once.

Definition 241 (Euler tour, Dt.: Eulersche Tour)

An Euler tour is a tour that contains all edges of a graph exactly once. A graph is an Eulerian graph if it has an Euler tour.



Definition 240 (Euler trail, Dt.: Eulerscher Weg)

An *Euler trail* is a trail that contains all edges of a graph exactly once.

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An *Euler tour* is a tour that contains all edges of a graph exactly once. A graph is an *Eulerian graph* if it has an Euler tour.

Definition 242 (Hamilton path, Dt.: Hamiltonscher Pfad)

A Hamilton path is a path that passes through all vertices of a graph exactly once.



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Definition 242 (Hamilton path, Dt.: Hamiltonscher Pfad)

A *Hamilton path* is a path that passes through all vertices of a graph exactly once.

Definition 243 (Hamilton cycle, Dt.: Hamiltonscher Kreis)

A *Hamilton cycle* is a cycle that passes through all vertices of a graph exactly once.



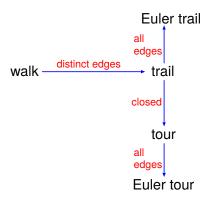
walk





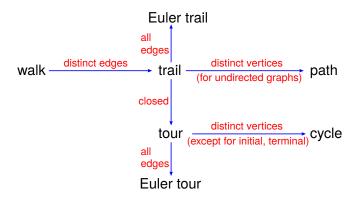




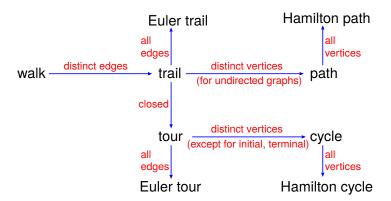




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Euler Tour

Theorem 244

Suppose that every node of a graph $\mathcal G$ has degree at least one. Then $\mathcal G$ has an Euler tour if and only if $\mathcal G$ is connected and every vertex of $\mathcal G$ has even degree.



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Euler Tour

Theorem 244

Suppose that every node of a graph $\mathcal G$ has degree at least one. Then $\mathcal G$ has an Euler tour if and only if $\mathcal G$ is connected and every vertex of $\mathcal G$ has even degree.

Theorem 245

Suppose that every node of a graph $\mathcal G$ has degree at least one. Then $\mathcal G$ has an Euler trail (but no Euler tour) if and only if $\mathcal G$ is connected and exactly two vertices of $\mathcal G$ have odd degrees.



Euler Tour

Theorem 244

Suppose that every node of a graph $\mathcal G$ has degree at least one. Then $\mathcal G$ has an Euler tour if and only if $\mathcal G$ is connected and every vertex of $\mathcal G$ has even degree.

Theorem 245

Suppose that every node of a graph $\mathcal G$ has degree at least one. Then $\mathcal G$ has an Euler trail (but no Euler tour) if and only if $\mathcal G$ is connected and exactly two vertices of $\mathcal G$ have odd degrees.

Corollary 246

An Euler tour or trail in a graph $\mathcal{G}:=(V,E)$ can be determined in O(|E|) time, if it exists. Otherwise, again in O(|E|) time, we can determine that neither an Euler tour nor an Euler trail exists in \mathcal{G} .



Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.



Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.

Suppose that G has an Euler tour T. It is obvious that G is connected.



Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.

Suppose that \mathcal{G} has an Euler tour T. It is obvious that \mathcal{G} is connected. Every occurrence of a vertex $v \in V$ in T is preceded and followed by an edge. Thus, each time T passes through v, two of the edges incident to v are consumed. Since T does neither start nor end in v, it is necessary that $\deg(v)$ is even.



Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.

Suppose that \mathcal{G} has an Euler tour T. It is obvious that \mathcal{G} is connected. Every occurrence of a vertex $v \in V$ in T is preceded and followed by an edge. Thus, each time T passes through v, two of the edges incident to v are consumed. Since T does neither start nor end in v, it is necessary that $\deg(v)$ is even.

Now suppose that every vertex of $\mathcal G$ has even degree, and, of course, that $\mathcal G$ is connected. We give a constructive proof that $\mathcal G$ admits an Euler tour.



Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.

Suppose that $\mathcal G$ has an Euler tour T. It is obvious that $\mathcal G$ is connected. Every occurrence of a vertex $v \in V$ in T is preceded and followed by an edge. Thus, each time T passes through v, two of the edges incident to v are consumed. Since T does neither start nor end in v, it is necessary that $\deg(v)$ is even.

Now suppose that every vertex of $\mathcal G$ has even degree, and, of course, that $\mathcal G$ is connected. We give a constructive proof that $\mathcal G$ admits an Euler tour. Pick any vertex v to start with and trace out a trail T. Every edge that is being traversed is marked. As above, we observe that passing through a vertex that is neither the start nor the end vertex of T consumes two edges.



Constructive Proof of Theorem 244

Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.

Suppose that \mathcal{G} has an Euler tour T. It is obvious that \mathcal{G} is connected. Every occurrence of a vertex $v \in V$ in T is preceded and followed by an edge. Thus, each time T passes through v, two of the edges incident to v are consumed. Since T does neither start nor end in v, it is necessary that $\deg(v)$ is even.

Now suppose that every vertex of $\mathcal G$ has even degree, and, of course, that $\mathcal G$ is connected. We give a constructive proof that $\mathcal G$ admits an Euler tour. Pick any vertex v to start with and trace out a trail T. Every edge that is being traversed is marked. As above, we observe that passing through a vertex that is neither the start nor the end vertex of T consumes two edges.

We realize that, eventually, T will get us back to v. (We cannot be stuck in some other vertex w since w has even degree.)



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We realize that, eventually, T will get us back to v. (We cannot be stuck in some other vertex w since w has even degree.) If at the time when we are back at v every vertex of T has no unmarked incident edge then we are done. Otherwise, we start a new trail T' at a vertex w of T which has an unmarked incident edge and follow it until we get back to w.



Constructive Proof of Theorem 244

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Now suppose that every vertex of $\mathcal G$ has even degree, and, of course, that $\mathcal G$ is connected. We give a constructive proof that $\mathcal G$ admits an Euler tour. Pick any vertex v to start with and trace out a trail T. Every edge that is being traversed is marked. As above, we observe that passing through a vertex that is neither the start nor the end vertex of T consumes two edges.

We realize that, eventually, T will get us back to v. (We cannot be stuck in some other vertex w since w has even degree.) If at the time when we are back at v every vertex of T has no unmarked incident edge then we are done. Otherwise, we start a new trail T' at a vertex w of T which has an unmarked incident edge and follow it until we get back to w.

This process continues until no unmarked edges remain. At the end the trails are spliced together appropriately.

Hamilton Cycle

Theorem 247

It is \mathcal{NP} -complete to determine whether a Hamilton cycle or Hamilton path exists in a general graph.

- Informally, Theorem 247 says that no (deterministic sequential) algorithm is known which determines the existence of a Hamilton cycle or path in an n-vertex graph in a time that is a polynomial function of n.
- Even worse, an efficient (polynomial-time) algorithm will never be found unless $\mathcal{P} = \mathcal{N}\mathcal{P}$ holds, which seems rather unlikely.



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Theorem 248 (Dirac, 1952)

If the degree of every vertex of an n-vertex graph \mathcal{G} , with $n \ge 3$, is at least $\lceil \frac{n}{2} \rceil$ then \mathcal{G} has a Hamilton cycle.



Hamilton Cycle

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Theorem 248 (Dirac, 1952)

If the degree of every vertex of an n-vertex graph \mathcal{G} , with $n \geqslant 3$, is at least $\lceil \frac{n}{2} \rceil$ then \mathcal{G} has a Hamilton cycle.

Theorem 249 (Ore, 1960)

If the sum of the degrees of every pair of non-adjacent vertices of an n-vertex graph \mathcal{G} , with $n \geqslant 3$, is at least n then \mathcal{G} has a Hamilton cycle.

- Graph Theory
 - What is a (Directed) Graph?
 - Paths
 - Trees
 - Basic Definitions
 - Elementary Properties
 - Binary Trees
 - Balance and Height
 - Spanning Trees
 - Recursion Trees
 - Special Graphs
 - Graph Coloring



Definition 250 (Acyclic, Dt.: zyklenfrei)

A graph is called acyclic if it contains no cycles.



