# Combinatorial polynomials 

University of Fribourg
Spring term 2020

Disclaimer: this is a draft-in-the-making, it will be updated during the semester.

Please report to me any correction or typos you might find - thank you for your cooperation!

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September 14, 2020

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## Chapter 0

## Some recurring notations

$\mathscr{P}(\mathrm{X})$ or $2^{\mathrm{X}}$ the set of all subsets of a set X .
$[n]$ The set $\{1,2, \ldots, n\}$. Set $[0]:=\emptyset$.
$B_{n}$ The boolean poset on $n$ elements, i.e., the set $\mathscr{P}([n])$ ordered by inclusion.
$\binom{X}{k}$, where $X$ is a finite set and $k \in \mathbb{N}$, is the set of all $k$-element subsets of $X$. If $X$ is finite and has, say, $n$ elements, then $\left|\binom{X}{k}\right|=\binom{n}{k}$
$B^{A}$, where $A, B$ are sets, denotes the set of all functions $A \rightarrow B$. If $A, B$ are finite, the cardinality of this set is $|B|^{|A|}$

## Chapter 1

## Graphs, colorings and flows


#### Abstract

Disclaimer. We will only briefly review the basics of graph theory that are strictly necessary for our purposes. For more background or further discussion of some topics the reader can refer to one of the many excellent dedicated textbooks. Our formal setup follows mostly Tutte's book [7], while we occasionally update some terminology in order to facilitate access to more modern literature such as [3, 4].


### 1.1 Graphs

Definition 1.1.1. A graph $G=(V, E, h, t)$ is a quadruple consisting of a set of vertices $V$, a set of Edges $E$ and two functions $h, t: E \rightarrow V$ that assign to every edge its "ends". Given any set $A \subseteq E$ of edges we let $V(A):=h(A) \cup t(A)$ be the set of all ends of edges in $A$.

We will often omit braces when designing one-element sets, if no need for specification arises. For instance, given $e \in E$ we will write $V(e)$ for $V(\{e\})$.

A loop in $G$ is any $e \in E$ with $|V(e)|=1$. Two edges $e, e^{\prime} \in E$ are called parallel if $\mathrm{V}(e)=\mathrm{V}\left(e^{\prime}\right)$. The graph G is called simple if it has no loops nor parallel edges. A trail in $G$ is any sequence $v_{0}, e_{1}, v_{1}, \ldots, e_{k}, v_{k}$ of vertices and edges such that $\left\{v_{i-1}, v_{i}\right\}=\mathrm{V}\left(e_{i}\right)$ for all $i=1, \ldots, k$. It is called closed (or a "cycle") if $k>0$ and $v_{0}=v_{\mathrm{k}}$. A path is a trail where all edges and all vertices are pairwise distinct (in this case we will talk about a "path from $v_{0}$ to $v_{\mathrm{k}}$ ". A circuit in $G$ is a minimal closed trail, i.e., a closed trail which, after removal of any edge, is a path (in particular, every loop is a circuit).


Figure 1: A trail, a path, a circuit and a closed trail that is not a circuit.

Let $\mathrm{T} \subseteq \mathrm{V}$ be a set of vertices of G the vertex-induced subgraph defined by T is the graph $G(T):=\left(T, E^{\prime}, h, t\right)$ where $E^{\prime}=\{e \in E \mid\{h(e), t(e)\} \subseteq T\}$ is the set of edges with both endpoints in T .
Definition 1.1.2. Let $G$ be a graph. We call $G$ connected if for any two vertices $v, w \in \mathrm{~V}$ there is a path from $v$ to $w$ in G . A connected component of G is any maximal connected vertex-induced subgraph, i.e., any vertex-induced subgraph $\mathrm{G}(\mathrm{T})$ that is connected and such that, for every $v \in \mathrm{~V} \backslash \mathrm{~T}, \mathrm{G}(\mathrm{T} \cup\{v\})$ is not connected. We define
$c(G):=$ the number of connected components of $G$.

(a)

(b)

Figure 2: Two connected graphs

Two fundamental operations on graphs are deletion and contraction of edges.

Definition 1.1.3 (Edge deletion). Let $G=(V, E, h, t)$ be a graph and let $A \subseteq E$. The deletion of $A$ from $G$ is the graph $G \backslash A:=\left(\mathrm{V}, \mathrm{E} \backslash A, \mathrm{~h}_{\mid \mathrm{E} \backslash A}, \mathrm{t}_{\mid \mathrm{E} \backslash \mathcal{A}}\right)$ on the same vertex set as $G$ but without the edges in $A$, and with the functions $h, t$ restricted accordingly. If $A=\{e\}$ consists of a single element, we sometimes write $\mathrm{G} \backslash e$ for $\mathrm{G} \backslash\{e\}$. The "restriction" of G to A is $\mathrm{G}[\mathrm{A}]:=\mathrm{G} \backslash(\mathrm{E} \backslash \mathcal{A})$.

Remark-Notation 1.1.4. An edge $e \in \mathrm{E}$ is called an isthmus if $\mathfrak{c}(\mathrm{G})<\mathfrak{c}(\mathrm{G} \backslash e)$.
Remark 1.1.5 (On the word "subgraph"). Every graph of the form $G[A]$ we will call a "subgraph" of G. Notice the difference with the notion of "vertex induced subgraph" discussed earlier on. The latter will not appear in the following, so we feel safe in our terminological choice.
Definition 1.1.6 (Edge contraction). Let $\mathrm{G}=(\mathrm{V}, \mathrm{E}, \mathrm{h}, \mathrm{t})$ be a graph and let $A \subseteq E$. The contraction of $A$ in $G$ is the graph $G / A:=\left(V^{\prime}, E^{\prime}, h^{\prime}, t^{\prime}\right)$ with edge set $\mathrm{E}^{\prime}:=\mathrm{E} \backslash \mathrm{A}$ and vertex set $\mathrm{V}^{\prime}:=\mathrm{V} / \sim$ given as the set of equivalence classes of the equivalence relation on V generated by $v \sim w$ if $\{v, w\}=\mathrm{V}(e)$ for some $e \in A$. The functions $h^{\prime}$ and $t^{\prime}$ are given by $h^{\prime}(e)=[h(e)]_{\sim}$, resp. $t^{\prime}(e)=[t(e)]$.

Intuitively, deleting $A$ from $G$ means removing every edge in $A$, while contracting $A$ means "shrinking" every edge in $A$ and identifying its endpointvertices.

(a)

(b)

Figure 3: Deletion (l.-h.s.) and contraction (r.-h.s.) of an edge in the graphs depicted in Figure 2(a), resp. Figure 2(b).

### 1.2 Colorings

Definition 1.2.1. Let $G=(V, E, h, t)$ be a graph. A coloring of $G$ is any function $\gamma: \mathrm{V} \rightarrow \mathbb{N}_{>0}$ such that $\gamma(\mathrm{h}(\mathrm{e})) \neq \gamma(\mathrm{t}(\mathrm{e}))$ for every $\mathrm{e} \in \mathrm{E}$. Given $\mathrm{k} \in \mathbb{N}$, we call $\gamma$ a k-coloring of G if $\gamma(\mathrm{V}) \leqslant \mathrm{k}$. The set of all k -colorings of a given graph G is denoted by $\Gamma_{\mathrm{G}}^{(\mathrm{k})}$.
Definition 1.2.2. Let $G=(V, E, h, t)$ be a finite graph. The chromatic function of $G$ is the function

$$
\chi_{\mathrm{G}}^{\mathbb{N}}: \mathbb{N} \rightarrow \mathbb{N}, \quad \mathrm{k} \mapsto\left|\Gamma_{\mathrm{G}}^{(\mathrm{k})}\right|
$$

## Example 1.2.3.

(1) If $G$ has a loop, then obviously $\chi_{G}^{\mathbb{N}}(k)=0$ for all $k$.
(2) If $e, e^{\prime}$ are two parallel edges of $G$, then $\chi_{G}^{\mathbb{N}}(k)=\chi_{G}^{\mathbb{N}} \backslash e(k)=\chi_{G}^{\mathbb{N}} \backslash e^{\prime}(k)$ for all $k$.
(3) If $G$ has no edges, i.e., $E=\emptyset$, then $\chi_{G}^{\mathbb{N}}(k)=k^{|V|}$ for all $k$.
(4) The complete graph on $n$ vertices is the simple graph $K_{n}$ on an $n$-element vertex set with one edge connecting any two distinct vertices. Then, $\chi_{K_{n}}^{\mathbb{N}}(k)=k!/(k-n)!$ for all $k$.

(a) $\chi_{G}^{\mathbb{N}}(k)=k(k-1)(k-2)^{2}$

(b) $\chi_{G}^{\mathbb{N}}(k)=k(k-1)(k-2)$

Figure 4: Two graphs with their chromatic functions.

Proposition 1.2.4. Let G be a finite graph, and let $\mathrm{e} \in \mathrm{E}$ be any edge of G that is not a loop. Then,

$$
\chi_{\mathrm{G}}^{\mathbb{N}}(\mathrm{k})=\chi_{\mathrm{G} \backslash e}^{\mathbb{N}}(\mathrm{k})-\chi_{\mathrm{G} / \mathrm{e}}^{\mathbb{N}}(\mathrm{k}) \text { for all } \mathrm{k} \in \mathbb{N}
$$

Together with Examples 1.2.4.(1) and (3), this determines the chromatic function uniquely.

Proof. We clearly have $\Gamma_{\mathrm{G}}^{(\mathrm{k})} \subseteq \Gamma_{\mathrm{G} \backslash e}^{(\mathrm{k})}$. Moreover every $\gamma \in \Gamma_{\mathrm{G} \backslash e}^{(\mathrm{k})} \backslash \Gamma_{\mathrm{G}}^{(\mathrm{k})}$ has $\gamma(\mathrm{h}(e))=\gamma(\mathrm{t}(e))$ and thus induces a valid coloring $\gamma^{\prime}$ of $\mathrm{G} / \mathrm{e}$ via $\gamma^{\prime}: \mathrm{V}^{\prime} \rightarrow \mathbb{N}$ with $\gamma^{\prime}\left([v]_{\sim}\right)=\gamma(v)$ for all $v$. Conversely, every $\gamma^{\prime} \in \Gamma_{\mathrm{G} / e}^{(\mathrm{k})}$ induces a unique $\gamma \in \Gamma_{\mathrm{G} \backslash e}^{(\mathrm{k})}$. Therefore $\Gamma_{\mathrm{G} \backslash e}^{(\mathrm{k})} \backslash \Gamma_{\mathrm{G}}^{(\mathrm{k})}$ and $\Gamma_{\mathrm{G} / e}^{(\mathrm{k})}$ have the same cardinality, and thus $\left|\Gamma_{\mathrm{G} \backslash e}^{(\mathrm{k})}\right|=\left|\Gamma_{\mathrm{G}}^{(\mathrm{k})}\right|+\left|\Gamma_{\mathrm{G} / \mathrm{e}}^{(\mathrm{k})}\right|$. The claim follows.
Corollary-Definition 1.2.5. For every given finite graph G there is a polynomial $\chi_{G} \in \mathbb{Z}[x]$ such that $\chi_{G}(k)=\chi_{G}^{\mathbb{N}}(\mathrm{k})$ for all $\mathrm{k} \in \mathbb{N}$. This $\chi_{G}$ is called the chromatic polynomial of G .

We are going to give what is called a "combinatorial interpretation" of the coefficients of the chromatic polynomial of a graph. This means that, although we know that we can write

$$
\begin{equation*}
\chi_{G}(t)=\sum_{i=0}^{n} a_{i} t^{i} \tag{1.1}
\end{equation*}
$$

for some integers $a_{i}$, we would like to be able to determine such coefficients. We present one of the earliest results in this vein, obtained by Hassler Whitney in 1934 [8]. We will follow Whitney's original argument, using an elementary but extremely versatile tool in enumerative combinatorics.

Principle of inclusion-exclusion. Let $\left(X_{i}\right)_{i \in I}$ be a family of (not necessarily distinct) subsets of a finite set $X$. Then

$$
\left|X \backslash\left(\bigcup_{i \in I} x_{i}\right)\right|=|X|-\sum_{i \in I}\left|X_{i}\right|+\sum_{\substack{i, j \in I \\ i \neq j}}\left|X_{i} \cap X_{j}\right|-\ldots \ldots+(-1)^{|I|}\left|\bigcap_{i \in I} x_{i}\right| \cdot
$$

Proposition 1.2.6. Let G be a graph with n vertices. Then the chromatic polynomial satisfies

$$
\begin{equation*}
X_{G}(t)=\sum_{A \subseteq E}(-1)^{|\mathcal{A}|} t^{c}(G[A]) . \tag{1.2}
\end{equation*}
$$

In particular, if we let $\mathrm{g}_{\mathrm{p}, \mathrm{q}}$ denote the number of $\mathrm{A} \subseteq \mathrm{E}$ such that $|\mathcal{A}|=\mathrm{p}$ and $\mathrm{c}(\mathrm{G}[\mathrm{A}])=\mathrm{q}$, then the coefficients of the expansion in Equation (1.1) can be written as $a_{i}=\sum_{j=0}^{\mathfrak{n}}(-1)^{j} g_{i, j}$ for all $i$.
Proof. Fix a positive integer k and let us call k -precoloring of G any function $\tilde{\gamma}: \mathrm{V} \rightarrow[\mathrm{k}]$, i.e., an arbitrary assignment of one of the k colors to the vertices of G . There are $\mathrm{k}^{|\mathrm{V}|}$ precolorings of G . Such a precoloring will sometimes not be a valid coloring of G , as there might be some "bad" edges of G joining
two vertices of the same color. Given $A \subseteq E$, we let $\operatorname{Pre}(A)$ be the set of all precolorings $\widetilde{\gamma}$ whose set of bad edges is $A$.
Formally, we set $\operatorname{Pre}(A):=\left\{\widetilde{\gamma}: V \rightarrow \mathbb{N}_{>0} \mid \widetilde{\gamma}(h(e))=\widetilde{\gamma}(t(e)) \forall e \in A\right\}$. Notice that
(1) $\operatorname{Pre}\left(A_{1}\right) \cap \operatorname{Pre}\left(A_{2}\right)=\operatorname{Pre}\left(A_{1} \cup A_{2}\right)$ for all $A_{1}, A_{2} \subseteq E$,
(2) $|\operatorname{Pre}(A)|=k^{|c(G[A])|}$.

The idea is now to count all precolorings of $G$ that do not have any bad edges, i.e., to enumerate the set

$$
\Gamma_{\mathrm{G}}(\mathrm{k})=\mathrm{k}^{\mathrm{V}} \backslash \bigcup_{\mathrm{e} \in \mathrm{E}} \operatorname{Pre}(\mathrm{e})
$$

We can apply the principle of inclusion-exclusion to our situation, and write

$$
\begin{align*}
\chi_{G}(k)=|\operatorname{Pre}(\emptyset)| & =k^{|V|}-\sum_{\emptyset \subseteq A \subseteq E}(-1)^{|A|}|\operatorname{Pre}(A)| \\
& =\sum_{A \subseteq E}(-1)^{|A|} k^{c(A)}, \tag{1.3}
\end{align*}
$$

where in the last equality we used that $|\mathrm{V}|=\mathrm{c}(\mathrm{G}[\emptyset])$, hence $\mathrm{k}^{|\mathrm{V}|}=(-1)^{\emptyset^{\mathrm{g}}} \mathrm{k}^{\mathrm{c}(\mathrm{G}[\emptyset])}$. The claim follows.

Definition 1.2.7. Let $G$ be a finite graph and fix any total order $<$ on its set $E$ of edges. A broken circuit of G is any set of edges that is obtained by deleting the maximal element of any circuit. More precisely, a broken circuit is any set of the form $C \backslash \max _{<} C$ where $C$ is a circuit of $G$.

A no-broken-circuit set (or just an "nbc-set") of G is any $A \subseteq E$ that does not contain any broken circuit. Let $\operatorname{nbc}(\mathrm{G})$ denote the family of all nbc-sets of $G$ and for every $i$ let $n b c_{i}(G)$ denote the number of nbc-sets of cardinality $i$.


Figure 5: The graph for Example 1.2 .8

Example 1.2.8. Consider the graph in Figure 5 Its chromatic polynomial has been computed in Figure 4 and can be expanded to

$$
\chi_{G}(t)=t(t-1)(t-2)^{2}=t^{4}-5 t^{3}+8 t^{2}-4 t
$$

We identify the set of edges with $\{1,2, \ldots, 5\}$ according to the numbering displayed in Figure 5. We consider on this set the natural total order of integer numbers.

This graph has three circuits, with edge-sets 124,345 and 1235 . Therefore, there are three broken circuits: $12,34,123$. The no-broken-circuit sets are then

| Size 0: | $\emptyset$ | $\mathrm{nbc}_{0}(\mathrm{G})=1$ |
| :--- | :--- | :--- |
| Size 1: | $1,2,3,4,5$ | $\operatorname{nbc}_{1}(\mathrm{G})=5$ |
| Size 2: | $13,14,15,23,24,25,3545$ | $\operatorname{nbc}_{2}(\mathrm{G})=8$ |
| Size 3: | $135,145,235,245$ | $\operatorname{nbc}_{3}(\mathrm{G})=4$ |
| Size 4: | none | $\operatorname{nbc}_{4}(\mathrm{G})=0$ |

Theorem 1.2.9. Let G be a graph on n vertices. Then the coefficients of the chromatic polynomial of G as written in Equation (1.1) satisfy

$$
\mathrm{a}_{\mathrm{n}-\mathrm{i}}=(-1)^{\mathrm{i}} \mathrm{nbc}_{\mathfrak{i}}(\mathrm{G})
$$

for all i.
Proof. Recall the total ordering $<$ on E . We extend this to a total ordering on the set of all broken circuits, with the property that $B<B^{\prime}$ implies max $<(B) \leqslant$ $\max _{<}\left(\mathrm{B}^{\prime}\right)$, for any two broken circuits $B, B^{\prime}$. $\left.{ }^{1}\right]$ We can now enumerate all broken circuits in this ordering as $\mathrm{B}_{1}, \mathrm{~B}_{2}, \ldots, \mathrm{~B}_{\mathrm{s}}$ and define a partition of the set of subsets of the edge set $E$ as

$$
\begin{gathered}
2^{\mathrm{E}}=\mathscr{E}_{1} \sqcup \ldots \sqcup \mathscr{E}_{\mathrm{s}+1}, \text { with } \\
\mathscr{E}_{1}:=\left\{A \subseteq E \mid \mathrm{B}_{1} \subseteq A\right\} \text { and } \mathscr{E}_{i}:=\left\{A \subseteq \mathrm{E} \mid \mathrm{B}_{\mathrm{i}} \subseteq A\right\} \backslash \mathscr{E}_{\mathrm{i}-1} \text { for } i>1,
\end{gathered}
$$

so the elements of $\mathscr{E}_{i}$ are all subsets that contain $B_{i}$ but none of the $B_{j}$ for $j<i$. Let us now consider the summands in Equation (1.3) subdivided according to the $\mathscr{E}_{\mathrm{i}} \mathrm{s}$ :

$$
\chi_{G}(k)=\left(\sum_{A \in \mathscr{E}_{1}}(-1)^{|\mathcal{A}|^{c} k^{\mathcal{c}(G[A])}}\right)+\ldots+\left(\sum_{A \in \mathscr{E}_{s}+1}(-1)^{|\mathcal{A}| k^{\mathcal{c}(G[A])}}\right)
$$

For every $\mathfrak{i}=1, \ldots, s$ now let $e_{i}$ be the edge that has been removed in order to form the broken circuit $B_{i}$. The ordering that we have chosen is such that

[^0]$e_{i} \notin B_{j}$ for all $\mathfrak{j}<i$ (in fact, $e_{i}>\max B_{i} \geqslant \max B_{\mathfrak{j}}$ ). Therefore, if we consider the bipartition
$$
\mathscr{E}_{i}=\underbrace{\left\{A \in \mathscr{E}_{i} \mid e_{i} \in A\right\}}_{=: \mathscr{E}_{i}^{+}} \sqcup \underbrace{\left\{A \in \mathscr{E}_{i} \mid e_{i} \notin A\right\}}_{=: \mathscr{E}_{i}^{-}},
$$
the function $\epsilon: \mathscr{E}_{i}^{+} \rightarrow \mathscr{E}_{i}^{-}, \epsilon(A):=A \backslash\left\{e_{i}\right\}$, is a bijection.