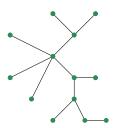
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Definition 250 (Acyclic, Dt.: zyklenfrei)

A graph is called *acyclic* if it contains no cycles.

Definition 251 (Tree, Dt.: Baum)

A tree is an undirected graph that is acyclic and connected.





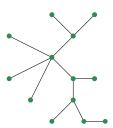
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Definition 251 (Tree, Dt.: Baum)

A tree is an undirected graph that is acyclic and connected.

• For trees most authors prefer to speak about *nodes* rather than vertices.





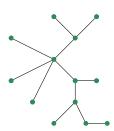
Definition 250 (Acyclic, Dt.: zyklenfrei)

A graph is called *acyclic* if it contains no cycles.

Definition 251 (Tree, Dt.: Baum)

A tree is an undirected graph that is acyclic and connected.

- For trees most authors prefer to speak about *nodes* rather than vertices.
- Unless explicitly stated otherwise, we will only deal with trees that have at least one node. (Some authors call a tree with $V = E = \emptyset$ a *null tree*.)

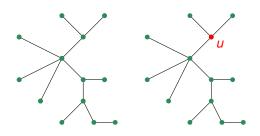




Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

A rooted tree is a directed graph with a node u such that

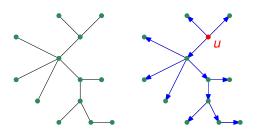
• the graph contains *u* as node ("*root*"),





Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

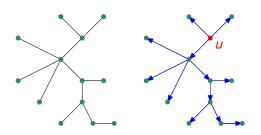
- the graph contains *u* as node ("root"),
- 2 paths from *u* to all other nodes of the graph exist,





Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

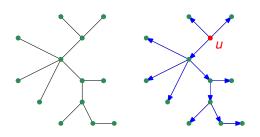
- the graph contains *u* as node ("*root*"),
- 2 paths from *u* to all other nodes of the graph exist,
- the in-degree of u is zero,





Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

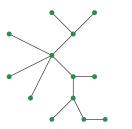
- the graph contains *u* as node ("root"),
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- 4 the in-degree of every other node of the graph is one.

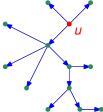


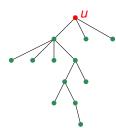


Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

- the graph contains *u* as node ("*root*"),
- 2 paths from u to all other nodes of the graph exist,
- the in-degree of *u* is zero,
- the in-degree of every other node of the graph is one.
 - It is common practice to draw rooted trees from the root downwards such that the (downwards) orientations of the edges are implied by the positions of the nodes.



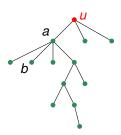






Definition 253 (Child and parent, Dt.: Kind und Eltern)

For a rooted tree $\mathcal{T} := (V, E)$ and nodes $a, b \in V$, the node b is a *child* of the node a, and a is the *parent* of b, if the edge ab belongs to E. *Siblings* are nodes which share the same parent.



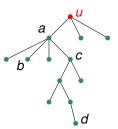


Definition 253 (Child and parent, Dt.: Kind und Eltern)

For a rooted tree $\mathcal{T} := (V, E)$ and nodes $a, b \in V$, the node b is a *child* of the node a, and a is the *parent* of b, if the edge ab belongs to E. *Siblings* are nodes which share the same parent.

Definition 254 (Descendant and ancestor, Dt.: Nachfahre und Vorfahre)

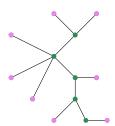
In a rooted tree $\mathcal{T} := (V, E)$, with $c, d \in V$, a node d is a *descendant* of a node c, and c is an *ancestor* of d, if $c \neq d$ and if the path from the root to d contains c.

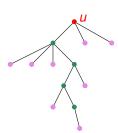




Definition 255 (Leaf, Dt.: Blatt)

A *leaf* of a rooted tree is a node without children. For a tree (that is not rooted) a leaf is a node with degree 1. All non-leaf nodes of a (rooted) tree are called *inner nodes*.



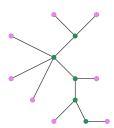


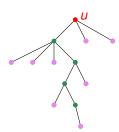


Definition 255 (Leaf, Dt.: Blatt)

A *leaf* of a rooted tree is a node without children. For a tree (that is not rooted) a leaf is a node with degree 1. All non-leaf nodes of a (rooted) tree are called *inner nodes*.

ullet Of course, the root of a rooted tree ${\mathcal T}$ may also be the (only) leaf of ${\mathcal T}$.



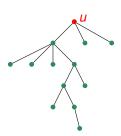




Definition 256 (Subtree, Dt.: Teilbaum)

A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node u if

 $\mathbf{0} \ \mathcal{T}'$ is a subgraph of \mathcal{T} ,

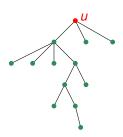




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A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node u if

- \bullet \mathcal{T}' is a subgraph of \mathcal{T} ,
- ② T' is rooted at a node v that is a descendant of u,

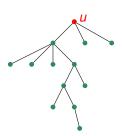




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- 2 \mathcal{T}' is rooted at a node v that is a descendant of u, and
- **3** \mathcal{T}' contains all descendants of v in \mathcal{T} , together with the appropriate edges of E.



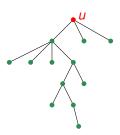


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A subtree rooted at v is called a *proper subtree* if v is a child of u.





Definition 256 (Subtree, Dt.: Teilbaum)

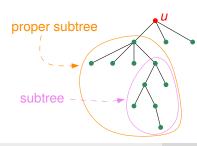
A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node u if

- $\mathbf{0}$ \mathcal{T}' is a subgraph of \mathcal{T} ,
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A subtree rooted at v is called a *proper subtree* if v is a child of u.

Warning

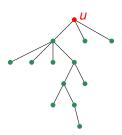
Some authors do not make the distinction between the node v being a child of u or some arbitrary descendant of u.





Definition 257 (Ordered tree, Dt.: geordneter Baum)

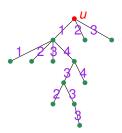
An ordered tree is a rooted tree $\mathcal T$ such that the children of every node of $\mathcal T$ are arranged in some specific order, e.g., by means of a numbering scheme.





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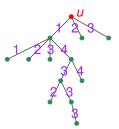


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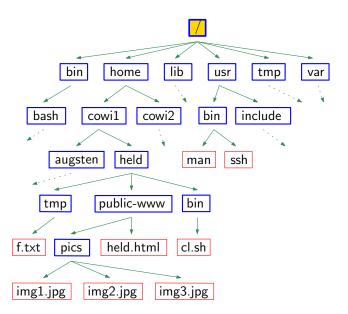
Definition 258 (Forest, Dt.: Wald)

A *forest* is a graph such that all its connected components are trees.



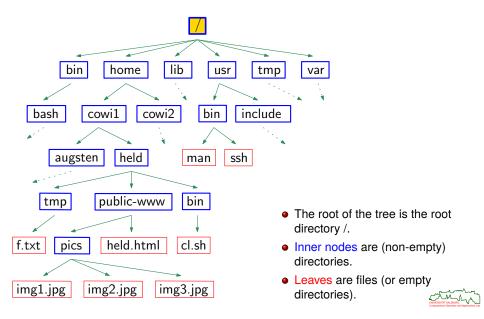


Real-World Application: File System as a Rooted Tree





Real-World Application: File System as a Rooted Tree



Theorem 259

Every pair of nodes in a tree is connected by exactly one path.



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Theorem 259

Every pair of nodes in a tree is connected by exactly one path.

Theorem 260

In a rooted tree there exists exactly one path from the root to any node.



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Theorem 260

In a rooted tree there exists exactly one path from the root to any node.

Lemma 261

Removing an edge from a (rooted) tree results in a graph with two connected components, each of which is a (rooted) tree.



Theorem 262

For every (rooted) tree $\mathcal{T}:=(V,E)$ we get |E|=|V|-1.



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Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

Proof of Theorem 262 for rooted trees: We use structural induction relative to proper subtrees.



Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

Proof of Theorem 262 for rooted trees: We use structural induction relative to proper subtrees. Obviously, the claim holds for the minimal elements, i.e., for trees that contain no proper subtrees and, thus, have only a root and no edges.



Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

Proof of Theorem 262 for rooted trees: We use structural induction relative to proper subtrees. Obviously, the claim holds for the minimal elements, i.e., for trees that contain no proper subtrees and, thus, have only a root and no edges.

Now consider an arbitrary but fixed rooted tree $\mathcal{T} := (V, E)$ and suppose that the equality claimed holds for all its k > 0 proper subtrees $(V_1, E_1), \dots, (V_k, E_k)$. (We do not need to assume explicitly that it holds for all subtrees of \mathcal{T} .)



Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

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$$|E| = k + \sum_{i=1}^k |E_i|$$



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For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

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$$|E| = k + \sum_{i=1}^{k} |E_i| = k + \sum_{i=1}^{k} (|V_i| - 1) = k + (-k) + \sum_{i=1}^{k} |V_i| = \sum_{i=1}^{k} |V_i|$$



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= $|V| - 1$,

thus establishing the claim also for T = (V, E).



Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

Proof of Theorem 262 for rooted trees: We use structural induction relative to proper subtrees. Obviously, the claim holds for the minimal elements, i.e., for trees that contain no proper subtrees and, thus, have only a root and no edges.

Now consider an arbitrary but fixed rooted tree $\mathcal{T}:=(V,E)$ and suppose that the equality claimed holds for all its k>0 proper subtrees $(V_1,E_1),\ldots,(V_k,E_k)$. (We do not need to assume explicitly that it holds for all subtrees of \mathcal{T} .) We get

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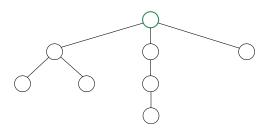
thus establishing the claim also for T = (V, E).

Corollary 263

If |V| > 1 holds for a (rooted) tree $\mathcal{T} := (V, E)$, then \mathcal{T} has at least one leaf.

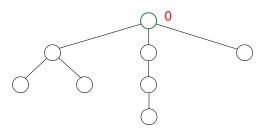
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Definition 264 (Depth, Dt.: Tiefe)



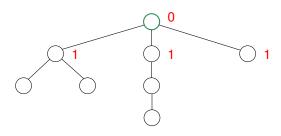


Definition 264 (Depth, Dt.: Tiefe)



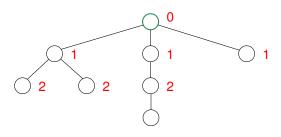


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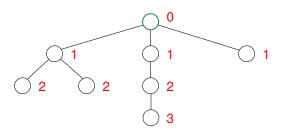


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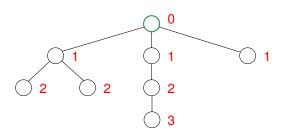


Definition 264 (Depth, Dt.: Tiefe)

The *depth* of the root u of a rooted tree $\mathcal{T}:=(V,E)$ is 0, and the depth of a node $v\neq u$ of \mathcal{T} is k if the depth of the parent of v is k-1, for all $v\in V$.

Warning

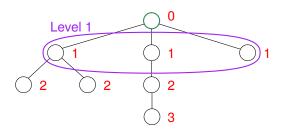
Some authors prefer to regard the root as a node at depth 1. Hence, make sure to check how depth is defined in a textbook prior to using the results stated!





Definition 265 (Level, Dt.: Niveau)

A *level* of a rooted tree $\mathcal T$ comprises all nodes of $\mathcal T$ which have the same depth.





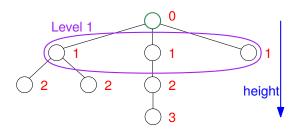
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Definition 265 (Level, Dt.: Niveau)

A *level* of a rooted tree \mathcal{T} comprises all nodes of \mathcal{T} which have the same depth.

Definition 266 (Height, Dt.: Höhe)

The *height* of a rooted tree \mathcal{T} is the maximum depth of nodes of \mathcal{T} .





Definition 267 (Binary tree, Dt.: Binärbaum)

A binary tree is an ordered tree \mathcal{T} with a root node u and at most two proper subtrees that are called left subtree, L, and right subtree, R. If \mathcal{T} has a left (right, resp.) subtree then L (R, resp.) is in turn a binary tree rooted in the left (right, resp.) child of u.



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Definition 268 (Complete binary tree, Dt.: vollständiger Binärbaum)

A *complete binary tree* is a binary tree in which every level, except possibly the last level, is completely filled, and the last level is filled from left to right.



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• E.g., a (binary) heap is a complete binary tree.



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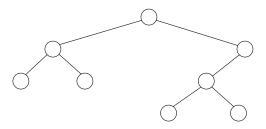
• E.g., a (binary) heap is a complete binary tree.

Definition 269 (Perfect binary tree, Dt.: perfekter Binärbaum)

A *perfect binary tree* is a binary tree that has the maximum number of nodes (relative to its height).



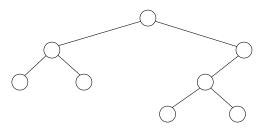
Definition 270 (Binary search tree, Dt.: binärer Suchbaum)





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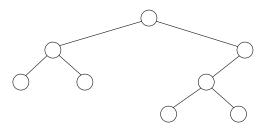
- if it has a proper left subtree L then
 - all values of nodes in *L* are less than the root value,





Definition 270 (Binary search tree, Dt.: binärer Suchbaum)

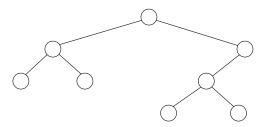
- if it has a proper left subtree L then
 - all values of nodes in L are less than the root value,
 - 2 L is a binary search tree itself,





Definition 270 (Binary search tree, Dt.: binärer Suchbaum)

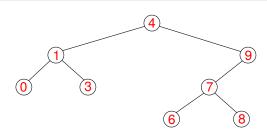
- if it has a proper left subtree *L* then
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- if it has a proper right subtree R then
 - all values of nodes in R are greater than the root value,
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Definition 270 (Binary search tree, Dt.: binärer Suchbaum)

- if it has a proper left subtree L then
 - all values of nodes in *L* are less than the root value,
 - L is a binary search tree itself,
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 - 3 all values of nodes in R are greater than the root value,
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Definition 271 (k-balanced tree, Dt.: k-balanzierter Baum)

A binary tree is height-balanced with balance factor k if it either has no proper subtrees



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A binary tree is height-balanced with balance factor k if it either has no proper subtrees or if

 it has two proper subtrees and the heights of both subtrees differ by not more than k,



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and if

- \odot all proper subtrees are height-balanced with balance factor k.
 - E.g., for k := 1: AVL tree.
 - Trees with balance factor 1 are simply called balanced or self-balancing.



Definition 272 (Perfectly balanced binary tree, Dt.: perfekt balanz. Binärbaum)

A binary tree $\mathcal T$ is *perfectly balanced* if all inner nodes of $\mathcal T$, except possibly on the second-last level, have exactly two children.



Definition 272 (Perfectly balanced binary tree, Dt.: perfekt balanz. Binärbaum)

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Lemma 273

A complete binary tree is perfectly balanced.



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• E.g., a (binary) heap is a perfectly balanced binary tree.

Lemma 273

A complete binary tree is perfectly balanced.

Lemma 274

A perfectly balanced binary tree has leaves only at its two bottom-most levels.



Lemma 275

For $i \in \mathbb{N}_0$, level i of a binary tree contains at most 2^i nodes.



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Lemma 275

For $i \in \mathbb{N}_0$, level i of a binary tree contains at most 2^i nodes.

Sketch of proof by induction: The claim holds for i := 0. If we have at most 2^k nodes on level k then we have at most $2 \cdot 2^k = 2^{k+1}$ nodes on level k+1.



Lemma 275

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Sketch of proof by induction: The claim holds for i := 0. If we have at most 2^k nodes on level k then we have at most $2 \cdot 2^k = 2^{k+1}$ nodes on level k + 1.

Lemma 276

Let h be the height and n be the number of nodes of a binary tree. Then $h \ge \lceil \log(n+1) \rceil - 1$, i.e., $h \in \Omega(\log n)$.