Figure 6: In the situation of Example 1.2.8 the construction of the proof of Theorem 1.2.9 would consider the broken circuits $\mathrm{B}_{1}=12, \mathrm{~B}_{2}=123, \mathrm{~B}_{3}=34$, leading to the following families of sets: $\mathscr{E}_{1}$ as shaded blue in the picture, $\mathscr{E}_{2}=\emptyset, \mathscr{E}_{3}$ as shaded orange in the picture, $\mathscr{E}_{4}=\operatorname{nbc}(\mathrm{G})$.

Notice that, since the ends of $e_{i}$ are already connected by the path $B_{i}$ in each $A \in \mathscr{E}_{i}^{-}$, by applying $\epsilon$ the number of connected components does not change: $c(G[\epsilon(A)])=c(G[A])$, but the cardinality decreases by one. Therefore, the contributions of the summands $A$ and $\epsilon(A)$ cancel out in the sum:

$$
\sum_{A \in \mathscr{E}_{i}}(-1)^{|A|} k^{c(G[A])}=\sum_{A \in \mathscr{E}_{i}^{+}}\left((-1)^{|A|} k^{\mathcal{c}(G[A])}+(-1)^{|\epsilon(A)|} k^{c(G[\epsilon(A)])}\right)=0
$$

Therefore, only the sum associated to $\mathscr{E}_{S}+1$ contributes to the expression of $\chi_{\mathrm{G}}(\mathrm{k})$. Now notice that in fact $\mathscr{E}_{\mathrm{S}+1}=\operatorname{nbc}(\mathrm{G})$. Moreover, no $A \in \mathscr{E}_{\mathrm{S}+1}$ contains any circuit of $G$, and thus $c(G[A])=|V|-|A|$ for such $A$. We can now rewrite

$$
\chi_{G}(k)=\sum_{A \in \mathscr{E}_{s+1}}(-1)^{|A|} k^{c(G[A])}=\sum_{i=0}^{n}(-1)^{i} n b c_{i}(G) k^{n-i}
$$

as desired.

### 1.3 Flows

Definition 1.3.1. Let $G=(V, E, h, t)$ be a graph. An integer flow on $G$ is any function $\mathrm{f}: \mathrm{E} \rightarrow \mathbb{Z}$ such that for every $v \in \mathrm{~V}$ we have

$$
\begin{equation*}
\sum_{e \in h^{-1}(v)} f(e)-\sum_{e \in \mathfrak{t}^{-1}(v)} f(e)=0 \tag{1.4}
\end{equation*}
$$

Let $k \in \mathbb{N} \backslash\{0\}$. We call $f$ an integer $k$-flow if $|f(e)|<k$ for all $e \in E$.
A modular k -flow is any $\mathrm{f}: \mathrm{E} \rightarrow \mathbb{Z}_{\mathrm{k}}$ satisfying Equation (1.4) understood as a congruence modulo $k$ (i.e., as an equation in the Abelian group $\mathbb{Z}_{k}$ ).

Any flow $f$ is called nowhere zero if $0 \notin f(E)$. The set of all nowhere zero integer k-flows, resp. nowhere zero modular k-flows, is denoted by $\Phi_{\mathrm{G}}^{(\mathrm{k})}$, $\operatorname{resp}$. $\bar{\Phi}_{\mathrm{G}}^{(\mathrm{k})}$.

## Remark 1.3.2.

- The definition of a flow on a graph can be given with values in an arbitrary Abelian group. Relevant examples include real-valued flows, of interest in optimization of networks, or complex-valued flows, arising from the analysis of electrical networks. In fact, Equation 1.4 is easily seen to correspond to Kirchoff's conservation condition for flows in networks.

Definition 1.3.3. Let $G$ be a finite graph. The integer flow function and the modular flow function of G are defined as

$$
\varphi_{\mathrm{G}}^{\mathbb{N}}: \mathbb{N} \rightarrow \mathbb{N}, \quad \mathrm{k} \mapsto\left|\Phi_{\mathrm{G}}^{(\mathrm{k})}\right|, \quad \bar{\varphi}_{\mathrm{G}}^{\mathbb{N}}: \mathbb{N} \rightarrow \mathbb{N}, \quad \mathrm{k} \mapsto\left|\bar{\Phi}_{\mathrm{G}}^{(\mathrm{k})}\right|
$$

and they "count" the number of nowhere-zero integer k-flows, resp. nowhere zero modular k-flows on G.

## Example 1.3.4.

(1) If $G$ has no edges, clearly $\bar{\varphi}_{G}^{\mathbb{N}}(k)=\varphi_{G}^{\mathbb{N}}(k)=1$ for every $k$
(2) If $e$ is a loop of $G$, then $\bar{\varphi}_{G}^{\mathbb{N}}(k)=(k-1) \bar{\varphi}_{G}^{\mathbb{N}} \backslash e(k)$.

Lemma 1.3.5. If G has an isthmus, then $\varphi_{\mathrm{G}}^{\mathbb{N}}(\mathrm{k})=\bar{\varphi}_{\mathrm{G}}^{\mathbb{N}}(\mathrm{k})=0$ for all k .
Proof. Let e be an isthmus of G and let f be a flow on G . Let K be the connected component of $G$ containing $e$, and let $S$ be the vertex set of one of the two connected components of $K$ that arise deleting $e$. Then, up to sign, $f(e)$ is the total flow $f(S, V \backslash S)$ and thus, by Exercise I.7, equals zero.

Theorem 1.3.6. Let G be a finite graph and let e be an edge of G . If e is an isthmus, then $\bar{\varphi}_{\mathrm{G}}^{\mathbb{N}} \equiv 0$. If e is not a loop nor an isthmus of G ,

$$
\begin{equation*}
\bar{\varphi}_{\mathrm{G}}^{\mathbb{N}}(\mathrm{k})=\bar{\varphi}_{\mathrm{G} / e}^{\mathbb{N}}(\mathrm{k})-\bar{\varphi}_{\mathrm{G} \backslash e}^{\mathbb{N}}(\mathrm{k}) \quad \text { for every } \mathrm{k} . \tag{1.5}
\end{equation*}
$$

Together with Lemma 1.3.5 and Example 1.3.6(1) this relation determines the function $\bar{\varphi}_{\mathrm{G}}^{\mathbb{N}}$ uniquely.

Proof. Let $v, w$ be the ends of $e$, and set

$$
\Sigma_{v}:=\sum_{\substack{e \in E \backslash e \\ h(e)=v}} f(e)-\sum_{\substack{e \in E \backslash e \\ t(e)=v}} f(e), \quad \Sigma_{w}:=\sum_{\substack{e \in E \backslash e \\ h(e)=w}} f(e)-\sum_{\substack{e \in E \backslash e \\ t(e)=w}} f(e) .
$$

Now any $f \in \bar{\Phi}_{G}^{(k)}$, viewed as a function on the edges of $G / e$, satisfies trivially Equation (1.4) at every vertex of $G / e$ other than the vertex $[v]_{\sim}$, i.e., the one formed merging $v$ and $w$. At this vertex, Equation 1.4 in $G / e$ is $\Sigma_{v}+\Sigma_{w}=0$ and this is trivially satisfied since f , being a flow on $\mathrm{G} \backslash e$, satisfies both $\Sigma_{v}=0$ and $\Sigma_{w}=0$. Thus, $\bar{\Phi}_{\mathrm{G} \backslash e}^{(\mathrm{k})} \subseteq \bar{\Phi}_{\mathrm{G} / e}^{(\mathrm{k})}$.

On the other hand, an $f \in \bar{\Phi}_{\mathrm{G} / e}^{(k)}$ that is not in $\bar{\Phi}_{\mathrm{G} \backslash e}^{(\mathrm{k})}$ is one for which $\Sigma_{v}=-\Sigma_{w} \neq 0$. This means that it extends uniquely to a nowhere-zero flow $f^{\prime} \in \bar{\Phi}_{G}^{(k)}$ by setting, for every $g \in E$,

$$
f^{\prime}(g):= \begin{cases}f(g) & g \neq e \\ \Sigma_{t(e)} & g=e\end{cases}
$$

It remains to check that, in fact, $\Sigma_{v}=-\Sigma_{w} \neq 0$ must hold also for every nowhere-zero flow on $G$, so that indeed there is a bijection between $\bar{\Phi}_{\mathrm{G} / e}^{(k)} \backslash \bar{\Phi}_{\mathrm{G} \backslash e}^{(\mathrm{k})}$ and $\bar{\Phi}_{\mathrm{G}}^{(\mathrm{k})}$. The claim follows now in view of the already proved containment $\bar{\Phi}_{\mathrm{G} \backslash e}^{(\mathrm{k})} \subseteq \bar{\Phi}_{\mathrm{G} / e}^{(\mathrm{k})}$.

Corollary-Definition 1.3.7. For every given finite graph G there is a polynomial $\bar{\varphi}_{\mathrm{G}} \in \mathbb{Z}[\mathrm{x}]$ such that $\bar{\varphi}_{\mathrm{G}}(\mathrm{k})=\bar{\varphi}_{\mathrm{G}}^{\mathbb{N}}(\mathrm{k})$ for all $\mathrm{k} \in \mathbb{N}$. This $\bar{\varphi}_{\mathrm{G}}$ is called the flow polynomial of G .

Example 1.3.8. We can compute the flow polynomial for the graphs in Figures 2 and 4 and compare them with the associated chromatic polynomial.

| $G$ | Figure2a | Figure2b | Figure 4a | Figure 4b |
| :--- | :---: | :---: | :---: | :---: |
| $\chi_{G}(t)$ | $t(t-1)(t-2)$ | $t(t-\overline{1})$ | $t(t-1)(t-2)^{2}$ | $t(t-1)(t-2)$ |
| $\bar{\varphi}_{G}(t)$ | $(t-1)$ | $(t-1)(t-2)$ | $(t-1)(t-2)$ | $(k-1)(k-2)^{2}$ |

We have seen that for some graphs, for example $K_{4}$, the chromatic and flow polynomials are related up to a multiplication by the variable. In the table above we see two pairs graphs where the chromatic polynomial of the one is the flow polynomial of the other (again, up to a multiplication of the variable), and vice-versa. Those graphs are related by what is called "planar duality".

## Digression: planar duality of graphs.

Let $\mathrm{G}=(\mathrm{V}, \mathrm{E}, \mathrm{t}, \mathrm{s})$ be a graph. A planar drawing of G is a set of points and paths in the plane satisfying some properties. More precisely, a drawing consists of

- one point $\mathrm{p}_{v} \in \mathbb{R}^{2}$ for each $v \in \mathrm{~V}$, so that the function $v \mapsto \mathrm{p}_{v}$ is injective and $p_{V}$ is discrete in $\mathbb{R}^{2}$.
- a simple, regular path $\omega_{e}$ for every $e \in E$, i.e., an injective and continuous function $\omega_{e}:[0,1] \rightarrow \mathbb{R}^{2}$ with $\omega_{e}(0)=p_{t(e)}$ and $\omega_{e}(1)=p_{h(e)}$, and for which $\nabla \omega_{e}(t) \neq 0$ for all $\left.t \in\right] 0,1[$.

We require that the paths representing edges only intersect at points labeled by the vertices incident to the edge under consideration. Precisely: for every $e, e^{\prime} \in E$, we require that $\omega_{e}(t)=\omega_{e}\left(t^{\prime}\right)$ implies $\left\{t, t^{\prime}\right\} \subseteq\{0,1\}$.

This definition entails that if $C$ is the set of edges of a circuit of $G$, the set $\bigcup_{e \in C} \omega_{e}([0,1])$ is a Jordan curve. Thus, by the Jordan curve theorem every circuit of $G$ determines a bipartition of $\mathbb{R}^{2}$ in two open and connected regions (one bounded and one unbounded). The (open, connected) parts of the partition of $\mathbb{R}^{2}$ that refines all bipartitions obtained by circuits of G are called the regions determined by the given drawing of G . They are all bounded and contractible, with exception of one unbounded (and in general not-contractible) region. The boundary of every region is a circuit of G drawn in the plane.

We can define now a planar dual $\mathrm{G}^{*}$ to G as follows. The vertices of $\mathrm{G}^{*}$ are the regions of the drawing, and the edge-set can be taken to be the same as the edge-set of $G$, but for clarity we will call it $E^{*}$, with a fixed bijection $e \mapsto e^{*}$ between $E$ and $E^{*}$. The ends of any $e^{*}$ are determined as follows: look at the drawing $\omega_{e}$ of the edge $e$ and notice that $\omega_{e}([0,1])$ is contained in the (topological) boundary of at most two regions, say $R_{1}, R_{2}$ (where we will allow $R_{1}=R_{2}$ ). As we go along the boundary of $R_{1}$ and $R_{2}$ in positive direction (counterclockwise), we'll meet in one case (say $R_{1}$ ) first $t(e)$ and then $h(e)$, in the other (say $R_{2}$ ) the reverse will happen. We will then say that $t^{*}\left(e^{*}\right)=R_{1}, h^{*}\left(e^{*}\right)=R_{2}$.

Alternatively (and most commonly), one obtains a drawing of $\mathrm{G}^{*}$ by choosing a point (representing a vertex of $\mathrm{G}^{*}$ ) in each region of the given drawing of $G$ and drawing a simple regular path (representing an edge $e^{*}$ of $G^{*}$ ) across each edge $e$ of $G$, whose endpoints are the vertices of the region(s) in whose boundary $\omega_{e}([0,1])$ lays, and that crosses the path representing $e$ at a unique crossing point $p(e)$. The ends $t\left(e^{*}\right)$ and $h\left(e^{*}\right)$ are assigned so that the pair $\left(\nabla \omega_{e}, \nabla \omega_{e^{*}}\right)$ taken at $p(e)$ is a positively oriented basis of the plane. (intuitively: " $e^{*}$ crosses $e$ from left to right", see Figure 77.

Different drawings of the same abstract graph G can give rise to nonisomorphic duals. However, it is always true that $\left(G^{*}\right)^{*}$ is isomorphic to G.


Figure 7: Illustration of the "orientation" of primal-dual edges

Proposition 1.3.9. Let G be a finite, connected planar graph and let $\mathrm{G}^{*}$ be its planar dual. For every $\mathrm{k} \in \mathbb{N} \backslash\{0\}$, there is a $k$-to-one surjective map between the set $\Gamma_{\mathrm{G}}^{(\mathrm{k})}$ of k -colorings of G and the set $\bar{\Phi}_{\mathrm{G}^{*}}^{(\mathrm{k})}$ of nowhere-zero k -flows on $\mathrm{G}^{*}$.
Corollary 1.3.10. Let G be a finite, connected planar graph. Then $\mathrm{X}_{\mathrm{G}}(\mathrm{t})=\mathrm{t} \bar{\varphi}_{\mathrm{G}^{*}}(\mathrm{t})$.
Proof of Proposition 1.3 .9 Let $G$ be any given finite planar graph and consider its planar dual $G^{*}$. In particular, every edge $e$ of $G$ is paired with an edge $e^{*}$ of $G^{*}$, i.e., the one that crosses $e$ (at a unique crossing point $p(e)$ ) in the given planar drawing of $G$.

Now, given a k-coloring $\gamma$ of $G$ define a function $f^{*}: E^{*} \rightarrow \mathbb{Z}_{k}$ as follows:

$$
f_{\gamma}\left(e^{*}\right):=\gamma(t(e))-\gamma(h(e)),
$$

where on the r.h.s. we take the congruence class modulo $k$. It is easy to see that $\gamma$ being a proper coloring implies that $f_{\gamma}$ is nowhere zero. In order to check the condition for $\mathrm{f}_{\gamma}$ to be a flow, consider any vertex $v^{*}$ of $\mathrm{G}^{*}$. This vertex corresponds to a region of the plane that is bounded by a circuit of G. Call $C: v_{1}, e_{1}, v_{2}, \ldots$ this circuit, with vertices and edges listed in counterclockwise direction. Then every edge of $\mathrm{G}^{*}$ that insists on $v^{*}$ is of the form $\left(e_{i}\right)^{*}$ for some $e_{i}$ in C , and by construction of the planar dual we have that $v^{*}=\mathrm{t}\left(e_{i}^{*}\right)$ if and only if $v_{i}=t\left(e_{i}\right)$ (and thus $v_{i+1}=h\left(e_{i}\right)$ ), see Figure 8 ,

Now we can compute Equation

$$
\begin{aligned}
& \sum_{v^{*}=\mathfrak{t}\left(e^{*}\right)} f_{\gamma}\left(e^{*}\right)-\sum_{v^{*}=h\left(e^{*}\right)} f_{\gamma}\left(e^{*}\right)= \\
& \sum_{v_{i}=\mathfrak{t}\left(e_{i}\right)} \gamma\left(t\left(e_{i}\right)\right)-\gamma\left(h\left(e_{i}\right)\right)+\sum_{v_{i}=h\left(e_{i}\right)} \gamma\left(h\left(e_{i}\right)\right)-\gamma\left(t\left(e_{i}\right)\right)= \\
& \sum_{v_{i}=t\left(e_{i}\right)} \gamma\left(v_{i}\right)-\gamma\left(v_{i+1}\right)+\sum_{v_{i}=h\left(e_{i}\right)} \gamma\left(v_{i}\right)-\gamma\left(v_{i+1}\right)=\sum_{v_{i}} \gamma\left(v_{i}\right)-\gamma\left(v_{i}\right)=0
\end{aligned}
$$

Now, fix a vertex $v_{0} \in \mathrm{~V}$. Let $\Gamma$ be the set of all k-colorings of $G$ and let $\Phi$ be the set of all nowhere-zero k-flows on $\mathrm{G}^{*}$. We already know that the function

$$
\Gamma_{\mathrm{G}}^{(\mathrm{k})} \rightarrow \bar{\Phi}_{\mathrm{G}}^{(\mathrm{k})} \times\{1, \ldots, \mathrm{k}\}, \quad \gamma \mapsto\left(\mathrm{f}_{\gamma}, \gamma\left(v_{0}\right)\right)
$$



Figure 8: Here $v^{*}=\mathrm{t}\left(e_{\mathrm{j}}^{*}\right)=\mathrm{h}\left(e_{\mathrm{i}}^{*}\right)$
is well-defined. We are left with proving that it is bijective, which we will do by providing an inverse. Let f be a nowhere-zero k -flow on $\mathrm{G}^{*}$ and let $\mathfrak{i}_{0} \in[\mathrm{k}]$. For every vertex $v$ of $G$ choose a path $v_{0}, e_{1}, v_{1}, \ldots, v_{l}=v$. Define

$$
\gamma_{f}(v) \in\{1, \ldots k\}, \quad \gamma_{f}(v) \equiv i_{0}+\sum_{h\left(e_{i}\right)=v_{i}} f\left(\left(e_{i}\right)^{*}\right)-\sum_{\mathfrak{f}\left(e_{i}\right)=v_{i}} f\left(\left(e_{i}\right)^{*}\right) \bmod k
$$

Notice that $\gamma_{f}(v)$ does not depend on the choice of paths (we prove this as Exercise [I.8) and is thus well-defined. Now notice that for any edge $e$ we have $\gamma_{\mathrm{f}}(\mathrm{h}(\mathrm{e}))-\gamma_{\mathrm{f}}(\mathrm{t}(e))=\mathrm{f}\left(e^{*}\right)$ and thus the fact that f is nowhere-zero implies that $\gamma_{f}$ is a proper $k$-coloring.

Since the double dual of a planar graph is the graph itself, Corollary 1.3.10 can be rewritten as saying that $\bar{\varphi}_{G}(t)=t^{-1} \chi_{G^{*}}(t)$ for any connected planar graph G. Thus Proposition 1.2 .6 can give us an explicit expression for the flow polynomial of a connected planar graph, provided we can interpret $\mathfrak{c}\left(\mathrm{G}^{*}[A]\right)$ in terms of G.