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Proof: Lemma 275 implies that a binary tree with height h contains at most

$$\sum_{i=0}^{h} 2^{i} = 2^{h+1} - 1$$

nodes.



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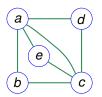
Theorem 277

If \mathcal{T} is a balanced binary tree with n nodes and height h then $h \in \Theta(\log n)$.

Definition 278 (Spanning tree, Dt.: spannender Baum)

A spanning tree of a connected graph ${\mathcal G}$ is a subgraph of ${\mathcal G}$ that

- is a tree,
- 2 includes all vertices of \mathcal{G} .

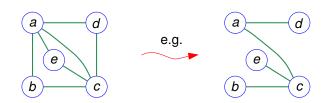




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A spanning tree of a connected graph ${\mathcal G}$ is a subgraph of ${\mathcal G}$ that

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- \bigcirc includes all vertices of \mathcal{G} .





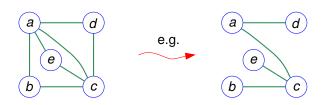
Definition 278 (Spanning tree, Dt.: spannender Baum)

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Theorem 279

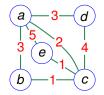
Every connected graph ${\cal G}$ contains a spanning tree.





Definition 280 (Weighted graph, Dt.: gewichteter Graph)

An *(edge-)weighted graph* is a graph in which every edge is assigned a (non-negative) real number, the so-called *weight* or *cost*.



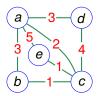


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Definition 281 (Minimum spanning tree, Dt.: minimal spannender Baum)

A minimum spanning tree (MST) of a weighted connected graph $\mathcal G$ is a spanning tree $\mathcal T$ of $\mathcal G$ such that the sum of the weights of the edges of $\mathcal T$ is minimum over all spanning trees of $\mathcal G$.





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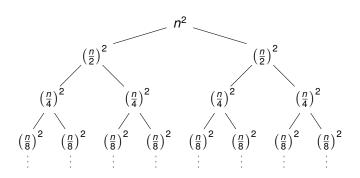
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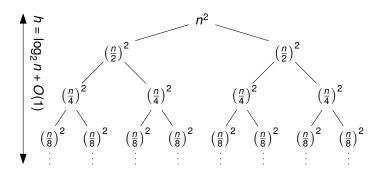


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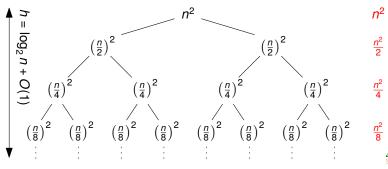


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- E.g., consider $T(n) = 2T(n/2) + n^2$. We get the following recursion tree with height $h = \log_2 n + O(1)$.

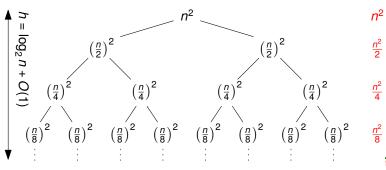




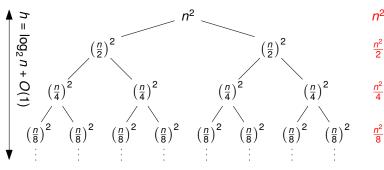
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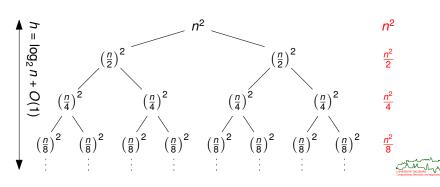
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- Summing across every level gives the total work done per level.
- Summing over all levels yields T(n): This is a geometric series, with $T \in \Theta(n^2)$.
- Master Theorem 215: We have a = b = k = 2 and, thus, $a < b^k$.



 Note that in this case the height of the tree does not really matter: The amount of work done at every level decreases so quickly that the total work is only a constant factor more than the work done at the root.

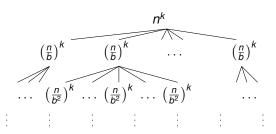


• For the recurrence relation $T(n) = a \cdot T\left(\frac{n}{b}\right) + n^k$ we get an a-ary recursion tree:



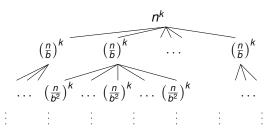
351/406

- For the recurrence relation $T(n) = a \cdot T\left(\frac{n}{b}\right) + n^k$ we get an a-ary recursion tree:
 - The problem size at level i is n/b^i .



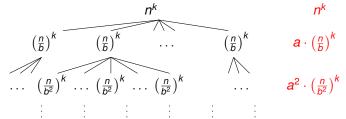


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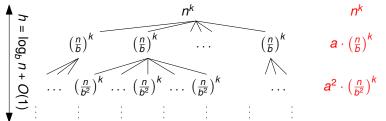




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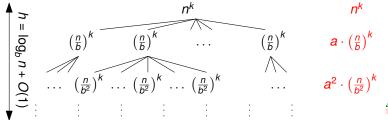


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 - The tree has $\log_b n + O(1)$ levels, i.e., a height of $O(\log n)$.

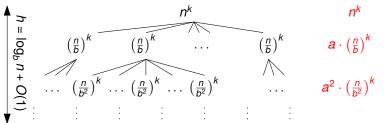




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 - The total number of leaves is $a^{\log_b n} = n^{\log_b a}$. (Recall $\log_b x = \log_a x / \log_a b$.)

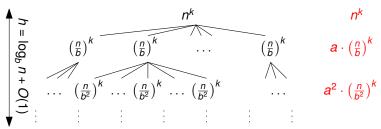


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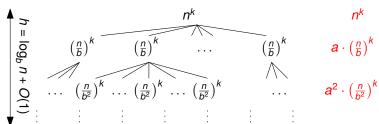
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 - The work done is constant per leaf.
 - Total work:

$$T(n) = \sum_{0 \leqslant i < \log_b n} a^i \cdot \left(\frac{n}{b^i}\right)^k + O(n^{\log_b a}) = \sum_{0 \leqslant i < \log_b n} n^k \cdot \left(\frac{a}{b^k}\right)^i + O(n^{\log_b a}).$$



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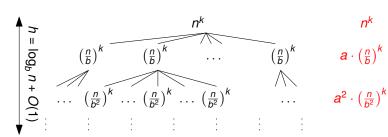


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• If $a = b^k$, i.e., if $k = \log_b a$, then

$$n^k \cdot \left(\frac{a}{b^k}\right)^i = n^k = n^{\log_b a}.$$





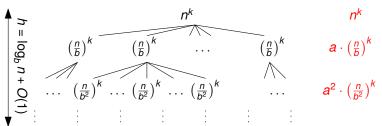
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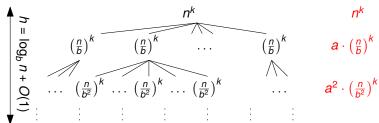
$$n^k \cdot \left(\frac{a}{b^k}\right)^i = n^k = n^{\log_b a}.$$

 Hence, the same order of work is done on every level, and since the tree has $O(\log n)$ levels, we get $T \in \Theta(n^{\log_b a} \log n)$; recall Thm. 215.



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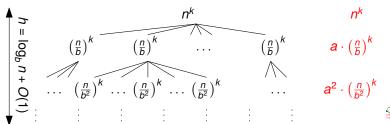




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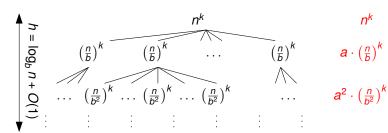
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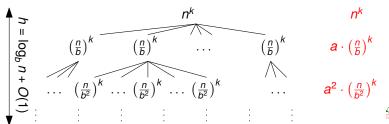




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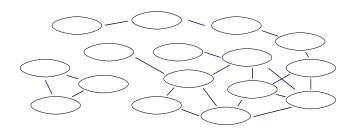
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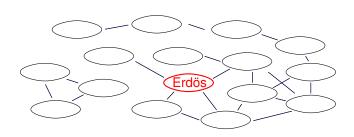


• A *collaboration graph* for a set of *n* scientists is a graph with *n* vertices such that two vertices are connected by an edge if the corresponding scientists have at least one joint publication.