Definition 1.3.11. Let $G$ be any finite graph. Define the cyclomatic number of G to be

$$
\beta_{1}(\mathrm{G}):=|\mathrm{E}|-|\mathrm{V}|+\mathrm{c}(\mathrm{G}) .
$$

As the name suggests, this is the first Betti number of the graph $G$ viewed as a 1 -dimensional simplicial complex. If G is planar, this is one less than the number of faces.

Lemma 1.3.12. Let G be any finite graph, e an edge of G . Then

$$
\beta_{1}(\mathrm{G} \backslash e)= \begin{cases}\beta_{1}(\mathrm{G}) & \text { if } \mathrm{e} \text { is an isthmus, }  \tag{1.6}\\ \beta_{1}(\mathrm{G})-1 & \text { otherwise. }\end{cases}
$$

$$
\beta_{1}(G / e)= \begin{cases}\beta_{1}(G)-1 & \text { if } e \text { is a loop },  \tag{1.7}\\ \beta_{1}(G) & \text { otherwise. }\end{cases}
$$

Theorem 1.3.13. Let G be a finite graph. The modular flow polynomial of G can be written as

$$
\bar{\varphi}_{G}(\mathrm{t})=(-1)^{|\mathrm{E}|} \sum_{\mathrm{A} \mathrm{\subseteq E}}(-1)^{|\mathrm{A}|} \mathrm{t}^{\beta_{1}(\mathrm{G}[\mathrm{~A}])}
$$

Proof. Let $\mathrm{P}_{\mathrm{G}}(\mathrm{t})$ denote the sum on the r.h.s. in the claim. We have to show that it equals $\bar{\varphi}_{G}(\mathrm{t})$. We argue using Theorem 1.3.6, which determines the flow polynomial fully.

- If $G$ has no edges, then the sum that defines $P_{G}(t)$ has only one summand (for $A=\emptyset$ ) and the cyclomatic number is 0 . Thus, $\mathrm{P}_{\mathrm{G}}(\mathrm{t})=$ $(-1)^{0}(-1)^{0} t^{0}=1$.
- Let $e$ be an isthmus of G . With Lemma 1.3.12 we compute

$$
\begin{aligned}
P_{G}(t) & =(-1)^{|E|} \sum_{A \subseteq E \backslash\{e\}}\left((-1)^{|A|} t^{\beta_{1}(G[A])}+(-1)^{|A \cup\{e\}|} t^{\beta_{1}(G[A \cup\{e\}])}\right) \\
& =(-1)^{|E|} \sum_{A \subseteq E \backslash\{e\}}\left((-1)^{|A|} t^{\beta_{1}(G[A])}-(-1)^{|A|} t^{\beta_{1}(G[A])}\right)=0
\end{aligned}
$$

- Let e be a loop of G. Again, we compute with Lemma 1.3.12

$$
\begin{aligned}
P_{G}(t) & =(-1)^{|E|} \sum_{A \subseteq E \backslash\{e\}}\left((-1)^{|A|} t^{\beta_{1}(G[A])}+(-1)^{|A \cup\{e\}|} t^{\beta_{1}(G[A \cup\{e\}])}\right) \\
& =(-1)^{|E|} \sum_{A \subseteq E \backslash\{e\}}(1-t)(-1)^{|A|} t^{\beta_{1}(G[A])} \\
& =(t-1)(-1)^{|E \backslash\{e\}|} \sum_{A \subseteq E \backslash\{e\}}(-1)^{|A|} t^{\beta_{1}(G[A])}=(t-1) P_{G \backslash e}(t)
\end{aligned}
$$

- Now suppose that $e$ is no loop. Then, contracting $e$ does not change the cyclomatic number. We compute

$$
\begin{aligned}
P_{G}(t) & =(-1)^{|E|}\left(\sum_{e \in A \subseteq E}(-1)^{|A|} t^{\beta_{1}(G[A] / e)}+\sum_{A \subseteq E \backslash\{e\}}(-1)^{|A|} t^{\beta_{1}(G[A])}\right) \\
& =(-1)^{|E|-1}\left(\sum_{A \subseteq E \backslash\{e\}}(-1)^{|A|^{\prime} \boldsymbol{t}_{1}(G[A] / e)}-\sum_{A \subseteq E \backslash\{e\}}(-1)^{|A|} t^{\beta_{1}(G[A])}\right) \\
& =P_{G / e}(t)-P_{G \backslash e}(t) .
\end{aligned}
$$

This completes the proof.

## Exercises I

I. 1 Show that a finite graph $G$ with $n$ vertices and $m$ edges is cycle-free if and only if $c(G)=n-m$.
I. 2 The degree of a vertex $v$ of a finite graph G is the "total number of vertexedge incidences at $v^{\prime \prime}$, i.e.,

$$
\mathrm{d}(v)=|\{e \in \mathrm{E} \mid \mathrm{h}(e)=v\}|+|\{e \in \mathrm{E} \mid \mathrm{t}(e)=v\}| .
$$

Prove that a finite graph G contains a cycle if and only if it contains a subgraph $G[A]$ whose vertex-degrees are all even.
I. 3 Let $G$ be a finite graph and $A, B \subseteq E$ two disjoint subsets of edges of $G$. Prove that $(G / A) \backslash B=(G \backslash B) / A$.
I. 4 Give a formal proof of the principle of inclusion-exclusion.
I. 5 Determine the chromatic polynomial of the graph in Figure 9 using deletion-contraction. Check your result using no-broken-circuit sets.


Figure 9
I. 6 Determine the modular flow polynomial of the graph in Figure 9 using deletion-contraction. Check your result using no-broken-circuit sets of the planar dual.
I. 7 Let $G=(V, E, h, t)$ be a finite graph and let $S, T \subseteq$ be two subsets of its vertex set. Let $f$ be a flow on $G$ with values in an abelian group $\mathbb{A}$, written additively (i.e., we write " + " for the group operation and consider a function $\mathrm{f}: \mathrm{E} \rightarrow \mathbb{A}$ satisfying Equation (1.4) as an identity between group elements with respect to the group operation).
The total S-T flow is defined as

$$
f(S, T):=\sum_{\substack{e \in E \\ t \\ h(e) \in S \\ h(e) \in T}} f(e)-\sum_{\substack{e \in E \\ h(e) \in S \\ t(e) \in T}} f(e)
$$

Prove that, if $V=S \sqcup T$ is a bipartition of $V$, then $f(S, T)=0$.
(Hint: as a first step, try to prove that $f(S, S)=0$ and $f(S, V)=0$ for all $\mathrm{S} \subseteq \mathrm{V}$ )
I. 8 Prove that the definition of $\gamma_{f}$ in the proof of Proposition 1.3.9 does not depend on the choice of paths. (Hint: you may use the previous exercise.)

## Chapter 2

## The Tutte polynomial

### 2.1 Definition

From our discussion of colorings and modular flows, we see the importance of two graph statistics: the number of connected components and the cyclomatic number. We can then consider the associated two-variable generating function.

Definition 2.1.1.

$$
\begin{equation*}
\mathrm{Q}_{\mathrm{G}}(v, w):=\sum_{A \subseteq E} v^{\mathrm{c}(\mathrm{G}[\mathcal{A}])} w^{\beta_{1}(\mathrm{G}[\mathcal{A}])} \tag{2.1}
\end{equation*}
$$

Remark 2.1.2. This polynomial, called the "dichromate" of G, generalizes both the chromatic and the flow polynomial of the graph. In fact we have

$$
\chi_{\mathrm{G}}(\mathrm{t})=(-1)^{|\mathrm{V}|} \mathrm{Q}_{\mathrm{G}}(-\mathrm{t},-1), \quad \bar{\varphi}_{\mathrm{G}}(\mathrm{t})=(-1)^{|\mathrm{E}|+|\mathrm{V}|} \mathrm{Q}_{\mathrm{G}}(-1,-\mathrm{t})
$$

Example 2.1.3. - If G has no edges, then $\mathrm{Q}_{\mathrm{G}}(v, w)=v^{|\mathrm{V}|}$

- If G has one vertex and a single loop, then $\mathrm{Q}_{\mathrm{G}}(v, w)=v(1+w)$
- If G has two vertices and a single isthmus, then $\mathrm{Q}_{\mathrm{G}}(v, w)=v(1+v)$

Proposition 2.1.4. Let G be a finite graph, e an edge of G .

- If e is a loop or an isthmus of G , then $\mathrm{Q}_{\mathrm{G}}(v, w)=v^{-1} \mathrm{Q}_{\mathrm{G}[e]}(v, w) \mathrm{Q}_{\mathrm{G} \backslash e}(v, w)$
- If e is neither a loop nor an isthmus of G , then $\mathrm{Q}_{\mathrm{G}}(v, w)=\mathrm{Q}_{\mathrm{G} \backslash e}(v, w)+$ $\mathrm{Q}_{\mathrm{G} / e}(v, w)$

Proof. We first expand

$$
\mathrm{Q}_{\mathrm{G}}(v, w)=\sum_{A \subseteq \mathrm{E} \backslash e} v^{\mathrm{c}(\mathrm{G}[\mathrm{~A}])} w^{\beta_{1}(\mathrm{G}[\mathrm{~A}])}+\sum_{e \in A \subseteq \mathrm{E}} v^{\mathrm{c}(\mathrm{G}[\mathrm{~A}])} w^{\beta_{1}(\mathrm{G}[\mathcal{A}])}
$$

Now if $e$ is a loop, the second term can be written as

$$
w \sum_{e \in A \subseteq E} v^{c(G[A \backslash e])} w^{\beta_{1}(G[A \backslash e])},
$$

whence $\mathrm{Q}_{\mathrm{G}}(v, w)=(1+w) \mathrm{Q}_{\mathrm{G} \backslash e}(v, w)$ in this case. If e is an isthmus the first term can be written as

$$
v \sum_{A \subseteq E \backslash e} v^{c(\mathrm{G}[\mathrm{~A}])} w^{\beta_{1}(\mathrm{G}[\mathrm{~A}])},
$$

and then $\mathrm{Q}_{\mathrm{G}}(v, w)=(v+1) \mathrm{Q}_{\mathrm{G} \backslash e}(v, w)$.
Otherwise, the second term equals

$$
\sum_{e \in A \subseteq E} v^{c(G / e[A \backslash e])} w^{\beta_{1}(G / e[A \backslash e])}
$$

and we conclude $\mathrm{Q}_{\mathrm{G}}(v, w)=\mathrm{Q}_{\mathrm{G} \backslash e}(v, w)+\mathrm{Q}_{\mathrm{G} / e}(v, w)$.

The idea of the dichromate $\mathrm{Q}_{\mathrm{G}}$ goes back to W . T. Tutte, but nowadays it is the following variation on it that is commonly called the "Tutte polynomial" of the graph G.

Definition 2.1.5. The Tutte polynomial of a finite graph G is

$$
\begin{aligned}
\mathrm{T}_{\mathrm{G}}(x, y) & :=(x-1)^{-c(G)} \mathrm{Q}_{\mathrm{G}}(x-1, y-1) \\
& =\sum_{A \subseteq E}(x-1)^{r(G)-r(G[A])}(y-1)^{n(G[A])}
\end{aligned}
$$

Where we write r and n for the rank and nullity of a finite graph, defined as

$$
\begin{equation*}
r(G)=|V|-c(G), \quad n(G)=|E|-r(G) \tag{2.2}
\end{equation*}
$$

Remark 2.1.6. Exercise ... shows that the rank is is the cardinality of any maximal cycle-free set of edges of G, and the cyclomatic number equals the number of edges that do not participate in such a cycle-free set (i.e., the nullity of G). In particular, rank and nullity are properties that can be detected without referring to connectivity properties such as isolated vertices etc.
Remark 2.1.7. We can translate Remark 2.1 .2 in terms of the Tutte polynomial.

$$
\chi_{\mathrm{G}}(\mathrm{t})=(-1)^{\mathrm{r}(\mathrm{G})} \mathrm{t}^{\mathrm{c}(\mathrm{G})} \mathrm{T}_{\mathrm{G}}(1-\mathrm{t}, 0), \quad \bar{\varphi}_{\mathrm{G}}(\mathrm{t})=(-1)^{|\mathrm{E}|+\mathrm{r}(\mathrm{G})} \mathrm{T}_{\mathrm{G}}(0,1-\mathrm{t})
$$

### 2.2 Duality and matroids

Lemma 2.2.1. Let G be a finite planar graph, and let $\mathrm{G}^{*}$ a planar dual to G . Then,

$$
\mathrm{r}\left(\mathrm{G}\left[A^{\mathrm{c}}\right]\right)=\mathrm{r}(\mathrm{G})-\mathfrak{n}\left(\mathrm{G}^{*}\left[A^{*}\right]\right)
$$

Proof. The equality holds tautologically if $A$ is empty. Now suppose that for some $A \subsetneq E$ the equality holds and choose any $e \in E \backslash A$. Now, $G\left[(A \cup\{e\})^{c}\right]=$ $G\left[A^{c}\right] \backslash e$ and we know that deleting an element drops the rank if and only if that element is an isthmus, i.e., there was a connected component K of $\mathrm{G}\left[A^{c}\right]$ that is disconnected when removing $e$. Now notice that $K$ is the part of $G$ that lies in one of the bounded regions determined by $G^{*}[A]$, call it R. Now removing e disconnects $K$ exactly when adding $e^{*}$ to $G^{*}\left[A^{*}\right]$ subdivides $R$ into two parts, making $\mathfrak{n}\left(\mathrm{G}^{*}\left[\mathrm{~A}^{*}\right]\right)$ increase. This concludes the inductive step and proves the claim.

Corollary 2.2.2. If G is a finite planar graph and $\mathrm{G}^{*}$ a dual to G , then

$$
\mathrm{T}_{\mathrm{G}^{*}}(x, y)=\mathrm{T}_{\mathrm{G}}(\mathrm{y}, \mathrm{x})
$$

We can now readily re-prove the analog of Corollary 1.3.10 for general planar graphs.

Corollary 2.2.3. For any finite planar graph G we have

$$
\chi_{\mathrm{G}}(\mathrm{t})=\mathrm{t}^{\mathrm{c}(\mathrm{G})} \bar{\varphi}_{\mathrm{G}^{*}}(\mathrm{t})
$$

These last results depend on the fact that planar graphs have duals. The crucial relationship between rank and nullity of primal and dual graphs is expressed in Lemma 2.2.1. From there, recalling the definition of nullity and the fact that the dual to a planar graph has as many edges as the primal graph, we can state that the main "feature" of dual graphs is to provide an interpretation of the quantity

$$
r(G[E \backslash A])-r(G)+|E \backslash A|
$$

as the rank function of another graph (i.e., as $r\left(G^{*}\left[A^{*}\right]\right)$ ).
Our next goal is to introduce a class of integer-valued functions that includes rank functions of graphs and that is closed under the above "duality" operation.

Definition 2.2.4. Given any finite set $E$ and any function $r: 2^{E} \rightarrow \mathbb{N}$, let us write

$$
r^{*}: 2^{E} \rightarrow \mathbb{N}, \quad r^{*}(A)=r(E \backslash A)-r(E)+|A| .
$$

Remark 2.2.5. The operation $\mathrm{r} \mapsto \mathrm{r}^{*}$ is an involution on the set of all functions $r: 2^{\mathrm{E}} \rightarrow \mathbb{N}$ with $\mathrm{r}(\emptyset)=0$. In fact:

$$
\begin{aligned}
\left(r^{*}\right)^{*}(A) & =r^{*}(E \backslash A)-r^{*}(E)+|A| \\
& =r(A)-r(E)+|A|-r(\emptyset)+r(E)-|E|+|E \backslash A|=r(A)-r(\emptyset)=r(A)
\end{aligned}
$$

Definition 2.2.6. Let $E$ be a finite set. We call $\mathscr{M}_{\mathrm{E}}$ the set of all functions $r: 2^{\mathrm{E}} \rightarrow \mathbb{N}$ satisfying:
(r0) for every $A \subseteq E: 0 \leqslant r(A) \leqslant|A|$;
(r1) for every $A \subseteq B \subseteq E: r(A) \leqslant r(B)$;
(r2) for all $A, B \subseteq E: r(A)+r(B) \geqslant r(A \cup B)+r(A \cap B)$.
The elements of $\mathscr{M}_{\mathrm{E}}$ are called "matroid rank functions" on E (or just "matroids on the ground set $\mathrm{E}^{\prime \prime}$ ).

Remark 2.2.7. Let r be a matroid rank function on the set E and $\mathrm{r}^{\prime}$ a matroid rank function on the set $E^{\prime}$. We call these matroids isomorphic if there is a bijection $\mathrm{b}: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ with $\mathrm{r}=\mathrm{r}^{\prime} \circ \mathrm{b}$.

Theorem 2.2.8. Let E be a finite set. The set $\mathscr{M}_{\mathrm{E}}$ is closed under duality, i.e., if r is in $\mathscr{M}_{\mathrm{E}}$ so is also $\mathrm{r}^{*}$.

Proof. Suppose that r satisfies (r0)- (r2) and consider $\mathrm{r}^{*}$.

- $r^{*}$ satisfies ( $r 0$ ). Let $A \subseteq E$. By ( $r 2$ ) applied to $A$ and $E \backslash A$ and ( $r 0$ ) we have $r(A)+r(E \backslash A) \geqslant r(E)$, hence with $(r 1) r(E \backslash A)-r(E) \geqslant-r(A) \geqslant$ $-|\mathcal{A}|$, thus $r^{*}(A)=r(E \backslash A)-r(E)+|A| \geqslant 0$. On the other hand, ( $r 1$ ) applied to $E \backslash A \subseteq E$ gives $r(E \backslash A)-r(E) \leqslant 0$ and thus $r^{*}(A)=r(E \backslash$ $A)-r(E)+|A| \leqslant|A|$.
$-r^{*}$ satisfies ( $r 1$ ). Let $A \subseteq B \subseteq E$. Then ( $r 2$ ) applied to $E \backslash B$ and $B \backslash A$, together with ( $r 0$ ), gives $r(E \backslash B)+r(B \backslash A) \geqslant r(E \backslash A)$. With this, we compute

$$
\begin{aligned}
r^{*}(B)-r^{*}(A) & =r(E \backslash B)+|B|-r(E \backslash A)-|A| \\
& \geqslant|B|-|A|-r(B \backslash A)=|B \backslash A|-r(B \backslash A) \geqslant 0
\end{aligned}
$$

where in the last inequality we used (r0) for $\mathrm{B} \backslash \mathrm{A}$.

- $r^{*}$ satisfies ( $r 2$ ). Let $A, B \subseteq E$. We compute directly

$$
\begin{aligned}
r^{*}(A)+r^{*}(B) & =r(E \backslash A)-r(E)+|A|+r(E \backslash B)-r(E)+|B| \\
& \geqslant r((E \backslash A) \cap(E \backslash B))+|A|+r((E \backslash A) \cup(E \backslash B))+|B|-2 r(E) \\
& =r(E \backslash(A \cup B))-r(E)+|A|+r(E \backslash(A \cap B))-r(E)+|B| \\
& =r^{*}(A \cup B)+r^{*}(A \cap B)+\underbrace{|A|-|A \cup B|-|A \cap B|+|B|}_{=0}
\end{aligned}
$$

where the inequality is (r2) for $r$ applied to $E \backslash A$ and $E \backslash B$.

Proposition-Definition 2.2.9. Let G be any finite graph on the edge set E and, for every $\mathrm{A} \subseteq \mathrm{E}$ write, for brevity, $\mathrm{r}_{\mathrm{G}}(\mathrm{A}):=\mathrm{r}(\mathrm{G}[\mathrm{A}])$. Then, $\mathrm{r}_{\mathrm{G}} \in \mathscr{M}_{\mathrm{E}}$, and we call this the "rank function of the cycle matroid of G ".

Remark 2.2.10. We distinguish two special rank functions. Call $r_{L}$ the rank function of the graph consisting of a single loop (so $r_{L}: 2^{\{e\}} \mapsto \mathbb{N}$ with $r_{L}$ constant equal to 0 ). Moreover, call $r_{I}$ the rank function of the graph consisting of a single isthmus (so $\mathrm{r}_{\mathrm{I}}: 2^{\{e\}} \mapsto \mathbb{N}$ with $\mathrm{r}_{\mathrm{I}}(\{e\})=1$ ).
Remark 2.2.11. Lemma 2.2.1 says that, for every planar graph $G, r_{G}^{*}=r_{G^{*}}$. The class $\mathscr{M}_{\mathrm{E}}$ however contains $\mathrm{r}_{\mathrm{G}}^{*}$ for all graphs with edge set E . Thus the notion of duality for matroids extends the planar graph duality to general finite graphs.