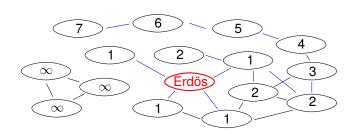


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- The *Erdös number* of a scientist is the "collaborative distance" of a scientist to the extremely prolific Hungarian mathematician Paul Erdös (1913–1996, more than 500 co-authors and more than 1500 publications): Erdös has 0, and a scientist has Erdös number k+1 if k is the lowest Erdös number of his/her co-authors.



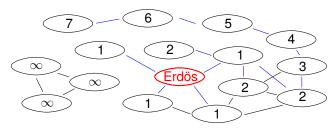


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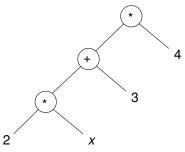
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- One's Erdös number can be obtained by computing minimum-weight paths on a collaboration graph.





Real-World Application: Algebraic Expression Trees

- An algebraic expression tree is a rooted tree that corresponds to an expression.
- E.g., an in-order traversal of the tree

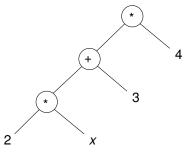


produces the standard (infix) expression $(2x + 3) \cdot 4$.



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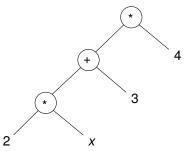
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A post-order traversal yields the postfix expression 2 x ⋅ 3 + 4 ⋅ , while a pre-order traversal yields the prefix expression ⋅ (+(⋅(2 x) 3) 4).



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- The analysis of expression trees is a central task for the simplification and parallel evaluation of an expression.



- Graph Theory
 - What is a (Directed) Graph?
 - Paths
 - Trees
 - Special Graphs
 - Complete and Bipartite Graphs
 - Hypercube
 - Isomorphic Graphs
 - Planar Graphs
 - Graph Coloring



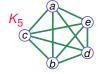
Definition 282 (Complete graph, Dt.: vollständiger Graph)

For $n \in \mathbb{N}$, the *complete graph* on n vertices, commonly denoted by K_n , is an undirected graph with n vertices in which every pair of vertices is adjacent.



Definition 282 (Complete graph, Dt.: vollständiger Graph)

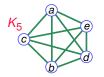
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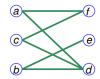




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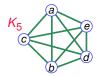
Definition 283 (Bipartite graph, Dt.: bipartiter Graph)

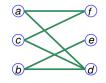
An undirected graph $\mathcal{G}:=(V,E)$ is a *bipartite graph* if V can be partitioned into two (non-empty) subsets V_1,V_2 such that $E\subseteq\{\{v_1,v_2\}:v_1\in V_1,v_2\in V_2\}.$

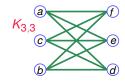


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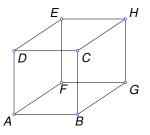
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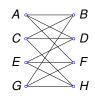
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Definition 284 (Complete bipartite graph, Dt.: vollständig-bipartiter Graph)

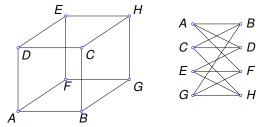
An undirected graph $\mathcal{G} := (V, E)$ is a *complete bipartite graph* if V can be partitioned into two (non-empty) subsets V_1, V_2 such that $E = \{\{v_1, v_2\} : v_1 \in V_1, v_2 \in V_2\}$. If $n := |V_1|$ and $m := |V_2|$ then \mathcal{G} is denoted by $K_{n,m}$.

• The edges and corners of a cube can be interpreted as a bipartite graph.

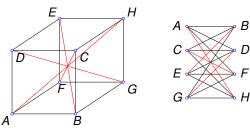




• The edges and corners of a cube can be interpreted as a bipartite graph.



• If we add all diagonals that cross the cube then we get $K_{4,4}$.





Lemma 285

Let $\mathcal{G} := (V, E)$ be a bipartite graph and let V_1, V_2 be the partition of V according to Def. 283. Then

$$\sum_{v_1\in V_1} \text{deg}(v_1) = \sum_{v_2\in V_2} \text{deg}(v_2) = |E|.$$



Lemma 285

Let $\mathcal{G}:=(V,E)$ be a bipartite graph and let V_1,V_2 be the partition of V according to Def. 283. Then

$$\sum_{\nu_1\in V_1} \text{deg}(\nu_1) = \sum_{\nu_2\in V_2} \text{deg}(\nu_2) = |E|.$$

Proof:

ullet As each edge has exactly one vertex from V_1 , we can write

$$\sum_{v_1\in V_1}\deg(v_1)=|E|.$$

Similarly,

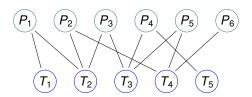
$$\sum_{v_2 \in V_2} \deg(v_2) = |E|.$$



 Suppose that we are given a set of tasks and a set of processors. We know which processor can carry out which tasks.

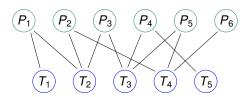


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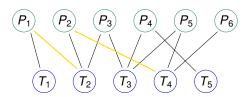




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Definition 286 (Matching, Dt.: Paarung)

• A matching in a simple graph $\mathcal{G} := (V, E)$ is a subset E' of E such that no two edges of E' are incident upon the same vertex of V.

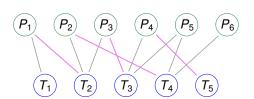




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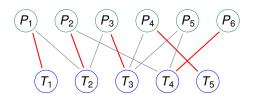




- Suppose that we are given a set of tasks and a set of processors. We know which processor can carry out which tasks.
- These relations can be represented as a bipartite graph.
- How can we get the maximum number of tasks processed concurrently?

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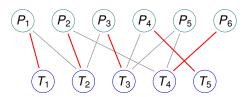




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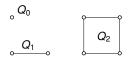


Hypercube

Definition 287 (Hypercube)

For $n \in \mathbb{N}_0$, the hypercube Q_n is defined recursively as follows:

- \bigcirc Q_0 is a single vertex;
- 2 Q_{n+1} is obtained by taking two disjoint copies of Q_n and linking each vertex in one copy of Q_n to the corresponding vertex in the other copy of Q_n .



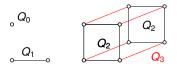


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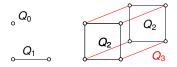


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For $n \in \mathbb{N}_0$, the hypercube Q_n is a regular graph of degree n with 2^n vertices and $n \cdot 2^{n-1}$ edges; it is bipartite for $n \ge 1$.

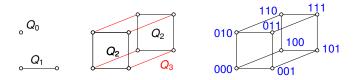


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For $n \in \mathbb{N}_0$, the hypercube Q_n is a regular graph of degree n with 2^n vertices and $n \cdot 2^{n-1}$ edges; it is bipartite for $n \ge 1$.

• We could also obtain Q_n by labeling 2^n vertices with distinct n-bit binary strings, and by connecting those vertices by edges whose strings differ in exactly one bit.

Real-World Application: Hamilton Cycles in Q_n Yield Gray Codes

Definition 289 (Gray code)

A (cyclic) Gray code of an ordered sequence of 2^n entities, for $n \in \mathbb{N}$, is a sequence of n-bit binary strings such that the encodings of two neighboring entities have Hamming distance one, i.e., differ by exactly one bit.

 Gray codes are widely used in position encoders and for error detection and correction in digital communication.



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Real-World Application: Hamilton Cycles in Q_n Yield Gray Codes

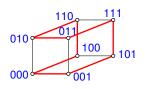
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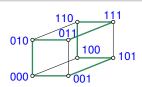
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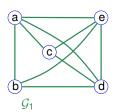
Definition 291 (Isomorphic, Dt.: isomorph)

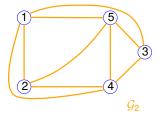
Two (directed) graphs $\mathcal{G}_1 = (V_1, E_1)$ and $\mathcal{G}_2 = (V_2, E_2)$ are *isomorphic*, denoted by $\mathcal{G}_1 \simeq \mathcal{G}_2$, if there exists a one-to-one mapping f between V_1 and V_2 that preserves adjacency; i.e., $uv \in E_1 \Leftrightarrow f(u)f(v) \in E_2$ for all $u, v \in V_1$. Such a suitable function f is called *graph isomorphism*.



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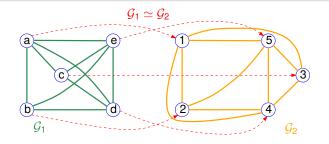






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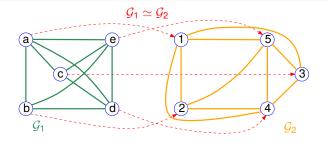
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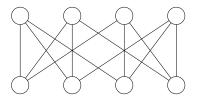


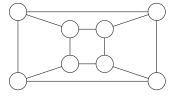
Lemma 292

The relation \simeq is an equivalence relation on graphs.

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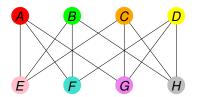
 Don't be fooled by drawings! Two graphs may be isomorphic even if their drawings look strikingly different.

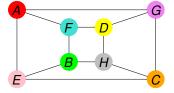






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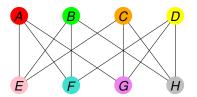


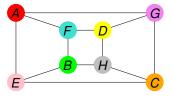




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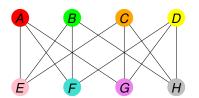


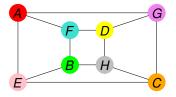


 Necessary (but not sufficient) conditions for two graphs to be isomorphic: same numbers of vertices and edges, same degrees.



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- Necessary (but not sufficient) conditions for two graphs to be isomorphic: same numbers of vertices and edges, same degrees.
- The complexity of the graph isomorphism problem for general n-vertex graphs is unknown. No polynomial-time algorithm is known, but the problem is also not known to be \mathcal{NP} -complete. In December 2015, Babai announced a deterministic algorithm that runs in time $2^{O(\log^c n)}$ time for some positive constant c, i.e., in quasi-polynomial time. In 2017, Helfgott claimed that one can take c:=3.
- Practically efficient algorithms for graph canonical labeling are known, though.



Real-World Application: Non-Isomorphic Trees Represent Molecules

- [Cayley 1857]: Molecules can be represented as graphs, where atoms are represented by vertices and bonds are represented by edges.
- Saturated hydrocarbons, C_nH_{2n+2} , are given by trees where each carbon atom is represented by a degree-four vertex and each hydrogen atom is a leaf.



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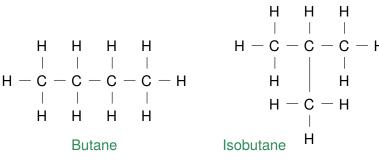
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- How many different isomers can exist for n := 4?
- We have exactly two non-isomorphic trees of this type and, thus, two different isomers of C₄H₁₀, namely butane and isobutane.



Planar Graphs

Definition 293 (Planar graph, Dt.: planarer oder plättbarer Graph)

A planar graph is a graph which can be drawn in the plane without edge crossings. A suitable drawing is called a (planar) embedding (Dt.: planare Einbettung).



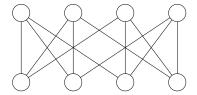
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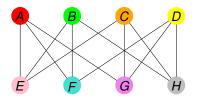


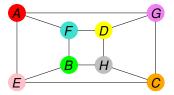


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- VLSI circuits are easier/cheaper to manufacture if their connections live in fewer layers.
- A scheme for a planetary gearset is compatible if and only if a suitably designed graph is planar.



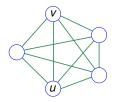
[Image credit: Rohloff AG, http://www.rohloff.de/]



Subdivision of a Graph

Definition 294 (Subdivision, Dt.: Unterteilung)

An *edge subdivision* of the edge $uv \in E$ by means of the vertex $w \notin V$ transforms the graph $\mathcal{G} := (V, E)$ into the graph $\mathcal{G}' = (V', E')$, where $V' = V \cup \{w\}$ and $E' = (E \setminus \{uv\}) \cup \{uw, wv\}$.

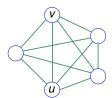


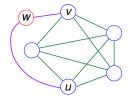


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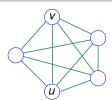
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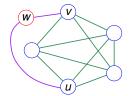
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Definition 295 (Subdivision graph, Dt.: Unterteilungsgraph)

A graph \mathcal{G}' is a *subdivision graph* of \mathcal{G} if \mathcal{G}' is obtained from \mathcal{G} via a finite sequence of edge subdivisions.







Theorem 296 (Kuratowski (1930))

A graph is planar if and only if it does not contain a subgraph that is isomorphic to a subdivision graph of K_5 or $K_{3,3}$.



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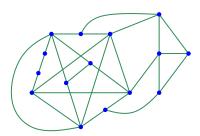
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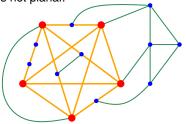
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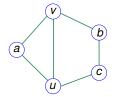
• Is the following graph planar? No: It contains a subdivision graph of K_5 as a subgraph. Hence, it is not planar.





Definition 298 (Edge contraction, Dt.: Kantenkontraktion)

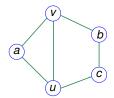
In a graph $\mathcal{G}:=(V,E)$, the *contraction* of an edge $e\in E$, with e=uv for some $u,v\in V$, replaces u and v by a new vertex $w\notin V$ such that edges incident to w are all edges other than e that were incident with u or v. All other nodes and edges are preserved.





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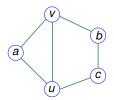


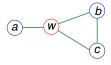




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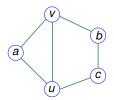


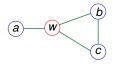




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A graph is planar if and only if it does not contain a subgraph that can be contracted to K_5 or $K_{3,3}$ via a finite sequence of edge contractions.



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Theorem 300 (Hopcroft&Tarjan (1974))

Testing whether a given graph with n vertices is planar can be done in O(n) time.



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Theorem 301 (Wagner (1936), Fáry (1948), Stein (1951))

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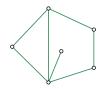


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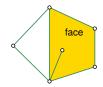
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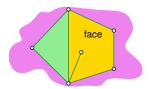
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 Note that one of the faces of a planar subdivision is unbounded: outer face.



Theorem 303 (Euler, Dt.: Eulerscher Polyedersatz)

Consider a planar subdivision induced by a connected planar graph \mathcal{G} . We denote

- the number of its vertices by v,
- the number of its edges by e,
- the number of its faces by *f*.

Then

$$v-e+f=2.$$



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Consider a planar subdivision induced by a connected planar graph $\mathcal{G}.$ We denote

- the number of its vertices by v,
- the number of its edges by e,
- the number of its faces by *f*.

Then

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• Euler's Formula generalizes to v-e+f=1+c for a planar graph with cconnected components.



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Corollary 304

Let v, e, f for a connected planar graph \mathcal{G} as defined in Theorem 303. If $v \ge 3$ then

$$e \leqslant 3v - 6$$
 and $f \leqslant 2v - 4$ and $f \leqslant \frac{2}{3}e$.



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We conclude $3f \leq 2e$.



Corollary 305

 K_5 is not planar.



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Proof: We get
$$v = 5$$
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Definition 306 (Triangle-free, Dt.: dreiecksfrei)

A triangle-free graph is a graph which does not contain a cycle of length three, i.e., in which no three vertices form a triangle of edges.



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Euler's Formula for Planar Graphs

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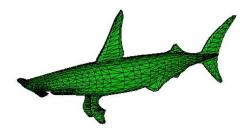
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Corollary 308

 $K_{3,3}$ is not planar.

Proof: $K_{3,3}$ is triangle-free and has six vertices and nine edges. If it were planar then, by Cor. 307, it could have at most $2 \cdot 6 - 4 = 8$ edges. Thus, $K_{3,3}$ is non-planar

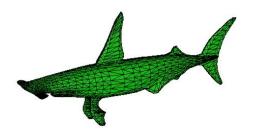
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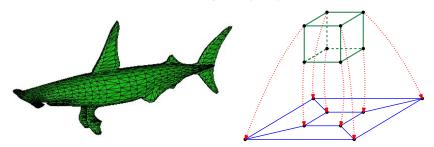


Theorem 309

The vertices and edges of a simple (bounded) polyhedron form a planar graph.