We will prove a more general statement than Proposition-Definition 2.2.9, motivated by the following construction.

Definition 2.2.12. Let $G$ be a finite graph on the edge set $E$ with vertex set $V$. For every $e \in E$ define a vector $w_{e} \in\left(\mathbb{F}_{2}\right)^{V}$ as follows. For every $v \in \mathrm{~V}$ let

$$
\left(w_{e}\right)_{v}:= \begin{cases}1 & \text { if } h(e) \neq \mathrm{t}(e) \text { and } v \in\{\mathrm{~h}(e), \mathrm{t}(e)\} \\ 0 & \text { otherwise }\end{cases}
$$

Lemma 2.2.13. Let G be any finite graph with edge set E and recall Definition 2.2.12 A given $A \subseteq E$ contains a cycle if and only if the set $\left\{w_{a} \mid a \in A\right\}$ is linearly dependent over $\mathbb{F}_{2}$.

Proof. Any $\mathbb{F}_{2}$-linear dependence in $\left\{w_{a}\right\}_{a \in A}$ is of the form $\sum_{i \in D} w_{i}=0$ for some $D \subseteq A$, and this means that every component of the sum $\sum_{i \in D} w_{i}$ contains an even number of 1 s . The parity of the number of 1 s in the $v$-th component is equivalent modulo 2 to the degree of the vertex $v$ in $G[D]$, and thus linear dependencies in $\left\{w_{a}\right\}_{a \in A}$ correspond to subgraphs of $G[A]$ with all-even degrees. Now the claim follows with Exercise I.2, where it is proved that $A \subseteq E$ contains a cycle if and only if $G[A]$ contains a subgraph $G[D]$ with all degrees even.

Corollary 2.2.14. Let $G$ be any finite graph with edge set E and recall Definition 2.2.12 Then, for every $A \subseteq E$ we have

$$
r_{G}(A)=\operatorname{dim}_{\mathbb{F}_{2}} \operatorname{span}\left\{w_{\mathrm{a}} \mid a \in A\right\} .
$$

Proof. The corollary follows by noting that $\mathrm{r}_{\mathrm{G}}(A)$ is the size of any maximal cycle-free subset of $A$ which, via the previous lemma, is exactly the size of a maximal linear independent subset of $\left\{w_{a}\right\}_{\mathfrak{a} \in \mathcal{A}}$.

Theorem 2.2.15. Let E be any finite set, $\mathrm{n} \in \mathbb{N}$ and let $\mathbb{K}$ be any field. Consider an E -tuple of vectors in the vectorspace $\mathbb{K}^{n}$, i.e., a matrix $w \in \mathbb{K}^{n \times E}$. Then the function

$$
r: 2^{\mathrm{E}} \rightarrow \mathbb{N}, \quad A \mapsto \operatorname{dim}_{\mathbb{K}} \operatorname{span}\left\{w_{a} \mid a \in A\right\}
$$

is in the set $\mathscr{M}_{\mathrm{E}}$. We say that the matroid defined by r is representable over the field $\mathbb{K}$.

Proof. Properties (r0) and (r1) can be proved immediately from the notion of dimension of the span of a finite set of vectors. For (r2) let $A, B \subseteq E$ and recall that, for any two subspaces $W_{1}, W_{2}$ of a given finite-dimensional vectorspace, we have the equality

$$
\operatorname{dim} W_{1}+\operatorname{dim} W_{2}=\operatorname{dim}\left(W_{1} \cap W_{2}\right)+\operatorname{dim}\left(W_{1} \cup W_{2}\right)
$$

Property (r2) now follows by applying this equality to $W_{1}=\operatorname{span}\left\{w_{a} \mid a \in A\right\}$, $W_{2}=\operatorname{span}\left\{w_{\mathrm{b}} \mid \mathrm{b} \in \mathrm{B}\right\}$ and noticing that

$$
\operatorname{span}\left\{w_{e} \mid e \in A \cup B\right\} \subseteq \operatorname{span}\left\{w_{a} \mid a \in A\right\} \cap \operatorname{span}\left\{w_{b} \mid b \in B\right\}
$$

We close this section with a discussion of operations on matroids that mimick deletion and contraction of edges in graphs.

Definition 2.2.16. Let $r: 2^{E} \rightarrow \mathbb{N}$ be a matroid rank function, and let $e \in E$. We call $e$ a loop of the matroid defined by $r \operatorname{if} r(e)=0$, and we call $e$ an isthmus if $r(E \backslash e)=r(E)-1$.

The rank function of the restriction of the matroid to any $A \subseteq E$, written $r_{[A]}$, is just the restriction of $r$, as a function, to $2^{A} \subseteq 2^{E}$ (i.e., $r_{[A]}(X)=r(X)$ for all $X \subseteq A$ ). The deletion of $A$ has rank function $r_{\backslash A}:=r_{[E \backslash A]}$ equal to the restriction of $r$ to $E \backslash A$. The rank function of the contraction of any $A \subseteq E$ is

$$
r_{/ A}: 2^{E \backslash A} \rightarrow \mathbb{N}, r_{/ A}(X):=r(X \cup A)-r(A)
$$

Any matroid that is obtained by a sequence of contractions and deletions from a given matroid is called a minor of the given matroid.

### 2.3 Universality

To every matroid rank function $r \in \mathscr{M}_{\mathrm{E}}$ we can associate the Tutte polynomial

$$
\mathrm{T}_{\mathrm{r}}(\mathrm{x}, \mathrm{y}):=\sum_{A \subseteq E}(x-1)^{r(E)-r(A)}(y-1)^{|A|-r(A)}
$$

Remark 2.3.1. $\mathrm{T}_{\mathrm{r}_{\mathrm{G}}}(\mathrm{x}, \mathrm{y})=\mathrm{T}_{\mathrm{G}}(\mathrm{x}, \mathrm{y})$ for every graph G , and $\mathrm{T}_{\mathrm{r}^{*}}(\mathrm{x}, \mathrm{y})=\mathrm{T}_{\mathrm{r}}(\mathrm{y}, \mathrm{x})$.
Example 2.3.2. For the two rank functions of Remark 2.2.10 we have

$$
\mathrm{T}_{\mathrm{r}_{\mathrm{I}}}(x, y)=x, \quad \mathrm{~T}_{\mathrm{r}_{\mathrm{L}}}(x, y)=y
$$

Remark 2.3.3. If r and $\mathrm{r}^{\prime}$ are isomorphic matroid rank functions, then $\mathrm{T}_{\mathrm{r}}(\mathrm{x}, \mathrm{y})=$ $\mathrm{T}_{\mathrm{r}^{\prime}}(\mathrm{x}, \mathrm{y})$.

The deletion-contraction behaviour of Tutte polynomials of graphs that is inherited from the dichromate generalizes to arbitrary matroids, as follows.

Proposition 2.3.4. Let $\mathrm{r} \in \mathscr{M}_{\mathrm{E}}$ be a matroid rank function and let $\mathrm{e} \in \mathrm{E}$. If e is either a loop or an isthmus, $\mathrm{T}_{\mathrm{r}}(\mathrm{x}, \mathrm{y})=\mathrm{T}_{\mathrm{r}_{[e]}}(\mathrm{x}, \mathrm{y}) \mathrm{T}_{\mathrm{r}_{\backslash e}}(\mathrm{x}, \mathrm{y})$. Otherwise,

$$
\mathrm{T}_{\mathrm{r}}(\mathrm{x}, \mathrm{y})=\mathrm{T}_{\mathrm{r}_{\backslash e}}(\mathrm{x}, \mathrm{y})+\mathrm{T}_{\mathrm{r}_{/ e}}(\mathrm{x}, \mathrm{y})
$$

Proof. We leave this proof as an exercise.
Definition 2.3.5. Let $\mathscr{K}$ be a class of matroid( rank function)s that is closed under minors. A Tutte-Grothendieck invariant on $\mathscr{K}$ is any function f defined on $\mathscr{K}$ and with values in a commutative ring $R$, satisfying
(I) $f(r)=f\left(r^{\prime}\right)$ whenever $r$ and $r^{\prime}$ are isomorphic.
(II) $f(r)=f\left(r_{\backslash e}\right)+f\left(r_{/ e}\right)$ if $e$ is not a loop nor an isthmus.
(III) $f(r)=f\left(r_{[e]}\right) f\left(r_{\backslash e}\right)$ if $e$ is a loop or an isthmus.

Such an $f$ is called a generalized ( $\sigma, \tau$ ) Tutte-Grothendieck invariant for some $\sigma, \tau \in R$ if it satisfies (I) and (III) as well as the following variation of (II):
(II') $f(r)=\sigma f\left(r_{\backslash e}\right)+\tau f\left(r_{/ e}\right)$ if $e$ is not a loop nor an isthmus.
Example 2.3.6. The Tutte polynomial is a Tutte-Grothendieck invariant on any minors-closed class of matroids.

Notation 2.3.7. In the following let $\mathscr{K}$ denote a class of matroids that is obtained from a minor-closed class by removing the empty matroid (i.e., the one with $E=\emptyset$ ).

Theorem 2.3.8. Let R be a commutative ring and let $\mathscr{K}$ be as in Notation 2.3.7 Suppose that $\mathrm{f}: \mathscr{K} \rightarrow \mathrm{R}$ is any function satisfying (I) always as well as (II) and (III) whenever $|\mathrm{E}| \geqslant 2$. Then, for all $\mathrm{r} \in \mathscr{K}$,

$$
f(r)=T_{r}\left(f\left(r_{I}\right), f\left(r_{L}\right)\right)
$$

Proof. For every matroid $r \in \mathscr{K}$, a finite number of applications of (II) and (III) allows to decompose both sides of the claimed equality into the same sum of terms associated to matroids with one element, i.e., with $|E|=1$. But in this case either $r=r_{I}$ or $r=r_{L}$. The claim now follows with Example 2.3.2, where we see that $T_{r_{L}}\left(f\left(r_{I}\right), f\left(r_{L}\right)\right)=f\left(r_{L}\right)$ and $T_{r_{I}}\left(f\left(r_{I}\right), f\left(r_{L}\right)\right)=f\left(r_{I}\right)$.

A similar argument proves the following universality result for generalized Tutte-Grothendieck invariants.

Theorem 2.3.9. Let $\sigma$ and $\tau$ be non-zero elements of a field $\mathbb{F}$. Then there is a unique $(\sigma, \tau)$-generalized Tutte-Grothendieck invariant $f$ on $\mathscr{K}$ with values in $\mathbb{F}[x, y]$ and such that $f\left(r_{I}\right)=x, f\left(r_{L}\right)=y$. This $f$ is given explicitly by

$$
\begin{equation*}
f(r)=\sigma^{|E|-r(E)} \tau^{r(E)} T_{r}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right) \tag{2.3}
\end{equation*}
$$

Proof. One checks explicitly that the right-hand side of Equation (2.3) is indeed a generalized $(\sigma, \tau)$ Tutte-Grothendieck invariant. I.e., for (II') let $e \in E$ be neither an isthmus nor a loop, so that $r(E)=r_{\backslash e}\left(E^{\prime}\right)=r_{/ e}\left(E^{\prime}\right)+1$, where we write $E^{\prime}:=E \backslash e$. Then,

$$
\begin{aligned}
& \tau\left[\sigma^{\left|E^{\prime}\right|-r_{/ e}\left(E^{\prime}\right)} \tau^{r_{/ e}\left(E^{\prime}\right)} \mathrm{T}_{r_{/ e}}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right)\right]+\sigma\left[\sigma^{\left|E^{\prime}\right|-r_{\backslash e}\left(E^{\prime}\right)} \tau^{r^{\text {le }}}{ }^{\left(E^{\prime}\right)} \mathrm{T}_{r_{\backslash e}}\left(\frac{x}{\tau^{\prime}}, \frac{y}{\sigma}\right)\right] \\
& =\tau\left[\sigma^{|E|-r(E)} \tau^{r(E)-1} T_{r_{/ e}}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right)\right]+\sigma\left[\sigma^{|E|-r(E)-1} \tau^{r(E)} T_{r_{\backslash e}}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right)\right] \\
& =\sigma^{|E|-r(E)} \tau^{r(E)}\left(\mathrm{T}_{r_{/ e}}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right)+\mathrm{T}_{r_{\backslash e}}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right)\right)=\sigma^{|E|-r(E)} \tau^{r(E)} \mathrm{T}_{r}\left(\frac{x}{\tau}, \frac{y}{\sigma}\right)
\end{aligned}
$$

We leave the verification of (III) as an exercise. The uniqueness part of the claim follows as in the proof of Theorem 2.3.8.

Example 2.3.10. We have seen that, on the class $\mathscr{K}$ of nonempty graphic matroids, the modular flow function is a $(-1,1)$-generalized Tutte-Grothendieck invariant, with value 0 on $r_{I}$ and value $k-1$ on $r_{L}$. The universality theorem gives us then directly $\bar{\varphi}(t)=(-1)^{|E|-r(E)} \mathrm{T}_{\mathrm{r}}(0,1-\mathrm{t})$, and thus the expression in Theorem 1.3.13

## Exercises II

II. 1 Compute the Tutte polynomial of the graph in Figure 9. Check your computation against the results you obtained in Exercise I. 5 .
II. 2 Show that the graphs in Figure 1 have the same Tutte polynomial but that the associated matroids are not isomorphic (i.e., there is no bijection between the edges of the two that preserves rank).


Figure 1
II. 3 Let $r$ be the rank function of a matroid on the set $E$ and let $A \subseteq E$. Prove that the dual of the contraction of $A$ equals the deletion of $A$ from the dual (i.e., prove that $\left.\left(r_{/ A}\right)^{*}=\left(r^{*}\right)_{\backslash A}\right)$.
II. 4 Prove Proposition 2.3.4.
II. 5 In analogy with the case of graphs, define the characteristic polynomial of a given matroid $r$ as $\chi_{r}(t):=(-1)^{r(E)} \mathrm{T}_{\mathrm{r}}(1-\mathrm{t}, 0)$.

## Chapter 3

## Geometric lattices and arrangements of hyperplanes

### 3.1 Arrangements of hyperplanes

Definition 3.1.1. Let $V$ be vectorspace of finite dimension d . An arrangement of hyperplanes in $V$ is a finite set

$$
\mathscr{A}:=\left\{\mathrm{H}_{1}, \ldots, \mathrm{H}_{\mathrm{n}}\right\}
$$

of codimension 1 linear subspaces of V . The poset of intersections of $\mathscr{A}$ is the set

$$
\mathcal{L}(\mathscr{A}):=\left\{\bigcap_{i \in \mathrm{I}} \mathrm{H}_{\mathrm{i}} \mid \mathrm{I} \subseteq[\mathfrak{n}]\right\}
$$

ordered by reverse inclusion: $X \leqslant Y$ if and only if $X \supseteq Y$.
Example 3.1.2. Let $\mathscr{A}$ be the arrengement in $\mathbb{R}^{3}$ consisting of the four planes

$$
\alpha:\{x=0\}, \quad \beta:\{y=0\}, \quad \gamma:\{x=y\}, \quad \delta:\{z=0\},
$$

depicted in Figure 1 Then $\mathcal{L}(\mathscr{A})$ is the poset represented on the r.h.s. of Figure 1.

A classical question is the following: given an arrangement of hyperplanes $\mathscr{A}$ in $\mathbb{R}^{n}$, how many regions are cut out of $\mathbb{R}^{n}$ by $\mathscr{A}$ ?

In order to treat this question via the universality of Tutte polynomials, let us set up some theory.

Definition 3.1.3. A partially ordered set P is a lattice if, for any two elements $p, q \in P$,


Figure 1

- the subposet $P_{\geqslant p} \cap P_{\geqslant q}$ of all upper bounds to $p$ and $q$ has a unique minimal element - called join of $p$ and $q$ and denoted $p \vee q$, and
- the subposet $P_{\leqslant p} \cap P_{\leqslant q}$ of all lower bounds to $p$ and $q$ has a unique maximal element - called meet of p and q and denoted $\mathrm{p} \wedge \mathrm{q}$.

Notice that every finite lattice must have a unique minimal element (denoted by $\hat{0}$ ) and a unique maximal element (written $\hat{1}$ ).

Definition 3.1.4. Let $P$ be a poset with a unique minimal element $\widehat{0}$ (we call $P$ "bounded below"). Then the atoms of P are the elements of the set

$$
A(P):=\{p \in P \mid p \gtrdot \widehat{0}\}
$$

Recall that every finite lattice has a unique minimal element.
Definition 3.1.5. Let $L$ be a finite lattice. We call $L$ geometric if, for all $x, y \in L$ :
(G) $x \lessdot y$ if and only if there is $p \in A(L), p \nless x$, such that $y=x \vee p$.

Example 3.1.6. Unique least upper bounds exist in $\mathcal{L}(\mathscr{A})$ (for $X, Y \in \mathcal{L}(\mathscr{A})$ take $X \vee Y:=X \cap Y)$. Moreover, since $\mathcal{L}(\mathscr{A})$ is finite, this implies that also unique greater lower bounds exist (take $X \wedge Y:=\vee\{Z \in \mathcal{L}(\mathscr{A}) \mid Z \leqslant X, Z \leqslant Y\}$ ). Thus, $\mathcal{L}(\mathscr{A})$ is a finite lattice.

Now, the atoms of $\mathcal{L}(\mathscr{A})$ are exactly the elements of $\mathscr{A}$, i.e., the hyperplanes. The other nontrivial elements of $\mathcal{L}(\mathscr{A})$ are subspaces of V obtained as intersections of the hyperplanes. Notice here that if $W$ is any linear subspace and $H$ is any hyperplane, the codimension of $H \cap W$ either equals that of $W$ (namely if $\mathrm{H} \supseteq W$ ) or else it surpasses it by one. Therefore, for $W_{1}, W_{2} \in \mathcal{L}(\mathscr{A})$, we have $W_{1} \lessdot W_{2}$ if and only if $W_{2}=W_{1} \cap H$ for some $\mathrm{H} \nsupseteq \mathrm{W}_{1}$ (i.e., $\mathrm{W}_{2}=\mathrm{W}_{1} \vee \mathrm{H}$ for some $\left.\mathrm{H} \in A(\mathcal{L}(\mathscr{A})), \mathrm{H} \nless \mathrm{W}_{1}\right)$. In summary, we see that if $\mathscr{A}$ is an arrangement of hyperplanes, then $\mathcal{L}(\mathscr{A})$ is a geometric lattice.

### 3.2 Matroids from geometric lattices

In what follows we will derive from the definition some properties of a geometric lattice that are "intuitively evident" for intersection posets of hyperplane arrangements. One of these properties is that intersection posets come with a function that assigns to every intersection its codimension as a subspace of V , and this function increases exactly by one along every covering relation. We say that intersection posets are ranked. More generally, we have the following definition.

Definition 3.2.1. Let $P$ be a poset. A rank function for $P$ is a function $\rho: P \rightarrow \mathbb{N}$ such that
(i) $\rho(x)=0$ if $x$ is a minimal element in $P$,
(ii) $\rho(x)+1=\rho(y)$ if $x \lessdot y$ in $P$.

Remark 3.2.2. Notice that, if a bounded-below poset admits a rank function, then this function is unique.

Before going forward, let us establish that the length of a chain $\omega=\left\{x_{0}<\right.$ $\left.\ldots<x_{k}\right\}$ in a partially ordered set $P$ is $\ell(\omega)=|\omega|-1=k$. The length of the poset $\ell(P)$ then is the maximum length of any chain in $P$.

Lemma 3.2.3. In a geometric lattice any two maximal chains between the same elements have the same length.

Proof. Let L be a geometric lattice. We prove by induction the following statement (note that in this proof, given $a, b \in L$, an $(a, b)$-chain is any chain in $L$ of the form $a=x_{0}<x_{1}<\ldots<x_{k}=b$ ).
$\left(*_{t}\right)$ For all $\mathrm{a}, \mathrm{b} \in \mathrm{L}$, if one maximal $(\mathrm{a}, \mathrm{b})$-chain has length t , then all of them do.
The premise of $\left(*_{1}\right)$ can only be satisfied if $a \lessdot b$. In this case there is only one maximal ( $a, b$ )-chain, hence $\left(*_{1}\right)$ holds.

Then let $t \geqslant 2$ and suppose that $\left(*_{r}\right)$ holds for all $r<t$. Consider two maximal ( $a, b$ )-chains

$$
\mathrm{a}=\mathrm{c}_{0} \lessdot \mathrm{c}_{1} \lessdot \ldots \lessdot \mathrm{c}_{\mathrm{t}}=\mathrm{b} \quad \mathrm{a}=\mathrm{d}_{0} \lessdot \mathrm{~d}_{1} \lessdot \ldots \lessdot \mathrm{~d}_{\mathrm{s}}=\mathrm{b} .
$$

Now, if $c_{1}=d_{1}$, then by induction hypothesis all maximal ( $\left.c_{1}, b\right)$-chains have $t-1$ elements, hence $s=t$ and we are done.

Suppose then $c_{1} \neq d_{1}$. By property (G) we can find $x, y \in A(L)$ with $c_{1}=a \vee x, d_{1}=a \vee y$. If $x \leqslant d_{1}$ (resp. $y \leqslant c_{1}$ ) we would have $c_{1} \leqslant d_{1}$ (resp. $d_{1} \leqslant c_{1}$ ), reaching a contradiction; hence, $x \nless \mathrm{~d}_{1}$ (resp. y $\notin \mathrm{c}_{1}$ ). Again by (G), we compute $c_{1} \vee d_{1}=a \vee x \vee y \gtrdot d_{1}, c_{1}$.
Now, by induction hypothesis applied to ( $c_{1}, b$ ), every maximal $\left(c_{1}, b\right)$-chain has length $t-1$, and in particular every maximal $\left(c_{1} \vee d_{1}, b\right)$-chain has length $t-2$. In the same way, induction hypothesis applied to $\left(d_{1}, b\right)$ gives that every $\left(c_{1} \vee d_{1}, b\right)$-chain has length $s-2$. We conclude $s=t$, and $\left(*_{t}\right)$ holds.


Corollary 3.2.4. Every geometric lattice admits a rank function.
Proof. Given a geometric lattice L a rank function is given by choosing, for every $x \in L$,

$$
\begin{equation*}
\rho(x):=\text { length of any maximal chain from } \widehat{0} \text { to } x . \tag{3.1}
\end{equation*}
$$

Lemma Lemma 3.2.3 ensures that this is well-defined, and one readily checks that the conditions of Definition 3.2.1 are satisfied.

Corollary 3.2.5. Let L be a geometric lattice with rank function $\rho$. For every $\mathrm{X} \subseteq$ $A(\mathrm{~L})$ we have $\rho(\mathrm{VX}) \leqslant \# X$.

Proof. First notice that by uniqueness of the rank function we know that $\rho$ can be expressed as in Equation (3.1). Induction on the cardinality of $X$. If $X=\emptyset$, $\rho(\vee X)=\rho(\widehat{0})=0$ and the claim holds.

If $\# X>0$, choose $x \in X$ and notice that either $V(X \backslash\{x\})=V X$ (when $x \leqslant V(X \backslash\{x\}))$ or, by $(G), V(X \backslash\{x\}) \lessdot \vee X$. In any case, a maximal chain from $\widehat{0}$ to $V X$ can be obtained by adding at most one new element to a maximal chain from $\hat{0}$ to $\vee(X \backslash\{x\})$. Therefore, $\rho(\vee X) \leqslant \rho(V(X \backslash\{x\}))+1$ and by induction hypothesis this is at most $\# X$.
Lemma 3.2.6. Let L be a geometric lattice and $\rho$ it $\}^{1}$ rank function. Then, for all $x, y \in L$,

$$
\rho(x)+\rho(y) \geqslant \rho(x \wedge y)+\rho(x \vee y) .
$$

Proof. Consider $z:=x \wedge y$ and any saturated chain $z=z_{0} \lessdot z_{1}, \lessdot z_{2}, \lessdot \cdots \lessdot$ $z_{\mathrm{k}}=\mathrm{y}$. Then,

$$
\begin{equation*}
k=\rho(y)-\rho(x \wedge y) \tag{3.2}
\end{equation*}
$$

By (G) we can choose atoms $a_{1}, \ldots, a_{k}$ so that $a_{i} \leqslant z_{i}, a_{i} \notin z_{i-1}$ and $z_{i}=$ $z_{i-1} \vee a_{i}$ for all $i=1, \ldots, k$.

[^1]Define now elements $w_{0}, \ldots, w_{k}$ by setting $w_{0}=x$ and $w_{i}:=w_{i-1} \vee a_{i}$ for all $i \geqslant 1$. Notice that $w_{k}=x \vee a_{1} \vee \ldots \vee a_{k}=x \vee z \vee a_{1} \vee \ldots \vee a_{k}=x \vee y$.

Then, by (G) we have either $w_{i}=w_{i-1}$ or $w_{i-1} \lessdot w_{i}$ for all $i$, so that $k \geqslant \rho\left(w_{k}\right)-\rho\left(w_{0}\right)=\rho(x \vee y)-\rho(x)$ and the claim follows by recalling Equation (3.2).

We have proved the following.
Proposition 3.2.7. Let E be a finite set and let $\mathcal{L} \subseteq 2^{\mathrm{E}}$ a family of subsets of E , partially ordered by inclusion and such that $\mathrm{E} \in \mathcal{L}$. Suppose further that $\mathcal{L}$ is a geometric lattice with rank function $\rho$.

Then for every $X \subseteq E$ there is a unique minimal $X^{\prime}$ in $\mathcal{L}$ such that $X \subseteq X^{\prime}$, and the extension r of $\rho$ on $2^{\mathrm{E}}$ given by $\mathrm{r}(\mathrm{X}):=\rho\left(\mathrm{X}^{\prime}\right)$ is a matroid rank function.

Proof. Axiom (r1) is trivially satisfied. For Axiom (r0) notice first that $\rho$ is never negative by definition. Moreover, given $X \subseteq E$ we can consider the family $A_{1}, \ldots, A_{k}$ of all atoms of $\mathcal{L}$ whose intersection with $X$ is not contained in the minimal element $B$ of $\mathcal{L}$. Then surely $k<|X|$, and

$$
X \subseteq B \cup \bigcup_{i} A_{i}=\bigcup_{i} A_{i}
$$

the second equality because $B \subseteq A_{i}$ for all $i$. In particular, we have $X^{\prime} \leqslant \bigvee_{i} A_{i}$ in $\mathcal{L}$, and by Corollary 3.2.5 $\rho\left(\bigvee_{i} A_{i}\right) \leqslant k$. Thus $r(X)=\rho\left(X^{\prime}\right) \leqslant k \leqslant|X|$ as desired.

We now turn to Axiom (r2). First notice that, trivially, $X^{\prime} \wedge Y^{\prime} \geqslant(X \cap Y)^{\prime}$. By definition, $X^{\prime} \vee Y^{\prime}$ is the minimal element of $\mathcal{L}$ containing $X^{\prime}$ and $Y^{\prime}$, while $(X \cup Y)^{\prime}$ is the minimal element of $\mathcal{L}$ containing $X$ and $Y$. Since $X \subseteq X^{\prime}$ and $Y \subseteq Y^{\prime}$, we have $X^{\prime} \vee Y^{\prime} \geqslant(X \cup Y)^{\prime}$. With the trivial inequality $X^{\prime} \vee Y^{\prime} \leqslant$ $(X \cup Y)^{\prime}$ we obtain $X^{\prime} \vee Y^{\prime}=(X \cup Y)^{\prime}$.

Now using Lemma 3.2.6 and the monotony of $\rho$ we can write

$$
\begin{aligned}
r(X)+r(Y) & \stackrel{\text { df }}{=} \rho\left(X^{\prime}\right)+\rho\left(Y^{\prime}\right) \\
& \geqslant \rho\left(X^{\prime} \wedge Y^{\prime}\right)+\rho\left(X^{\prime} \vee Y^{\prime}\right) \geqslant \rho\left((X \cap Y)^{\prime}\right)+\rho\left((X \cup Y)^{\prime}\right) \\
& \stackrel{d f}{=} r(X \cap Y)+r(X \cup Y)
\end{aligned}
$$

Corollary 3.2.8. Given any (abstract) geometric semilattice $\mathcal{L}$, we can associate to every $x \in \mathcal{L}$ the set $\mathcal{A}(x)$ of all atoms of $\mathcal{L}$ below $x$. Then, $\mathcal{L}$ is isomorphic to the set $\mathcal{L}^{\prime}:=\{A(x) \mid x \in \mathcal{L}\}$ ordered by inclusion (since $x<y$ if and only if $A(x) \subset A(y)$ ). The matroid constructed from the proposition, then, has the set $\mathcal{A}(\mathcal{L})$ of all atoms as a ground set and rank function given by $r(X)=\rho(V X)$ for all $X \subseteq A(\mathcal{L})$. This matroid has no loops, and it is referred to as the "simple matroid associated to $\mathcal{L}$.

Example 3.2.9. Let us consider the geometric lattice from Figure 1 The set of atoms is $\{\alpha, \beta, \gamma, \delta\}$, and the associated geometric lattice $\mathcal{L}^{\prime}$ in Corollary 3.2.8 is as follows.


The claim of Corollary 3.2 .8 is then that this is the lattice of flats of a matroid on $E=\{\alpha, \beta, \gamma, \delta\}$ with rank function given by $r(A)=|A|$ if $|A| \leqslant 2$, $r(\{\alpha, \beta, \gamma\})=2$, and $r(A)=3$ for all $A$ with $A \neq\{\alpha, \beta, \gamma\}$ and $|A| \geqslant 3$.

### 3.3 Geometric lattices from matroids

We aim at a "converse" of Proposition 3.2.7, constructing a geometric lattice for every given matroid.
Definition 3.3.1. Let $E$ be a finite set and $r: 2^{E} \rightarrow \mathbb{N}$ a matroid rank function.
Define a closure operator

$$
\mathrm{cl}: 2^{\mathrm{E}} \rightarrow 2^{\mathrm{E}}, \quad \mathrm{X} \mapsto\{\mathrm{e} \in \mathrm{E} \mid \mathrm{r}(\mathrm{X} \cup \mathrm{e})=\mathrm{r}(\mathrm{X})\}
$$

Call $X \subseteq E$ closed if $X=\operatorname{cl}(X)$, and let $\mathcal{L}_{r}$ be the poset of all closed sets ordered by inclusion (i.e., for $F, F^{\prime} \in \mathcal{L}_{r}$ we have $F \leqslant F^{\prime}$ if $F \subseteq F^{\prime}$ ).

## Example 3.3.2.



Example 3.3.3. Consider the rank function $r: 2^{[4]} \rightarrow \mathbb{N}$ defined by $r(X)=1$ if $|X| \leqslant 1$ and $r(X)=2$ otherwise. One can check that this is a matroid rank function - it is called "uniform of rank 2 on 4 elements". The associated poset of flats is depicted below.


Our next goal is to prove that, in general, $\mathcal{L}_{\mathrm{r}}$ is a geometric lattice.
Remark 3.3.4.
(1) For all $X \subseteq E$ we have obviously $X \subseteq \operatorname{cl}(X)$.
(2) For all $X \subseteq E, r(c l(X))=r(X)$.

Proof. Let $Y:=\operatorname{cl}(X) \backslash X$. We prove that $r\left(X \cup Y^{\prime}\right)=r(X)$ for all $Y^{\prime} \subseteq Y$, by induction on $k:=\left|\mathrm{Y}^{\prime}\right|$.
The case $k=1$ holds by definition of " $y \in \operatorname{cl}(X)$ ".
For $k>1$, assume that the claim holds for smaller sizes of $Y^{\prime}$, and choose $y_{1}, y_{2} \in Y^{\prime}$. Then

$$
\begin{aligned}
2 r(X) & \stackrel{\text { I.H. }}{=} r\left(\left(X \cup Y^{\prime}\right) \backslash\left\{y_{1}\right\}\right)+r\left(\left(X \cup Y^{\prime}\right) \backslash\left\{y_{2}\right\}\right) \\
& \stackrel{r 3}{\geqslant} r\left(X \cup Y^{\prime}\right)+r\left(X \cup Y^{\prime} \backslash\left\{y_{1}, y_{2}\right\}\right) \\
& \stackrel{r 2}{\geqslant} r(X)+r(X)
\end{aligned}
$$

and equality must hold throughout. Thus, $r(X)=r\left(X \cup Y^{\prime}\right)$.
(3) For all $\left.X \subseteq E, c l(X)=\max _{\subseteq} \subseteq Y \supseteq X \mid r(Y)=r(X)\right\}$.

Proof. Let $\Gamma:=\{\mathrm{Y} \supseteq \mathrm{X} \mid \mathrm{r}(\mathrm{X})=\mathrm{r}(\mathrm{Y})\}$. By (1) and (2) above we have $c l(X) \in \Gamma$. It is now enough to prove that $Y \subseteq \operatorname{cl}(X)$ for all $Y \in \Gamma$.
To this end, take $Y \in \Gamma$ and, for any $y \in Y$ compute $r(X) \leqslant r(X \cup\{y\}) \leqslant$ $r(Y)=r(X)$ (the last equality because $Y \in \Gamma)$, hence $r(X \cup\{y\})=r(X)$, thus $y \in \operatorname{cl}(X)$. We conclude $Y \subseteq \operatorname{cl}(X)$ as desired.
(4) For all $X \subseteq Y \subseteq E, c l(X) \subseteq \operatorname{cl}(Y)$.

Proof. Clear from (3).
Lemma 3.3.5. Let r be a matroid rank function. Then, meet and join of every $\mathrm{F}_{1}, \mathrm{~F}_{2} \in$ $\mathcal{L}_{\mathrm{r}}$ exist. In fact,
(1) $F_{1} \vee F_{2}=\operatorname{cl}\left(F_{1} \cup F_{2}\right)$
(2) $F_{1} \wedge F_{2}=F_{1} \cap F_{2}$

In particular, $\mathcal{L}_{\mathrm{r}}$ is a lattice.
Proof.
(1) By definition of the ordering, every element of $\left(\mathcal{L}_{r}\right) \geqslant F_{1} \cap\left(\mathcal{L}_{r}\right) \geqslant F_{2}$ must contain $F_{1} \cup F_{2}$. But by the definition of the closure operator, $\operatorname{cl}\left(F_{1} \cup F_{2}\right)$ is the (unique) smallest closed set containing $F_{1} \cup F_{2}$.
(2) It is enough to prove that $F_{1} \cap F_{2}$ is closed, i.e., that $\operatorname{cl}\left(F_{1} \cap F_{2}\right)=F_{1} \cap F_{2}$. We do this next.
$\supseteq$ This inclusion is clear (see Remark 3.3.4(1)).
$\subseteq$ For every $e \in \operatorname{cl}\left(F_{1} \cap F_{2}\right)$ we have $r\left(\left(F_{1} \cap F_{2}\right) \cup\{e\}\right)=r\left(F_{1} \cap F_{2}\right)$. Hence:

$$
\begin{array}{r}
\underbrace{r\left(\left(F_{1} \cap F_{2}\right) \cup\{e\}\right)}_{=r\left(F_{1} \cap F_{2}\right)}+r\left(F_{1}\right) \stackrel{(r 3)}{\geqslant} r\left(F_{1} \cup\{e\}\right)+r\left(F_{1} \cap F_{2} \text { or }\left(F_{1} \cap F_{2}\right) \cup\{e\}\right) \\
\\
=r\left(F_{1} \cup\{e\}\right)+r\left(F_{1} \cap F_{2}\right)
\end{array}
$$

and we conclude $r\left(F_{1}\right)=r\left(F_{1} \cup\{e\}\right)$, hence $e \in F_{1}$.
Analogously we prove $e \in F_{2}$, hence $e \in F_{1} \cap F_{2}$.

Remark 3.3.6. If $X<Y$ in $\mathcal{L}_{r}$, then $r(X)<r(Y)$. Otherwise, by $r 2$ we would have $r(X)=r(Y)$ and by definition of closure $Y \subseteq c l(X)=X$, proving $X=Y$.

Proposition 3.3.7. For any $F_{1}, F_{2} \in \mathcal{L}_{r}$, (G) holds. I.e.,

$$
F_{1} \lessdot F_{2} \Leftrightarrow \exists P \in A\left(\mathcal{L}_{r}\right), P \nless F_{1} \text {, s.t. } F_{2}=F_{1} \vee P .
$$

Proof.
$\Leftarrow$ Let $P$ be as in the claim. Since $P$ is an atom, $P=\operatorname{cl}(\{e\})$ for some element $e \in E$ and, since $P \nless F_{1}$, by Remark 3.3.4 (4) it must be $e \in E \backslash F_{1}$. Now we can write $F_{2}$ using Lemma 3.3.5 as

$$
F_{2}=F_{1} \vee P=\operatorname{cl}\left(F_{1} \cup P\right)=\operatorname{cl}\left(F_{1} \cup\{e\}\right)
$$

and we have

$$
r\left(F_{1}\right)+r(\{e\}) \geqslant r(\emptyset)+r\left(F_{1} \cup\{e\}\right)
$$

Now, $r(\{e\})=r(c l(e))$ and since $P=\operatorname{cl}(e)$ has rank 1, $r(e)=1$. Thus

$$
r\left(F_{2}\right)=r(F \cup\{e\}) \leqslant r\left(F_{1}\right)+1
$$

Moreover, since $F_{1}$ is closed and $e \notin F_{1}$ we have $r\left(F_{1} \cup\{e\}\right)>r\left(F_{1}\right)$, and we conclude that $r\left(F_{2}\right)=r\left(F_{1}\right)+1$.
Now by Remark 3.3.6 any $Z \in \mathcal{L}_{r}, F_{1}<Z<F_{2}$, would force $r\left(F_{2}\right) \geqslant$ $r\left(F_{1}\right)+2$, hence a contradiction. We conclude $F_{1} \lessdot F_{2}$.
$\Rightarrow F_{1} \lessdot F_{2}$ implies $F_{1} \subsetneq F_{2}$ and so we can choose $e \in F_{2} \backslash F_{1}$. Then $r(\{e\})=1$ since otherwise $e$ is in the closure of every flat, in particular we would have $e \in F_{1}$. It follows that $P:=\operatorname{cl}(\{e\})$ is an atom of $\mathcal{L}_{r}$, and $P \leqslant F_{2}$ by Remark 3.3.4(1). Now define

$$
F:=F_{1} \vee P=\operatorname{cl}\left(F_{1} \cup\{e\}\right)
$$

Then the following claim concludes the proof.
Claim. $\mathrm{F}_{2}=\mathrm{F}$.
Proof. We have

$$
\begin{equation*}
r(F) \geqslant r\left(F_{1}\right)+1=r\left(F_{2}\right) \tag{3.3}
\end{equation*}
$$

The inequality holds since $F \supseteq F_{1} \cup\{e\}, F_{1}$ is closed and $e \notin F_{1}$, the equality is immediate since $F_{1} \lessdot F_{2}$.
Now since $F_{1} \cup\{e\} \subseteq F_{2}$ Remark 3.3.4 (1) implies $F \subseteq F_{2}$. Together with Equation (3.3) this shows $F=F_{2}$.

Theorem 3.3.8. Let r be any matroid rank function. Then the poset $\mathcal{L}_{\mathrm{r}}$ is a geometric lattice whose rank function $\rho$ satisfies $\rho(F)=r(F)$ for every $F \in \mathcal{L}_{r}$.

Proof. That $\mathcal{L}_{\mathrm{r}}$ is a geometric lattice follows from Lemma 3.3.5 and Proposition 3.3.7 For the claim about rank consider any $F \in \mathcal{L}_{r}$ and let $\hat{0} \lessdot F_{1} \lessdot \ldots \lessdot$ $F_{k}=F$ be a maximal chain below $F$. Then, $\rho(F)=k$.