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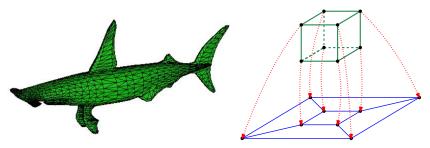


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Corollary 310

A simple (bounded) polyhedron with n vertices has at most 3n-6 edges and 2n-4 faces.

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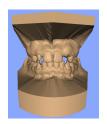
Real-World Application: Reducing the Face Count

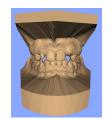
- Recent improvements in laser rangefinder technology allow the digitization of the shapes of physical objects at extremely high resolutions.
- The resulting polyhedral models are huge: E.g., a 0.25 mm model of Michelangelo's 5-meter statue of David contains about 1 billion polygonal faces!



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- Goal of multi-resolution modeling and level-of-detail modeling: Reduce the face count without sacrificing the visual appearance.
- E.g., the left dental model has 424 376 faces, while the other two models have only a few thousand faces.





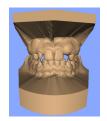


[Image credit: Michael Garland, Eurographics'99 STAR]



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- Edge contraction is one of the techniques used for reducing the face count.







[Image credit: Michael Garland, Eurographics'99 STAR]



- Graph Theory
 - What is a (Directed) Graph?
 - Paths
 - Trees
 - Special Graphs
 - Graph Coloring



Definition 311 (Coloring, Dt.: Färbung)

An assignment of colors to all vertices of a graph \mathcal{G} is called a *(vertex) coloring* if adjacent vertices are assigned different colors.



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Lemma 314

The chromatic number of a graph \mathcal{G} is two if and only if \mathcal{G} is bipartite.



• It is straightforward that every planar graph can be colored by six colors and that every triangle-free planar graph can be colored by four colors.



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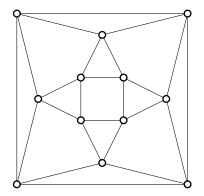
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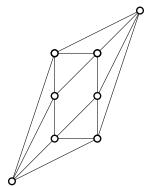
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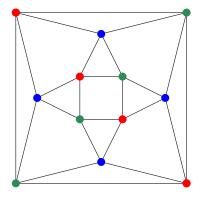
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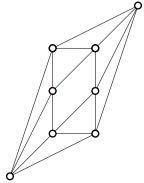
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- In 1996, Robertson et alii reduced the number of computer-checked cases to 633.
- In 2005, Werner and Gonthier used a general-purpose proof assistant ("Cog") to prove the theorem.



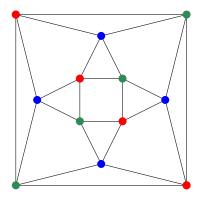


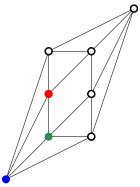


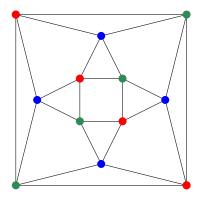


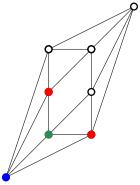




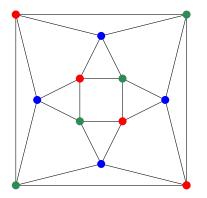


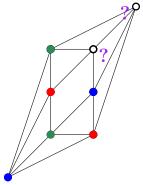






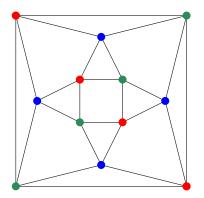


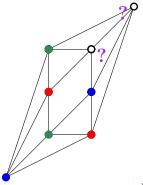




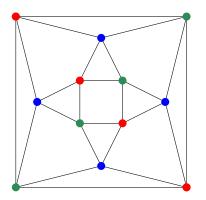


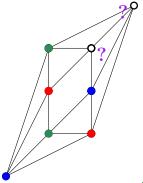
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- However, fairly efficient heuristics exist for approximate graph coloring.







[Image credit: Wikipedia]

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Corollary 316

If every entity of a topographic map is a connected area then four colors suffice to color the map such that no two entities that share a common border (other than a common point) are colored with the same color.





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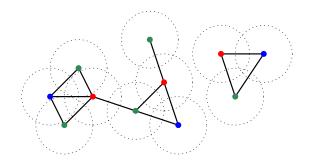
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• Note that this result holds only in the plane! E.g., on the surface of a torus seven colors are sufficient and may be necessary.

Real-World Application: Channel Assignment

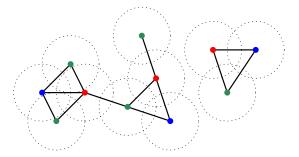
- We can solve the channel assignment problem by considering its so-called unit-disk graph, where
 - the vertices are given by the broadcast stations,
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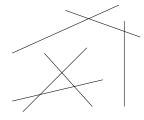
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- Obviously, the chromatic number of that graph equals the minimum number of frequencies needed.





Real-World Application: Minimum Plane Partition

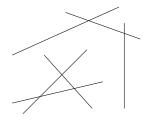
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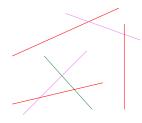
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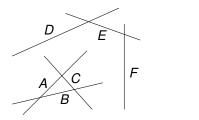
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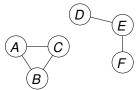




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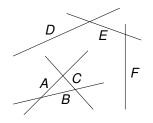


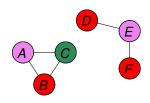




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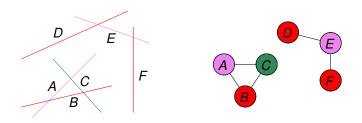






Real-World Application: Minimum Plane Partition

- CG:SHOP Geometric Optimization Challenge 2022: Given is a set *S* of line segments in the plane.
- We seek a partitioning of S into a minimum number of k subsets S_1, \ldots, S_k such that, for all $1 \le i \le k$, the line segments of S_i do not intersect pairwise.
- An obvious attempt to solve this problem is to construct the conflict graph G for S
 and then apply graph coloring to G.





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- How many index registers are needed for a given loop?



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- Other applications of graph coloring:
 - Scheduling consumer-producer interactions to allow concurrency.
 - Sudoku puzzles.



- - Cryptography
 - Introduction
 - Symmetric-Key Cryptography
 - Public-Key Cryptography



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- A sender A ("Alice") sends an encoded message to a receiver B ("Bob").
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Two main schemes in use nowadays

Symmetric-Key Cryptography (SKC): The same secret key is used for both encryption and decryption; aka secret-key cryptography.

Public-Key Cryptography (PKC): Different keys are used for encryption and decryption, with some keys being known publicly; aka asymmetric-key cryptography.

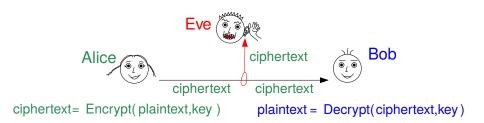


Basic Terms

- Plaintext original message.
- Ciphertext encoded/encrypted message.
- Encryption generating ciphertext from plaintext.
- Decryption / Deciphering generating plaintext from ciphertext.
- Cryptanalysis trying to break the encryption by applying various methods.
- Adversary, Spy the message thief.
- Eavesdropper a secret listener who listens to private conversations.
- Authentication the process of proving one's identity.
- Privacy ensuring that the message is read only by the intended receiver. (GnuPG: "Privacy is not a crime!")

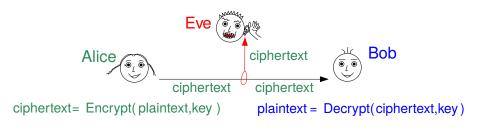


Eavesdropper Attacks





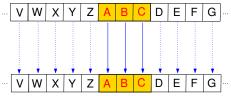
Eavesdropper Attacks



- Eve might attempt to
 - break the encryption,
 - replay the encrypted message (e.g., login) without breaking the encryption,
 - modify the message,
 - block the message,
 - fabricate a new message.

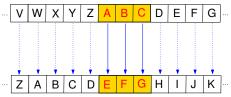


 According to Suetonius, Caesar (100–44 BCE) used an encryption scheme (for communication with his generals) that shifted the alphabet of the plaintext by some fixed position value n.