Choose atoms $A_{1}, \ldots, A_{k}$ with $F_{i}=F_{i-1} \vee A_{i}$ for all $i$. Since every $F_{i-1}$ is closed and $A_{i} \nsubseteq F_{i-1}$, we must have

$$
\begin{equation*}
r\left(F_{i-1}\right)>r\left(F_{i-1} \cup A_{i}\right)=r\left(F_{i}\right) \tag{3.4}
\end{equation*}
$$

(the last equality by Lemma 3.3.5.(1)). On the other hand, (r2) implies

$$
\begin{equation*}
r\left(F_{i-1}\right)+r\left(A_{i}\right) \geqslant r(\underbrace{F_{i-1} \cap A_{i}}_{=\widehat{0}})+r\left(F_{i-1} \cup A_{i}\right)=r\left(F_{i}\right) \tag{3.5}
\end{equation*}
$$

and since $r\left(A_{i}\right)=1$ from Equations 3.4 and 3.5 we conclude $r\left(F_{i}\right)=$ $r\left(F_{i-1}\right)+1$, thus $r(X)=r\left(F_{k}\right)=k=\rho(X)$.

### 3.4 Interlude: Arrangements

At this stage we have two, a priori different, rank functions associated to an arrangement $\mathscr{A}=\left\{\mathrm{H}_{1}, \ldots, \mathrm{H}_{\mathrm{m}}\right\}$ of hyperplanes:

- The rank function $r_{\text {lat }}$ of the simple matroid associated to the geometric lattice $\mathcal{L}(\mathscr{A})$ as in Corollary 3.2.8

$$
r_{\text {lat }}: 2^{[\mathrm{m}]} \rightarrow \mathbb{N}, \quad \mathrm{I} \mapsto \rho\left(\bigvee_{i \in \mathrm{I}} \mathrm{H}_{\mathrm{i}}\right)
$$

- The rank function $r_{\text {dep }}$ obtained via Theorem 2.2.15 from the [ $n$ ]-tuple of vectors $\left\{n_{1}, \ldots, n_{m}\right\}$, where $n_{i}$ is any choice of normal vector for the hyperplane $\mathrm{H}_{\mathrm{i}}$ :

$$
r_{\text {dep }}: 2^{[m]} \rightarrow \mathbb{N}, \quad \mathrm{I} \mapsto \operatorname{dim} \operatorname{span}\left\{v_{i} \mid i \in \mathrm{I}\right\}
$$

Our next goal is to show that they are the same.
Lemma 3.4.1. For every intersection $X \in \mathcal{L}(\mathscr{A})$ we have $\rho(X)=\operatorname{codim} X$
Proof. By definition $\rho(X)=k$ means that $k$ is the length of a maximal chain $\widehat{0} \lessdot X_{1} \lessdot \cdots \lessdot X_{k}=X$. Now consider the subspaces $X_{i}$. By property (G), every $X_{i}$ is of the form $X_{i-1} \cap H_{i}$ for some atom $H_{i}$ of $\mathcal{L}(\mathscr{A})$ (i.e., hyperplane in $\mathscr{A}$ ) with $H_{i} \notin X_{i-1}$ (i.e., $H_{i} \nsupseteq X_{i-1}$ ). Notice that the latter implies that $X_{i-1}+H_{i}=\mathbb{R}^{\mathrm{d}}$, the ambient space. Now, elementary linear algebra tells us that

$$
\operatorname{dim}(\underbrace{X_{i-1} \cap H_{i}}_{=X_{i}})+\underbrace{\operatorname{dim}\left(X_{i-1}+H_{i}\right)}_{=d}=\operatorname{dim}\left(X_{i-1}\right)+\underbrace{\operatorname{dim}\left(H_{i}\right)}_{d-1}
$$

and thus $\operatorname{dim}\left(X_{i}\right)=\operatorname{dim}\left(X_{i-1}\right)-1$. Therefore, $X$ has dimension $k$ less than $\widehat{0}=\mathbb{R}^{\mathrm{d}}$, and the proof is complete.

Proposition 3.4.2.

$$
\mathrm{r}_{\mathrm{lat}} \equiv \mathrm{r}_{\mathrm{dep}}
$$

Proof. Let $\mathrm{I} \subseteq[m]$ and write $X:=\bigvee_{i \in \mathrm{I}} \mathrm{H}_{\mathrm{i}}=\bigcap_{\mathrm{i} \in \mathrm{I}} \mathrm{H}_{\mathrm{i}}$ Then, with Lemma 3.4.1 we know that $r_{\text {lat }}(I)=\rho(X)=\operatorname{codim}(X)$. On the other hand, $r_{d e p}(I)$ equals the rank of the $d \times|I|$ matrix $M$ whose columns are $v_{i}$ for $i$ in $I$. Now, $X$ is the subspace of all points that are orthogonal to each $v_{i}, i \in I$, and therefore $X=$ $\operatorname{ker} M$. Now, again by elementary linear algebra we know that $\operatorname{dim} \operatorname{ker} M=$ $d-\operatorname{rank} M$. We summarize and conclude

$$
r_{d e p}(I)=\operatorname{rank} M=d-\operatorname{dim} \operatorname{ker} M=d-\operatorname{dim} X=\operatorname{codim}(X)=r_{l a t}(I)
$$

### 3.5 Minors

Proposition 3.5.1. Let $r$ be the rank function of a matroid on E and let $\mathrm{e} \in \mathrm{E}$. Call $\mathrm{cl}, \mathrm{cl}_{\backslash e}, \mathrm{cl}_{/ e}$ the closure operators associated to r, its deletion and contraction respectively. Then, for every $A \subseteq E \backslash e$ we have
(1) $\operatorname{cl}_{\backslash e}(A)=\operatorname{cl}(A) \backslash\{e\}$
(2) $\mathrm{cl}_{/ e}(A)=\operatorname{cl}(A \cup\{e\}) \backslash\{e\}$

Proof. For item (1) we have that $\operatorname{cl}_{\backslash e}(A)$ is the set of all $x \in E \backslash e$ such that $r_{\backslash e}(A \cup x)=r_{\backslash e}(A)$. The claim follows by noticing that $r_{\backslash e}=r$ on $E \backslash e$.

For item (2) notice that $\mathrm{cl}_{/ e}(A)$ is the set of all $x \in E \backslash e$ such that $r_{/ e}(A \cup$ $x)=r_{/ e}(A)$, i.e., such that $r(A \cup x \cup e)-r(e)=r(A \cup e)-r(e)$. This equality clearly holds if and only if $x \in \operatorname{cl}(A \cup e)$. Since $x \neq e$, claim (2) is proved.

Theorem 3.5.2. Let $r$ be the rank function of a matroid on the ground set E and let $e \in E$. Then, $\mathcal{L}_{r_{\backslash e}}$ is the sublattice of $\mathcal{L}_{r}$ generated by the joins of all $x \in E \backslash e$, and $\mathcal{L}_{r_{/ e}} \simeq\left(\mathcal{L}_{\mathrm{r}}\right)_{\geqslant \mathrm{cl}(e)}$.
Proof. Exercise.

### 3.6 Dissection theory

Let $\mathscr{A}$ be an arrangement of hyperplanes in $\mathbb{R}^{\mathrm{d}}$ and let $\mathcal{L}(\mathscr{A})$ be the associated geometric lattice of intersections. If we denote by $r_{\mathscr{A}}$ the rank function of the matroid defined by $\mathcal{L}(\mathscr{A})$ (or, equivalently, by the dependency relations of the normal vectors), we obtain a function

$$
\mathrm{r}: \operatorname{Arr} \rightarrow \text { Mat, } \quad \mathscr{A} \mapsto \mathrm{r}_{\mathscr{A}}
$$

from arrangements to matroids.
Now call $\mathcal{R}(\mathscr{A})$ the set of chambers of $\mathscr{A}\left(\right.$ so $\left.\mathcal{R}(\mathscr{A})=\pi_{0}\left(\mathbb{R}^{\mathrm{d}} \backslash \cup \mathscr{A}\right)\right)$. Our aim will be to apply matroid theory and the universality of the Tutte polynomial in order to give a formula for computing $|\mathcal{R}(\mathscr{A})|$.

For any given $\mathrm{H} \in \mathscr{A}$ consider $\mathscr{A}_{\mathrm{H}}:=\mathscr{A} \backslash\{\mathrm{H}\}$ and $\mathscr{A}^{\mathrm{H}}:=\left\{\mathrm{H}^{\prime} \cap \mathrm{H} \mid\right.$ $\left.\mathrm{H}^{\prime} \in \mathscr{A}_{\mathrm{H}}\right\}$. We see that the regions of $\mathscr{A}^{\mathrm{H}}$ are "walls" that bisect some of the chambers of $\mathscr{A}_{\mathrm{H}}$. In fact, for every chamber of $\mathscr{A}_{\mathrm{H}}$ that is subdivided in passing to $\mathscr{A}$ there is exactly one chamber of $\mathscr{A}^{\mathrm{H}}$ that performs this "cut". We have shown that

$$
\begin{equation*}
|\mathcal{R}(\mathscr{A})|=\left|\mathcal{R}\left(\mathscr{A}^{\mathrm{H}}\right)\right|+\left|\mathcal{R}\left(\mathscr{A}_{\mathrm{H}}\right)\right| . \tag{3.6}
\end{equation*}
$$

## Lemma 3.6.1.

(1) $\mathcal{L}\left(\mathscr{A}_{\mathrm{H}}\right)$ is the sublattice of $\mathcal{L}(\mathscr{A})$ generated by the set of atoms $\mathscr{A} \backslash\{\mathrm{H}\}$.
(2) $\mathcal{L}\left(\mathscr{A}^{\mathrm{H}}\right) \simeq \mathcal{L}(\mathscr{A}) \geqslant \mathrm{H}$.

Proof. Exercise.
Now notice that $\mathscr{A}=\{\mathrm{H}\}$ if and only if $\mathcal{L}_{\mathrm{r}}=\{\widehat{0}<\widehat{1}\}$, and in that case $|\mathcal{R}(\mathscr{A})|=2-$ in particular, this number only depends on the isomorphism type of $\mathcal{L}_{\mathrm{r}}$. With Equation 3.6 and Lemma 3.6.1 we see that this then is true also for bigger arrangements. Thus, we can write

$$
|\mathcal{R}(\mathscr{A})|=\mathrm{f}\left(\mathrm{r}_{\mathscr{A}}\right)
$$

where $\mathrm{f}:$ Mat $\rightarrow \mathbb{Z}$ is a function that is constant on every set of matroids that share the same lattice of flats and that satisfies the deletion-contraction recursion

$$
\begin{equation*}
f(r)=f\left(r_{/ e}\right)+f\left(r_{\backslash e}\right) \tag{3.7}
\end{equation*}
$$

Lemma 3.6.2. Let $\mathrm{f}: M a t \rightarrow \mathbb{Z}$ be a function that is constant on every set of matroids that share the same lattice of flats and satisfies Equation (3.7). Then, $f\left(r_{L}\right)=0$.

Proof. Let $r$ be the rank function of the matroid on $E=\{a, b\}$ with $r(\{a\})=$ $r(\{b\})=r(\{a, b\})=1$. Its flats are $\emptyset$ and $E$, and so its lattice of flats is isomorphic to that of $r_{I}$. Moreover, $r_{/ a}=r_{L}$, and $r_{\backslash a}=r_{I}$. Now Equation 3.7) gives $f\left(r_{I}\right)=f\left(r_{I}\right)+f\left(r_{L}\right)$, implying $f\left(r_{L}\right)=0$ as desired.

Thus, our f satisfies the definition of a Tutte-Grothendieck invariant in $\mathbb{Z}$ (for item (III) in Definition 2.3.5 notice Exercise III.2 and recall that that $f\left(r_{I}\right)=2$ ), with $f\left(r_{L}\right)=0$ and $f\left(r_{I}\right)=2$. We conclude immediately that

$$
|\mathcal{R}(\mathscr{A})|=\mathrm{T}_{\mathrm{r}_{\mathscr{A}}}(2,0) .
$$

### 3.6.1 An example: graphic arrangements

Let $G$ be a loopless graph with edge-set $E$ and set of vertices V. To $G$ is associated an arrangement $\mathscr{A}_{\mathrm{G}}$ in $\mathbb{R}^{V}$ with one hyperplane $\mathrm{H}_{e}$ for each $e \in E$, defined via its normal vector $n_{e}$ that we choose as

$$
\begin{equation*}
\mathrm{n}_{e}=(0, \ldots, 0,1,0, \ldots, 0,-1,0, \ldots, 0) \tag{3.8}
\end{equation*}
$$

with nonzero value at entries $v$ and $w$, where $\{v, w\}=\{h(e), \mathrm{t}(e)\}$. Notice that, identifying edges of the graph with hyperplanes of the arrangement, the rank function of the graphic matroid derived from $G$ is the same as that derived from the (lattice of flats of the) arrangement.

Recall further that, associated to a graph, we have a chromatic polynomial

$$
\chi_{G}(t)=(-1)^{r(G)} t^{c(G)} T_{G}(1-t, 0)
$$

where $T_{G}(x, y)$ is the Tutte polynomial of the (matroid associated to the) graph G.

By our considerations above, we have that the number of regions of $\mathscr{A}_{\mathrm{G}}$ is $\left|\chi_{G}(-1)\right|$. This allows us to easily compute the number of regions of certain classes of such "graphic arrangements". For instance, take $G=K_{n}$, the complete graph on $n$ vertices. The chromatic polynomial of $K_{n}$ is easy to compute, and so is the number of regions of the arrangement:

$$
\chi_{\mathrm{K}_{\mathrm{n}}}(\mathrm{t})=\mathrm{t}(\mathrm{t}-1)(\mathrm{t}-2) \ldots(\mathrm{t}-\mathrm{n}+1), \quad\left|\mathcal{R}\left(\mathscr{A} \mathrm{K}_{\mathrm{n}}\right)\right|=\left|\chi_{\mathrm{K}_{\mathrm{n}}}(-1)\right|=\mathrm{n}!
$$

You might recall from Coxeter theory that this is no coincidence - in fact, $\mathscr{A}_{\mathrm{K}_{\mathrm{n}}}$ is the set of reflecting hyperplanes of the standard linear representation of the symmetric group (or: Coxeter group of type $A_{n-1}$ ), and the chambers are in bijection with the elements of the group... of which there are exactly $n!$.

Question. Can we treat in the same way the other finite Coxeter types?

## Exercises III

III. 1 Let $r$ be the rank function of a matroid on $E$ and suppose that $e$ is a loop of this matroid. Prove that $e$ is contained in every flat.
III. 2 Let $r$ be the rank function of a matroid on $E$ and suppose that $e$ is an isthmus of this matroid. Prove that the lattice of flats of $r_{/ e}$ and $r_{\backslash e}$ are isomorphic.
III. 3 Let $\mathrm{d}, \mathrm{n} \in \mathbb{N}, \mathrm{d} \leqslant \mathrm{n}$. The uniform matroid $\mathrm{U}_{\mathrm{m}, \mathrm{d}}$ of rank d on n elements has ground set $[\mathrm{n}]$ and rank function defined by

$$
r(A)= \begin{cases}|A| & \text { if }|\mathcal{A}| \leqslant d \\ d & \text { otherwise }\end{cases}
$$

Prove that the dual to $\mathrm{U}_{\mathrm{n}, \mathrm{d}}$ is $\mathrm{U}_{\mathrm{n}, \mathrm{n}-\mathrm{d}}$.
III. 4 Let $\mathscr{A}$ be the arrangement of four hyperplanes $\mathrm{H}_{1}, \mathrm{H}_{2}, \mathrm{H}_{3}, \mathrm{H}_{4}$ in $\mathbb{R}^{2}$ with equations:

$$
H_{1}: x+y=0, H_{2}: x=0, H_{3}: y=0, H_{4}: x-y=0
$$

Draw a picture of this arrangement and of the Hasse diagram of $\mathcal{L}(\mathscr{A})$. Compute the Tutte polynomial of the associated matroid and verify the formula for the number of regions of $\mathscr{A}$ in this case.
III. 5 Consider the graph


Draw the Hasse diagram of the lattice of flats of the associated matroid. Compute the Tutte polynomial and the number of regions of the associated graphic arrangement. Check it with a sketch of the arrangement (it is an arrangement in $\mathbb{R}^{4}$, but all hyperplanes contain the line $\ell$ in direction ( $1,1,1$ ), thus by quotienting by $\ell$ one can attempt a 3 -dimensional sketch.
III. 6 Let $\mathcal{L}$ be a finite geometric lattice. The Möbius function of $\mathcal{L}$ is the function $\mu: \mathcal{L} \rightarrow \mathbb{N}$ defined via the following recursion:

$$
\mu(\widehat{0})=1, \quad \text { and } \sum_{y \leqslant x} \mu(y)=0 \text { for all } x \in \mathcal{L}_{>\hat{0}} .
$$

- Compute $\mu$ for the geometric lattice $\mathcal{L}$ in Example 3.1.2 and find an interpretation of the sum $\sum_{x \in \mathcal{L}}|\mu(x)|$ in this case.
- Compute the values of $\mu$ for the poset of Example 3.3.2 and consider the numbers you obtain by summing those values over every ranklevel of the poset. Then compute the chromatic polynomial of the graph of the same example. What do you notice?
- Let now $\mathcal{L}$ be any geometric lattice. Prove that, for all $x \in \mathcal{L}$,

$$
\mu(x)=\sum_{\substack{S \subseteq \mathcal{A}(\mathcal{L}) \\ V S=x}}(-1)^{|S|}
$$

where the sum is over all sets $S$ of atoms of $\mathcal{L}$ whose join equals $x$. (Hint: prove that $\mu(x)$ and the expression on the r.-h.s. satisfy the same recursion).

- Let $r$ be the rank function of a matroid with ground set $E$ that has no loops. Show that

$$
T_{r}(1-t, 0)=(-1)^{r(E)} \sum_{F \in \mathcal{L}_{r}} \mu(F) t^{r(E)-r(F)} .
$$

[Hint: use the previous item in this exercise.]

- Prove that, for every arrangement of hyperplanes $\mathscr{A}$, the number of regions determined by $\mathscr{A}$ equals $\sum_{x \in \mathcal{L}(\mathscr{A})}|\mu(x)|$.
- Prove that for every graph $G$ on $n$ vertices the chromatic polynomial of $G$ equals

$$
\chi_{G}(t)=(-1)^{r(G)} t^{c(G)} \sum_{i=0}^{r(G)} \mu_{i}\left(\mathcal{L}_{r_{G}}\right) t^{i}
$$

where for a geometric lattice $\mathcal{L}$ with rank function $\rho$ we write $\mu_{i}(\mathcal{L})=\sum_{\rho(x)=i} \mu(x)$ for the sum of the values of $\mu$ on all elements of $\mathcal{L}$ of rank $i$.
Check this formula with the lattice of Exercise III.5.

## Chapter 4

## Signed graphs

### 4.1 Definition

A loopless signed graph is, roughly speaking, a graph with two kinds of edges ("standard" and "half-" edges) and a "sign" associated to every edge.

Formally we will think of a tuple

$$
\Sigma=(G, \sigma), \quad G=\left(V, E, h: E \rightarrow V, t: E_{r} \rightarrow V\right), \quad \sigma: E_{r} \rightarrow\{ \pm\}
$$

where $E_{r} \subseteq E$ is the set of all "regular" edges, while $E_{1 / 2}:=E \backslash E_{r}$ is the set of half-edges. We will only need to consider loopless graphs, thus we assume $|\mathrm{V}(\mathrm{e})|=2$ for all $e \in \mathrm{E}_{\mathrm{r}}$.

Note that $\sigma$ labels every edge with a sign + or - . Therefore every trail has a sign as well, obtained as the product of the signs of all its edges. A circuit is called balanced if its sign + , unbalanced otherwise. A connected component of a signed graph is called balanced if every circuit in it is balanced, unbalanced otherwise. For every $A \subseteq E$ we call $b(A)$ the number of balanced connected components of the restriction of the signed graph on the set $A$ (as in the unsigned case, this is the graph with same vertex set V but only the edges in A).

Remark 4.1.1. In particular, any cycle in a balanced component of a signed graph has positive sign. This is easy to see by noticing that a cycle decomposes as a disjoint union of circuits (each with positive sign) and "repeated" edges (and the product of two negatives is a positive).
Lemma 4.1.2. Let $p_{1}, p_{2}$ be two paths with the same endpoints in a balanced component of a signed graph. Then the product of the signs of the edges of $p_{1}$ equals the product of the signs of the edges of $\mathrm{p}_{2}$.
Proof. The sign of the union of $p_{1}$ and $p_{2}$ is by definition the product of the signs of the two paths. But this union is a cycle which must have positive sign by the preceding remark. The claim follows.

### 4.2 Coloring

Let $k \in \mathbb{N}$ and let $[[k]]:=\{-k,-k+1, \ldots, 0,1, \ldots, k\}$. A $k$-coloring of a loopless signed graph $\Sigma$ is an assignment $\gamma: \mathrm{V} \rightarrow[[k]]$ such that:
(i) $\gamma(h(e)) \neq \sigma(e) \gamma(t(e))$ for all $e \in E_{r}$,
(ii) $\gamma(h(e)) \neq 0$ for all $e \in E_{1 / 2}$.

As in the case of unsigned graphs, we call precoloring any function $\mathrm{V} \rightarrow[[\mathrm{k}]]$ and, given any $A \subseteq E$, we define the set of precoloring which fail to satisfy the constraints imposed by the edges in $A$ :

$$
\operatorname{Pre}_{\Sigma}(A):=\left\{\begin{array}{l|l}
\gamma: V \rightarrow[[k]] & \begin{array}{ll}
\gamma(h(e))=\sigma(e) \gamma(t(e)) & \text { for } e \in A \cap E_{r} \\
\gamma(h(e))=0 & \text { for } e \in E_{1 / 2} \cap A
\end{array}
\end{array}\right\}
$$

Then any $\gamma \in \operatorname{Pre}_{\Sigma}(A)$ must have constant value across every connected component of $\mathrm{G}[A]$. Moreover, notice that every unbalanced component of $\mathrm{G}[A]$ must contain an unbalanced cycle, say with sequence of edges-vertices $v_{0}, e_{1}, v_{1} \ldots, e_{m}, v_{m}=v_{0}$, and we have

$$
\gamma\left(v_{0}\right)=\sigma\left(e_{1}\right) \gamma\left(v_{1}\right)=\ldots=\sigma\left(e_{1}\right) \cdots \sigma\left(e_{\mathfrak{m}}\right) \gamma\left(v_{\mathfrak{m}}\right)=-\gamma\left(v_{0}\right),
$$

the last equality because $\nu_{\mathrm{m}}=\nu_{0}$ and the cycle is unbalanced. In particular, $\gamma\left(v_{0}\right)=0$ and thus $\gamma$ must be 0 across every unbalanced component of $A$.

Since the (constant) value of $\gamma$ on any balanced connected component of $G[A]$ can be chosen freely, we conclude that

$$
\left|\operatorname{Pre}_{\Sigma}(A)\right|=(2 k+1)^{b(A)}
$$

Just as in the case of unsigned graphs, applying inclusion-exclusion on the family of all subsets of $E$ we obtain that the function $\chi_{\Sigma}(t)$ that for odd $t$ counts the number of $\frac{(t-1)}{2}$-colorings of $\Sigma$ can be written

$$
\chi_{\Sigma}(\mathrm{t})=\sum_{A \subseteq E}(-1)^{|\mathrm{A}|} \mathrm{t}^{\mathrm{b}(\mathrm{~A})}
$$

and is therefore a polynomial.
Coup de théâtre: we'll show next that this is in fact, up to a factor, the characteristic polynomial of a matroid!

### 4.3 Matroids

Definition 4.3.1. Let $\Sigma$ be a (loopless) signed graph on $n$ vertices. Define a function

$$
\mathrm{r}_{\Sigma}: 2^{\mathrm{E}} \rightarrow \mathbb{N}, \quad A \mapsto \mathrm{r}_{\Sigma}(A):=\mathrm{n}-\mathrm{b}(A)
$$

Theorem 4.3.2. The function $r_{\Sigma}$ is the rank function of a matroid on E .
Proof. We will verify the rank axioms, but first a general observation.
G.O. When adding an edge to a signed graph $\Sigma$, the number of balanced components can drop in two ways: either ome balanced component becomes unbalanced (i.e., by adding a half-edge, by adding an edge that creates an unbalanced cycle or by joining a balanced component with an unbalanced one) or two balanced components are joined via the new edge to one bigger balanced component. In any case the number of balanced components either stays the same or it decreases at most by one, i.e,

$$
\begin{equation*}
r_{\Sigma}(A) \leqslant r_{\Sigma}(A \cup\{e\}) \leqslant r_{\Sigma}(A)+1 . \tag{4.1}
\end{equation*}
$$

$r 0$ Clearly $r_{\Sigma}$ is never negative. In order to prove that $r_{\Sigma}(\mathcal{A}) \leqslant|\mathcal{A}|$ for all $A \subseteq E$, start with the graph $\Sigma_{0}$ on $V$ and with no edges: here every vertex is its own (balanced!) component. Therefore $r_{\Sigma}(\emptyset)=n-n=0$. Now add the edges in $A$ to $\Sigma_{0}$ one by one. By the G.O., the number of balanced components decreases at most by one at each step, thus the rank increases at most by one at each step (second inequality in Equation (4.1). Thus, $\mathrm{r}_{\Sigma}(\mathcal{A}) \leqslant|\mathcal{A}|$ as required.
$r 1$ Obviously it is enough to show that $r_{\Sigma}(\mathcal{A}) \leqslant r_{\Sigma}(A \cup\{e\})$ for all $A \subseteq E$ and all $e$. But this is implied directly from the G.O. above (first inequality of Equation (4.1)).
r2 Here it is enough to show axiom $r 2$ for the case where $X=A \cup\{e\}$ and $Y=A \cup\{f\}$, for any given $A \subseteq E$ and $e, f \in E$ (see Exercise IV.1). In this case we have to show

$$
r_{\Sigma}(A \cup\{e, f\})-r_{\Sigma}(A \cup\{f\}) \leqslant r_{\Sigma}(A \cup\{e\})-r_{\Sigma}(A) .
$$

Now, if the right-hand side equals 1 then the inequality holds trivially, since by ( r 1 ) the left-hand side is at most 1 .
Suppose then that $r_{\Sigma}(A \cup\{e\})=r_{\Sigma}(A)-$ in particular, if $e$ has one end in a balanced component, then it must be a regular edge with both ends in the same balanced component of $A$. By way of contradiction, suppose now that $r_{\Sigma}(A \cup\{e, f\})-r_{\Sigma}(A \cup\{f\})=1$. The latter means that the number of balanced components changes by adding $e$ to $A \cup\{f\}$. By the G.O., this means that either
(1) e connects two balanced components of $A \cup\{f\}$ - a contradiction because both ends of $e$ are in the same component of $A$, and $a$ fortiori of $A \cup f$, or
(2) $e$ "closes" an unbalanced circuit $C$ in a balanced component of $A \cup\{f\}$. But this can only happen if $C$ also contains $f$. Now, if $f$
connects two balanced components of $A$, then $e$ would, too - a contradiction. So the two ends of $f$ lie in the same balanced component of $A$ and (by the connectedness of this component) are connected by a path $p$ whose sign agrees with that of $f$ (by Lemma 4.1.2. Now we can replace $f$ by $p$ in $C$ and obtain a closed, unbalanced cycle containing $e$ in $A \cup e$, thus $e$ "closes" an unbalanced cycle in a balanced component of $\Sigma[A]$, contradicting the assumption.

In all possible cases we reached a contradiction and the desired inequality is proved.

### 4.4 Arrangements from signed graphs

We now show that the matroid defined by $r_{\Sigma}$ is in fact obtainable as the matroid of linear dependencies of a set of vectors.

Definition 4.4.1. Let $\Sigma$ be a loopless signed graph on the vertex set $V$ with set of edges $E$. For each edge $e \in E$ define a vector $m_{e} \in \mathbb{R}^{V}$ as follows:

$$
\left(m_{e}\right)_{v}:= \begin{cases}1 & \text { if } v=h(e) \\ \sigma(e)(-1) & \text { if } v=\mathrm{t}(e) \\ 0 & \text { otherwise }\end{cases}
$$

In particular, if $e \in E_{1 / 2}$ the vector $m_{e}$ has only one nonzero component. Let $M_{\Sigma}$ the $V \times E$ matrix with columns $m_{e}$ and, for every $A \subseteq E$, let $M_{\Sigma}[A]$ denote the matrix constructed with only the columns indexed by $A$.

The following theorem states that $\mathrm{r}_{\Sigma}$ is the rank function of the matroid obtained as in Theorem 2.2 .15 from the vectors $m_{e}$ where $e$ ranges over $E$.

Theorem 4.4.2. For every $A \subseteq E, r_{\Sigma}(A)=\operatorname{rank} M_{\Sigma}[A]$.
The proof will use the following lemma.
Lemma 4.4.3. Let $\Sigma$ be a connected signed graph on n vertices. Then $\operatorname{rank} \mathrm{M}_{\Sigma}=$ $\mathrm{n}-1$ if $\Sigma$ is balanced, and rank $M_{\Sigma}=\mathrm{n}$ otherwise.

Proof. Let T be a set of regular (non-half) edges that is a spanning tree of $\Sigma$. We can suppose without loss of generality that, if $\Sigma$ has half-edges, then one half-edge is attached to a leaf $v_{0}$ of $T$. Then $T$ has size $n-1$ and we can order the vertices of $\Sigma$ beginning with $v_{0}$, and continuing in such a way that, for all $i$, the subgraph of $T$ induced by the vertices $v_{0}, \ldots, v_{i}$ is connected. Then the matrix $M_{\Sigma}$ with the rows ordered this way can be written as follows (where the first row is only present in case $v_{0}$ has a half-edges attached to it)

$$
\left(\begin{array}{ccccccccc}
(1) & \mp 1 & * & * & \cdots & * & * & \cdots & * \\
(0) & \pm 1 & * & & & * & & * \\
(0) & 0 & \pm 1 & * & & * & & & * \\
(0) & 0 & 0 & \pm 1 & \ddots & * & & * \\
\vdots & \vdots & & \ddots & \ddots & * & & & * \\
(0) & 0 & & \cdots & 0 & \pm 1 & * & \cdots & *
\end{array}\right)
$$

and it has clearly rank at least $n-1$. Since $M_{\Sigma}$ has $n$ rows, its rank is at most $n$. Now, it is enough to show the following, which we leave as an exercise.

- If $\Sigma$ is balanced, then the parenthesized column is not present and every $m_{e}$, for $e \notin T$, is a linear combination of the set $\left\{m_{t} \mid t \in T\right\}$. (Hint: for every such $e$ consider the (unique) circuit contained in $\mathrm{T} \cup\{e\}$.
- If $\Sigma$ is not balanced, then either it has an half-edge, in which case the parenthesized column is present and so the matix has rank $n$, or there is no half-edge. In this last case, there is an $e \in E \backslash T$ that is linearly independent from $\left\{m_{t} \mid t \in T\right\}$.

Proof of Theorem 4.4.2 First notice that $M_{\Sigma}[A]$ is the matrix associated to the restriction of $\Sigma$ to the edge-set $A$. Therefore, it is enough to prove the claim for $A=E$, i.e., we have to prove that $\operatorname{rank} M_{\Sigma}=r_{\Sigma}(E)$. Moreover, if $\Sigma$ is disconnected we can let $V_{1}, \ldots, V_{k}$ and $A_{1}, \ldots, A_{k}$ be the vertex sets, resp. the edge sets of the connected components, and we see that the matrix $M_{\Sigma}$ has block-diagonal form

$$
\begin{gathered}
\\
\mathrm{V}_{1} \\
\mathrm{~V}_{2} \\
\vdots \\
\mathrm{~V}_{\mathrm{k}}
\end{gathered}\left(\begin{array}{cccc}
A_{1} & A_{2} & \cdots & A_{k} \\
\mathrm{M}_{\Sigma_{1}} & 0 & 0 & 0 \\
0 & M_{\Sigma_{2}} & 0 & 0 \\
0 & 0 & \ddots & 0 \\
0 & 0 & 0 & M_{\Sigma_{k}}
\end{array}\right)
$$

showing that $\operatorname{rank} M_{\Sigma}=\Sigma_{i} \operatorname{rank} M_{\Sigma_{i}}$ where $\Sigma_{i}$ is the signed graph on the vertex set $V_{i}$ with (signed) edges $A_{i}$. On the other hand, writing $n_{i}=\left|V_{i}\right|$ and setting $b_{i}=1$ if $\Sigma_{i}$ is balanced and $b_{i}=0$ otherwise, we have that

$$
\sum_{i} r_{\Sigma_{i}}\left(A_{i}\right)=\sum_{i}\left(n_{i}-b\left(A_{i}\right)\right)=\sum_{i}\left(n_{i}-b_{i}\right)=n-b(E)=r_{\Sigma}(E) .
$$

Therefore, it is enough to prove the claim for $\Sigma$ connected. This has been done precisely in Lemma 4.4.3.

The Tutte polynomial of the matroid $r_{\Sigma}$ is of course related to the chromatic polynomial $\chi_{\Sigma}$.

Lemma 4.4.4. Let $\Sigma$ be a loopless signed graph. Then

$$
\chi_{\Sigma}(t)=(-1)^{n} t^{b(E)} T_{r_{\Sigma}}(1-t, 0)
$$

Proof. Direct computation:

$$
\begin{aligned}
T_{r_{\Sigma}}(1-t, 0)= & \sum_{A \subseteq E}(-t)^{(n-b(E))-(n-b(A))}(-1)^{|A|-(n-b(A))} \\
& =(-1)^{n} t^{-b(E)} \sum_{A \subseteq E}(-1)^{|A|} t^{b(A)}
\end{aligned}
$$

### 4.5 Arrangements associated to root systems

For every fixed dimension $d$ there are only a finite number of finite (irreducible) groups of isometries of $\mathbb{R}^{\mathrm{d}}$ generated by reflections. Those are the (irreducible) finite Coxeter groups, and their theory is a mainstay of algebra and geometry since its bases were laid by H.S.M. Coxeter ${ }^{1}$ in the mid- 20. Century.

The classification of (irreducible) finite Coxeter groups identifies in every dimension one representative of each of four infinite families (indexed by $A_{n}, B_{n}, C_{n}, D_{n}$ ) and, in some (low) dimensions, a finite number of sporadic cases. One way to specify one of these groups, say $G$, is to give all isometries $f \in G$ that are reflections $s^{2}$ - by definition, these generate the group G. Now, a (orthogonal) reflection $f$ is fully specified by its fixed space $\operatorname{fix}(f)=\left\{x \in \mathbb{R}^{d} \mid f(x)=x\right\}$. Note that if $f$ is a reflection, then fix(f) is a hyperplane in $\mathbb{R}^{d}$. Therefore, $G$ is given by the associated "arrangement of reflection hyperplanes". The arrangements that arise this way are explicitly listed, e.g., by means of a choice of their normal vectors. The set of normal vectors is usually referred to as the set of (positive) roots of the corresponding Coxeter group ${ }^{3}$, (The word "positive" is there because root systems contain the opposite of each root, but opposite vectors define the same hyperplane.) Notice that the group acts freely and transitively on the set of regions of the corresponding arrangement, so there are as many regions as there are group elements.

We list here a set of normal vectors for the hyperplane arrangements of the four infinite families. We use the notation $\epsilon_{i}$ to denote the $i$-th standard basis

[^2]vector in $\mathbb{R}^{n}$, so $\epsilon_{i}$ has all components equal to 0 except the $i$-th component being equal to 1 .
$\left.A_{n-1}\right|^{4}$. The vectors $\epsilon_{i}-\epsilon_{j}$ in $\mathbb{R}^{n}$ for $n \geqslant i>j \geqslant 1$
$B_{n}$ : The vectors $\epsilon_{i}-\epsilon_{j}$ as well as $\epsilon_{i}+\epsilon_{j}$ and $\epsilon_{i}$ in $\mathbb{R}^{n}$ for $n \geqslant i>j \geqslant 1$
$C_{n}$ : The vectors $\epsilon_{i}-\epsilon_{j}$ as well as $\epsilon_{i}+\epsilon_{j}$ and $2 \epsilon_{i}$ in $\mathbb{R}^{n}$ for $n \geqslant i>j \geqslant 1$
$D_{n}$ : The vectors $\epsilon_{i}-\epsilon_{\mathfrak{j}}$ and $\epsilon_{\mathfrak{i}}+\epsilon_{\mathfrak{j}}$ in $\mathbb{R}^{n}$ for $\mathfrak{n} \geqslant \mathfrak{i}>\boldsymbol{j} \geqslant 1$
We first notice that the arrangement of type $B$ and $C$ are equal (since $\epsilon_{i}$ and $2 \epsilon_{i}$ define the same hyperplane). Then, we notice that the set of vectors in types $A, B, D$ can be obtained as vectors of the type $m_{e}$ for edges of suitable signed graphs, that we will call $\Sigma\left(A_{n-1}\right), \Sigma\left(B_{n}\right), \Sigma\left(D_{n}\right)$ and we describe below.
$\Sigma\left(A_{n-1}\right)$ is the graph $K_{n}^{+}$obtained from the complete graph on $n$ vertices by assigning a positive sign to each vertex.
$\Sigma\left(B_{n}\right)$ is the graph on $n$ vertices with a pair of opposite-signed edges between every pair of vertices, and a half-edge at every vertex.
$\Sigma\left(D_{n}\right)$ is the graph on $n$ vertices with a pair of opposite-signed edges between every pair of vertices.

Now the number of regions in each such arrangement is, by the general principle, given by the Tutte polynomial of the associated matroid $\mathrm{r}_{\Sigma}$ - and, by Lemma 4.4.4 as

$$
\mathrm{T}_{\mathrm{r}_{\Sigma}}(2,0)=(-1)^{\mathrm{n}-\mathrm{b}(\mathrm{E})} \chi_{\Sigma}(-1)
$$

In order to compute this number it is now enough to compute the polynomial $\chi_{\Sigma}(\mathrm{t})$ for the signed graphs associated to the different root systems.
$A_{n-1}$ The polynomial is $\chi_{\Sigma\left(A_{n-1}\right)}(t)=\prod_{i=0}^{n-1}(t-i)$
The constraints are as in the (unsigned) complete graph: given $t$ colors we can assign any of them to the first vertex, only $t-1$ to the second, and so on.
Thus the number of regions of the arrangement is $n$ ! (this agrees with the fact that the reflection group of type $A_{n-1}$ is the symmetric group $S_{n}$ of all permutations of $n$ elements - and of those there are notoriously n !.
$B_{n}$ The polynomial is $\chi_{\Sigma\left(B_{n}\right)}(t)=\prod_{i=1}^{n}(t-2 i+1)$
Here the constraints are as follows: between any two vertices we have two signed edges, that forbid the two vertices to have the same or the opposite color. Moreover the half-edges forbid the color 0
everywhere. So that for (odd) $t$ signed colors we have $(t-1)$ possibilities for vertex $1,(t-3)$ for vertex 2 and so on.

Thus the number of regions of the arrangement (and hence the order of the reflection group) is $2 n(2 n-2)(2 n-4) \cdots=2^{n} n$ !
$D_{n}$ The polynomial is $\chi_{\Sigma\left(D_{n}\right)}(t)=(t-n+1) \prod_{i=1}^{n-1}(t-2 i+1)$
Here the constraints are as in $B_{n}$ but 0 may be allowed: between any two vertices we have two signed edges, that forbid the two vertices to have the same or the opposite color. Therefore we can sum over the possibilities: either 0 is not used, or it is used once, and there are $n$ possibilities to assign it, while all other vertices must be colored differently, and according to the rules of $B_{n}$ :

$$
\chi_{\Sigma\left(D_{n}\right)}(t)=\prod_{i=1}^{n}(t-2 i+1)+n \prod_{i=1}^{n-1}(t-2 i+1)
$$

The number of regions of the arrangement (and hence the order of the reflection group) is $n(2 n-2)(2 n-4) \cdots=2^{n-1} n$ !