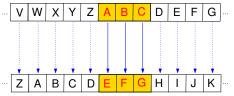


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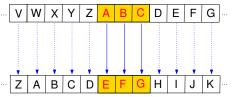


• With *n* := 4:

Plaintext: alea iacta est Ciphertext: epie megxe iwx



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With n := 4:

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- Suppose that the (Roman) letters are mapped to the numbers $0, 1, \ldots, 25$.
- Then Caesar's encryption and decryption with shift *n* can be computed as follows:

$$ciphertext := Encrypt_n(plaintext) = (plaintext + n) \mod 26$$

 $plaintext := Decrypt_n(ciphertext) = (ciphertext - n) \mod 26$



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- Nevertheless, Caesar's cipher with n := 13, aka ROT13, has been (mis-)used for serious applications even rather recently.
- However, it is used within more complex systems, e.g., the Vigenère cipher.
- On a Unix machine, the tr utility can be used for carrying out Caesar's cipher. E.g.,

```
echo "alea iacta est" | tr 'A-Za-z' 'E-ZA-De-za-d'
yields
     epie megxe iwx,
and
     echo "epie megxe iwx" | tr 'E-ZA-De-za-d' 'A-Za-z'
recovers the original text
     alea iacta est.
```



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Simple example: Suppose that Alice wants to encrypt a bit string A. Then Alice and Bob could choose a secret key B and apply a bit-wise XOR (exclusive OR, ⊕) — an output bit is 1 if exactly one of the two input bits is 1 — in order to transmit A ⊕ B. Then Bob would compute (A ⊕ B) ⊕ B and, thus, retrieve A.



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Α	В	$A \oplus B$	$(A \oplus B) \oplus B$
0	0	0	0
0	1	1	0
1	0	1	1
1	1	0	1



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- The key distribution problem is a major roadblock on the road to secure communication among folks who do not meet regularly.
- A second big disadvantage is the need for multiple keys in order to encrypt messages intended for different receivers.



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 - Diffie-Hellman Algorithm
 - RSA



Public-Key Cryptography

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 - exponentiation versus logarithms ("discrete log problem"):
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 - Again, finding a and b such that $\log_a 243 = b$ is considerably more difficult.



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- Diffie and his advisor Hellman were the first to publish a PKC scheme in 1976
 (They were the recipients of the 2015 ACM Turing Award.)

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- Diffie-Hellman key exchange is used by the Tor system to set-up secure communication links with onion routers.
- The Diffie-Hellman key exchange is vulnerable to man-in-the-middle attacks

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• Alice and Bob make p:=13 and g:=2 public. The number 2 is indeed a generator modulo 13 because the following powers of two taken modulo 13 yield the integers $1,2,\ldots,12$: $2^{12},2^1,2^4,2^2,2^9,2^5,2^{11},2^{15},2^8,2^{10},2^7,2^6$.



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- Finally, $T^s \mod p = 12^5 \mod 13 = (12140 \cdot 13 + 12) \mod 13 = 12$, and $S^t \mod p = 6^6 \mod 13 = (3588 \cdot 13 + 12) \mod 13 = 12$.
- Hence, Alice and Bob have managed to exchange 12 as a master key for their future communication.



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No toy numbers!

Of course, in practice considerably larger values are chosen for p!!

Lemma 317

Let $a, b \in \mathbb{N}$ such that gcd(a, b) = 1. Then there exists $x \in \mathbb{Z}$ such that $a \cdot x \equiv_b 1$.



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Proof: Since gcd(a,b) = 1, Cor. 125 tells us that there exist $x, y \in \mathbb{Z}$ such that $a \cdot x + b \cdot y = 1$. Hence, $a \cdot x = 1 - b \cdot y \equiv_b 1$.



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Definition 318 (Euler's Totient Function, Dt.: Eulersche φ -Funktion)

Euler's totient function $\varphi : \mathbb{N} \to \mathbb{N}$ is defined as

$$\varphi(n) := |U_n|, \text{ with } U_n := \{x \in \mathbb{N} \colon 1 \leqslant x \leqslant n \land \gcd(x, n) = 1\}.$$

The set U_n is called the *group of units* of n.



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- Hence, $\varphi(n)$ is the number of integers among 1, 2, ..., n that are coprime to n.
- We have $\varphi(4) = 2$, $\varphi(5) = 4$, $\varphi(6) = 2$.
- More generally, $\varphi(p) = p 1$ for every $p \in \mathbb{P}$.



Lemma 319

Let
$$p, q \in \mathbb{P}$$
. If $p \neq q$ then $\varphi(pq) = (p-1)(q-1)$.



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Proof: There are q multiples of p and p multiples of q within $\{1, 2, \ldots, pq\}$, and the only common multiple of both p and q is pq. Hence, by the Inclusion-Exclusion Principle (Thm. 167), $\varphi(pq) = pq - p - q + 1 = (p-1)(q-1)$.



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Let $n \in \mathbb{N}$ and $x \in U_n$. Then $x^{\varphi(n)} \equiv_n 1$.

Corollary 321

Let $n \in \mathbb{N}$ and $x \in U_n$. If n = pq, with $p, q \in \mathbb{P}$ and $p \neq q$, then $x^{(p-1)(q-1)} \equiv_n 1$.



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 - Select two distinct prime numbers p and q (each of which, in practice, has at least 150 digits) and compute $n = p \cdot q$.
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 - The numbers n and e are published (Bob's public key).
 - Compute a number d which is the inverse of e in $\mathbb{Z}_{\varphi(n)}$, i.e., such that $d \cdot e \equiv_{\varphi(n)} 1$. (Such a number exists due to Lem. 317.)
 - The number d is called Bob's private key and is kept secret.



• Hence, we have $n = p \cdot q$ and, thus, $\varphi(n) = (p-1) \cdot (q-1)$. Furthermore, $\gcd(e, \varphi(n)) = 1$ and $d \cdot e \equiv_{\varphi(n)} 1$.



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- Encoding the ciphertext:
 - Alice encodes a message $x \in \mathbb{N}$, with x < n to keep it in \mathbb{Z}_n and with gcd(x, n) = 1, by using Bob's public key e and n: $y := x^e \mod n$ with 0 < y < n.



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 - Alice encodes a message x ∈ N, with x < n to keep it in Zn and with gcd(x, n) = 1, by using Bob's public key e and n:
 y := xe mod n with 0 < y < n.
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 - Bob computes $z := y^d \mod n$ with 0 < z < n.



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• Suppose that p := 5 and q := 11. Hence n = 55 and $\varphi(n) = 40$. Suppose further that three users chose the following keys:

е	d
23	7
37	13
9	9
	23 37



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- Let x := 2 (and use Mathematica to do the arithmetic).

Alice:
$$2^{23} = 8388608 = 152520 \cdot 55 + 8 \equiv_{55} 8 =: y$$

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Bob:
$$2^{37} = 137438953472 = 2498890063 \cdot 55 + 7 \equiv_{55} 7 =: y$$

$$7^{13} = 96889010407 = 1761618371 \cdot 55 + 2 \equiv_{55} 2 =: z$$

Caesar:
$$2^9 = 512 = 9 \cdot 55 + 17 \equiv_{55} 17 =: y$$

 $17^9 = 118587876497 = 2156143209 \cdot 55 + 2 \equiv_{55} 2 =: z$



- Note that there are $\varphi(n)$ many messages that can be sent for n given.
- Since

$$\frac{\varphi(n)}{n} = \frac{(p-1)(q-1)}{pq} = (1 - \frac{1}{p})(1 - \frac{1}{q})$$

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- An eavesdropper who only knows n, e, and y cannot do much with this information. In particular, no efficient algorithm is known to factor n into p, q as a simple means to obtain $\varphi(n)$.
- It is also important to ensure that $x^e > n$, i.e., that y is obtained by exponentiation and then by a reduction modulo n.
 - If x^e < n then one could simply recover x by taking the e-th root of y. (After all, e is known publicly!)
 - Hence, it is wise to select e such that $2^e > n$.



The End!

I hope that you enjoyed this course, and I wish you all the best for your future studies.