## Exercises IV

IV. 1 Prove that axiom (r2) in the definition of a matroid rank function is equivalent to the statement that

$$
r(A \cup\{e\})+r(A \cup\{f\}) \geqslant r(A)+r(A \cup\{e, f\}) \quad \text { for all } A \subseteq E, e, f \in E
$$

IV. 2 Let $\Sigma$ be a connected, loopless signed graph and let $T$ be any spanning tree of $\Sigma$ (i.e., T is an acyclic set of regular edges that are incident to every vertex). Prove:

- For every regular $e \in E_{r} \backslash T$ there is a unique circuit contained in $T \cup\{e\}$. We call this circuit $C(T, e)$.
- If $\Sigma$ is unbalanced and has no half-edges, then there is an edge $e \in E$ such that $C(T, e)$ is unbalanced.


## Chapter 5

## A polyhedral interlude

### 5.1 Motivation: a loose end about flows on graphs

Recall that an integral $k$-flow on a graph $G=(V, E, h, t)$ is an assignment $\mathrm{f}: \mathrm{E} \rightarrow\{-\mathrm{k}, \ldots, \mathrm{k}\}$ such that

$$
\begin{equation*}
\sum_{h(e)=v} f(e)=\sum_{t(e)=v} f(e) \quad \text { (as an equation in } \mathbb{Z} \text { ) for all } v \in V \text {. } \tag{5.1}
\end{equation*}
$$

The task we set ourselves in Section 1.3 was to determine the function $\varphi_{\mathrm{G}}^{\mathrm{N}}$ that counts the number of nowhere-zero integral k-flows on $G$.

Definition 5.1.1. Call an integer $k$ - flow on $G$ positive if $f(E)>0$. Let $I_{G}(k)$ be the number of positive integral $k$-flows on $G$.

We can make any flow positive by "flipping" edges, i.e., performing the following operation: Given a graph $G=(V, E, h, t)$ and $A \subseteq E$, let

$$
G^{\neg A}:=\left(V, E, h^{\prime}, t^{\prime}\right), \quad \text { where } h^{\prime}(a)=t(a) \text { and } t^{\prime}(a)=h(a) \text { for all } a \in A .
$$

Intuitively, as the word says, we "flip" head and tail of all edges in A.
Lemma 5.1.2. If f is a nowhere-zero flow on G , let $\mathrm{A}:=\{e \in \mathrm{E} \mid \mathrm{f}(e)<0\}$. Then the assignment

$$
f^{\prime}(e):= \begin{cases}-f(e) & \text { if } e \in A \\ f(e) & \text { otherwise }\end{cases}
$$

is a positive integer k -flow on $\mathrm{G}^{\checkmark \mathrm{A}}$.
Proof. We only have to show that $\mathrm{f}^{\prime}$ is an integer k -flow.
In particular, for every nowhere-zero integral flow there is exactly one set $A \subseteq E$ such that $f^{\prime}$ is a positive integral flow on $G^{-A}$. This shows that

$$
\varphi_{\mathrm{G}}^{\mathbb{N}}(\mathrm{k})=\sum_{A \subseteq E} \mathrm{I}_{\mathrm{G}^{\neg \mathrm{A}}}(\mathrm{k})
$$

and therefore we are led to consider the function $I_{G}(k)$.
Recall the incidence matrix of the graph $G$, a matrix $A \in \mathbb{R}^{V \times E}$ whose $v$-th component in the e-th column is

$$
\left(A_{e}\right)_{v}= \begin{cases}1 & \text { if } v=h(e) \\ -1 & \text { if } v=\mathrm{t}(e) \\ 0 & \text { otherwise }\end{cases}
$$

Now a positive k-flow on $G$ is any $f \in \mathbb{Z}^{E}$ such that

$$
\begin{equation*}
\text { Af }=0, \quad \mathbf{0}<\mathrm{f}<\mathbf{k} \tag{5.2}
\end{equation*}
$$

where $\mathbf{0}=(0, \ldots, 0), \mathbf{k}=(k, \ldots, k)$ and the inequalities are understood componentwise.

Now consider the set

$$
\mathrm{U}:=\left\{x \in \mathbb{R}^{\mathrm{E}} \mid A x=0, \mathbf{0}<x<\mathbf{1}\right\}
$$

and its closure $\overline{\mathrm{U}}:=\left\{x \in \mathbb{R}^{\mathrm{E}} \mid A x=0, \mathbf{0} \leqslant x \leqslant \mathbf{1}\right\}$.
Then we have

$$
\mathrm{I}_{\mathrm{G}}(\mathrm{k})=\# \mathbb{Z}^{\mathrm{E}} \cap \mathrm{kU}
$$

The fact that $\overline{\mathrm{U}}$ is a convex polytope with integer vertices opens up the possibility to compute $\mathrm{I}_{\mathrm{G}}(\mathrm{k})$ using Ehrhart theory.

### 5.2 Polyhedra and polytopes

Recall that, when writing inequalities between vectors we intend the inequalities to be taken componentwise. Moreover, in this section $A$ will denote any $n \times m$ matrix with real entries.

Definition 5.2.1. A polyhedron is any subset of $\mathbb{R}^{n}$ of the form

$$
\mathbf{P}=\left\{x \in \mathbb{R}^{n} \mid A x \leqslant b\right\}
$$

for some matrix $A \in \mathbb{R}^{m \times n}$ and some vector $b \in \mathbb{R}^{m}$.
Example 5.2.2. In Figure 1 we depict two instances of polyhedra in $\mathbb{R}^{2}$.
Definition 5.2 .3 . A polytope is any subset of $\mathbb{R}^{n}$ of the form

$$
P=\operatorname{conv}\left\{x_{1}, \ldots, x_{k}\right\}:=\left\{\sum \lambda_{i} x_{i} \mid \lambda_{i} \in \mathbb{R}_{\geqslant 0}, \sum \lambda_{i}=1\right\}
$$

for some finite subset $\left\{x_{1}, \ldots, x_{k}\right\} \subseteq \mathbb{R}^{n}$.


(a) The (unbounded) polyhedron
$\left\{\left[\begin{array}{cc}1 & 2 \\ -1 & -1\end{array}\right]\left[\begin{array}{l}x \\ y\end{array}\right] \leqslant\left[\begin{array}{c}8 \\ -2\end{array}\right]\right\}$
(b) The (bounded) polyhedron

$$
\left\{\left[\begin{array}{cc}
1 & 2 \\
-1 & -1 \\
-1 & 1 \\
1 & -1
\end{array}\right]\left[\begin{array}{l}
x \\
y
\end{array}\right] \leqslant\left[\begin{array}{c}
8 \\
-2 \\
1 \\
-1
\end{array}\right]\right\}
$$

Figure 1


Figure 2: The polytope conv $\left\{(2,3),\left(\frac{5}{2}, \frac{3}{2}\right)\right\}$

Example 5.2.4. In Figure 2 we depict a polytope in $\mathbb{R}^{2}$, identical to the bounded polyhedron of Figure 1 b .

The goal of this section will be to prove the following fundamental theorem of polytope theory.

Theorem 5.2.5. A subset of $\mathbb{R}^{n}$ is a polytope if and only if it is a bounded polyhedron.
The next statement is our stepping stone towards proving Theorem 5.2.5

Theorem 5.2.6. Let $a_{1}, \ldots, a_{m}, v \in \mathbb{R}^{n}$. Then either
(1) $v=\lambda_{1} a_{1}+\ldots+\lambda_{m} a_{m}$ with $\lambda_{i} \geqslant 0$,or
(2) there is $c \in \mathbb{R}^{n}$ such that the hyperplane $c^{\perp}$ contains $t-1$ linearly independent vectors from the set $\left\{a_{1}, \ldots, a_{m}\right\}$ such that $c \cdot v<0$ and $c \cdot a_{1}, \ldots, c \cdot a_{m} \geqslant 0$, where $t=\operatorname{dim}\left\langle a_{1}, \ldots, a_{m}, v\right\rangle$.

Proof. First, notice that (1) and (2) are mutually exclusive: if both hold, then we would have $0>c \cdot v=\lambda_{1} c \cdot a_{1}+\ldots+\lambda_{m} c \cdot a_{m} \geqslant 0$, a contradiction. Thus, it will suffice to show that in any case at least one of the options holds.

Let $V:=\left\langle a_{1}, \ldots, a_{m}\right\rangle$. If $v \notin \mathrm{~V}$, then $\mathrm{V} \subsetneq \mathbb{R}^{n}$ is a proper subspace, and thus it is contained in a hyperplane $H$. Choose c to be any normal vector to $H$. We have $c \cdot v \neq 0$ and so, after possibly switching sign of $c$, we have $c \cdot v>0$. Moreover, as $V \subseteq H$, we also have $c \cdot a_{1}=\ldots=c \cdot a_{m}=0$, so case (2) holds.

Assume from now $v \in \mathrm{~V}$ and, without loss of generality, $\mathrm{V}=\mathbb{R}^{n}$, so that $t=n$. Choose then a linearly independent set $\left\{a_{i_{1}}, \ldots, a_{i_{n}}\right\}$ among the $a_{i}$ and set

$$
D_{0}:=\left\{i_{1}, \ldots, i_{n}\right\}
$$

and for all $i \geqslant 0$ iterate the following procedure.
(i) Consider the (unique) expression $v=\sum_{j \in D_{i}} \lambda_{j} a_{j}$. If $\lambda_{j} \geqslant 0$ for all $j \in D$, we are in case (1) and we stop.
(ii) Otherwise, choose $h:=\min \left\{j \in D \mid \lambda_{j}<0\right\}$ and let $H=\left\langle a_{j} \mid j \in D \backslash\{h\}\right\rangle$, a hyperplane of which we can choose a normal $c$ with $c \cdot a_{h}=1$, so that $c \cdot v=\lambda_{h}<0$.
(iii) If $C \cdot a_{i} \geqslant 0$ for all $i=1, \ldots, m$, we are in case (2) and we can stop.
(iv) Otherwise, let $s:=\min \left\{j \in[m] \mid c \cdot a_{j}<0\right\}$, set

$$
\mathrm{D}_{\mathfrak{i}+1}:=\left(\mathrm{D}_{\mathfrak{i}} \backslash\{\mathrm{h}\}\right) \cup\{\mathrm{s}\}
$$

and repeat.
We must prove that this iteration terminates; by way of contradiction, assume that this is not the case. Then, since there is only a finite number of options for choosing $D_{i} \subseteq[m]$, we must have $D_{k}=D_{l}$ for some $k<l$.
Now at every step $i$ between $k$ and $l$ one element is removed from (and one is added to) $D_{i}$. Let $r$ be the maximum among all elements that are removed at some step $i, k \leqslant i<l$ and let $p$ be a step where this happens (so $D_{p+1}=$ $D_{p} \backslash\{r\} \cup\{j\}$ for some $\mathfrak{j}$ ). Moreover, since $D_{k}=D_{l}$ there must be a $q \neq p$, $\mathrm{k} \leqslant \mathrm{q}<\mathrm{l}$, such that r is added to $\mathrm{D}_{\mathrm{q}}$ at the q -th step. In particular,

$$
\begin{equation*}
\mathrm{D}_{\mathrm{p}} \cap\{\mathrm{r}+1, \ldots, \mathrm{~m}\}=\mathrm{D}_{\mathrm{q}} \cap\{\mathrm{r}+1, \ldots, \mathrm{~m}\} \tag{5.3}
\end{equation*}
$$

Letting $c_{q}$ be the vector from item (ii) in the q-th iteration and considering the (unique) expression $v=\sum_{j \in D_{p}} \lambda_{j} a_{j}$ we reach the desired contradiction:

$$
0>c_{q} \cdot v=c_{q} \cdot\left(\sum_{j \in D_{p}} \lambda_{j} a_{j}\right)=\sum_{j \in D_{p}} \lambda_{j} c_{q} \cdot a_{j}>0
$$

the first inequality is by the choice of $\mathrm{c}_{\mathrm{q}}$ in item (ii). The second inequality holds because, in fact, every summand is nonnegative and the $r$-th is positive. To prove this, remember that

* since $r$ is removed from $D_{p}$, by (ii) $r=\min \left\{j \in D_{p} \mid \lambda_{j}<0\right\}$, and
$\star$ since $r$ is added to $D_{q}$ by (iv) $r=\min \left\{j \in[m] \mid c_{q} \cdot a_{j}<0\right\}$.
Now let $j \in D_{p}$ :
If $j<r$ then $\lambda_{j} \geqslant 0$ by $*$ and $c_{q} \cdot a_{j} \geqslant 0$ by $\star$, thus $\lambda_{j} c_{q} \cdot a_{j} \geqslant 0$.
If $j=r$ then $\lambda_{j}<0$ by $*$ and $c_{q} \cdot a_{j}<0$ by $\star$, thus $\lambda_{j} c_{q} \cdot a_{j}>0$.
If $\mathfrak{j}>r$ then $j \in D_{q}$ by Equation (5.3) and since the element removed from $D_{q}$ is smaller than $r$ (by maximality in the choice of $r$ ), then the definition of $c_{q}$ implies $c_{q} \cdot a_{j}=0$.

Let us review the two alternatives in the claim of Theorem 5.2.6. Claim (1) can be rephrased by saying that " $v$ is in the convex cone generated by $a_{1}, \ldots, a_{m}$. In fact, a convex cone is any $C \subseteq \mathbb{R}^{n}$ such that $\lambda x+\mu y \in C$ for all $x, y \in C$ and all $\lambda, \mu \geqslant 0$, and the smallest convex cone containing a given set $x_{1}, \ldots x_{m} \in \mathbb{R}^{n}$ is called the "cone generated by $x_{1}, \ldots x_{m}$ and written

$$
\operatorname{cone}\left\{x_{1}, \ldots, x_{m}\right\}:=\left\{\lambda_{1} x_{1}+\ldots+\lambda_{m} x_{m} \mid \lambda_{i} \geqslant 0\right\}
$$

and any convex cone arising this way is called finitely generated.
On the other hand we call polyhedral cone any nonempty $C \subseteq \mathbb{R}^{n}$ of the form

$$
\begin{equation*}
C=\left\{x \in \mathbb{R}^{n} \mid A x \leqslant 0\right\} \tag{5.4}
\end{equation*}
$$

for some matrix $A$.
Example 5.2.7. The cone in Figure 3 c . is generated, e.g., by the vectors $\left[\begin{array}{c}2 \\ -1\end{array}\right]$ and $\left[\begin{array}{c}1 \\ -1\end{array}\right]$.

The alternative expressed in Theorem 5.2.6 then can be understood to express the fact that any point $v$ that is not in the convex cone generated by $a_{1}, \ldots, a_{m}$ is also not in the smallest polyhedral cone containing the $a_{i} s$. In this vein, and more precisely, we have the following.

Proposition 5.2.8. A convex cone is polyhedral if and only if it is finitely generated.

## Proof.

Finitely generated $\Rightarrow$ polyhedral. Let $X:=\left\{x_{1}, \ldots, x_{m}\right\} \subset \mathbb{R}^{n}$ and assume without loss of generality that $x_{1}, \ldots, x_{n} \operatorname{span} \mathbb{R}^{n}$. Every nonzero vector $c \in \mathbb{R}^{n}$ defines a halfspace $\{x \mid c \cdot x \leqslant 0\}$ bounded by the hyperplane $\{x \mid c \cdot x=0\}$, and we consider the set $\mathcal{H}$ of all halfspaces containing $X$ and whose bounding hyperplanes are spanned by $n-1$ linearly independent elements among the $x_{i}$. For each halfspace $H \in \mathcal{H}$ choose a defining vector $\mathrm{c}_{\mathrm{H}}$. By Theorem 5.2.6. we have

$$
\operatorname{cone}\left\{\mathrm{x}_{1}, \ldots, \mathrm{x}_{\mathrm{m}}\right\}=\bigcap \mathcal{H}
$$

and since there are finitely many elements in $\mathcal{H}$, the cone is polyhedral (with respect to the matrix whose rows are the vectors $\mathrm{c}_{\mathrm{H}}, \mathrm{H} \in \mathcal{H}$ ).

Polyhedral $\Rightarrow$ finitely generated. Consider a polyhedral cone $C:=\{A x \leqslant 0\}$ and call $a_{1}, \ldots, a_{m}$ the rows of the matrix $A$. Then the cone generated by the $a_{1}, \ldots, a_{m}$ is polyhedral by the previous item ("the other direction"), and so there are $b_{1}, \ldots, b_{t} \in \mathbb{R}^{n}$ such that

$$
\operatorname{cone}\left\{a_{1}, \ldots, a_{m}\right\}=\{B x \leqslant 0\}
$$

for the matrix $B$ whose rows are the $b_{i} s$.
Claim. cone $\left\{\mathrm{b}_{1}, \ldots, \mathrm{~b}_{\mathrm{t}}\right\}=\mathrm{C}$, so C is finitely generated.
Proof. We call $K:=$ cone $\left\{\mathrm{b}_{1}, \ldots, \mathrm{~b}_{\mathrm{t}}\right\}$ and split the proof in two parts. First we prove $K \subseteq C$. In fact, by definition $a_{i} \in K$ and so $b_{j} \cdot a_{i} \leqslant 0$ for all $i, j$. Thus $A b_{j} \leqslant 0$, and consequently $b_{j} \in C$, for all $j$.
Now we prove $C \subseteq K$. If there was a $y \in C \backslash K$, then by Theorem 5.2.6 there would be $c$ with $c \cdot y<0$ and $c \cdot b_{i} \geqslant 0$ for all $i$. But then we would have

$$
(-c) \cdot b_{i} \leqslant 0 \text { for all } i \text {, so }-c \in \operatorname{cone}\left\{a_{1}, \ldots, a_{m}\right\}
$$

and we could write $-c=\sum_{i} \lambda_{i} a_{i}$ with all $\lambda_{i} \geqslant 0$. In particular, for every $x \in C$ we'd have

$$
(-c) \cdot x=\left(\sum_{i} \lambda_{i} a_{i}\right) x=\sum \lambda_{i} a_{i} \cdot x \leqslant 0
$$

the last inequality because $\lambda_{i} \geqslant 0$ and, since $x \in C, a_{i} \cdot x \leqslant 0$. In particular, since $y \in C$ we would have $c \cdot y \geqslant 0$, contradicting $c \cdot y<0$.

The last preparation for the proof of Theorem 5.2 .5 is a statement characterizing (possibly unbounded) polyhedra. Recall that we add subsets of $\mathbb{R}^{n}$ pointwise - i.e., for $A, B \in \mathbb{R}^{n}$ we write $A+B=\{a+b \mid a \in A, b \in B\}$.

Proposition 5.2.9. Let $\mathrm{P} \subseteq \mathbb{R}^{n}$. The following are equivalent.
(1) P is a polyhedron,
(2) $\mathrm{P}=\mathrm{Q}+\mathrm{C}$ for a polytope Q and a polyhedral cone C .


Figure 3: The polyhedron $P$ in (a) can be written as $P=Q+C$ for the polytope $Q$ from (b) and the polyhedral cone $C$ from (c).

Proof of Proposition 5.2 .9
$(2) \Rightarrow(1)$ Suppose $P=Q+C$ with $Q$ a polytope, say $Q=\operatorname{conv}\left\{x_{1}, \ldots, x_{m}\right\}$, and $C$ a polyhedral cone, say $C=$ cone $\left\{y_{1}, \ldots, y_{t}\right\}$. Then, $v \in P$ if and only if

$$
\begin{equation*}
\binom{v}{1} \in \text { cone }\left\{\binom{x_{1}}{1}, \ldots\binom{x_{m}}{1},\binom{y_{1}}{0}, \ldots\binom{y_{t}}{0}\right\} \tag{5.5}
\end{equation*}
$$

Call K the finitely generated cone in Equation (5.5). By Proposition 5.2.8. K is polyhedral, i.e., there is a matrix $\mathcal{A}^{\prime}$ such that

$$
K=\left\{A^{\prime}\binom{\nu}{\lambda} \leqslant 0\right\}=\{A x+\lambda \cdot b \leqslant 0\}
$$

where $b$ is the last column of $A^{\prime}$ and $A$ is $A^{\prime}$ without the column $b$.
Therefore, $v \in \mathrm{P}$ if and only if $A v \leqslant-b$, and thus $\mathrm{P}=\{A x \leqslant-b\}$ is a polyhedron.
$(1) \Rightarrow(2)$ Suppose now that $P=\{A x \leqslant b\}$ for some matrix $A$ and vector $b$. Consider the polyhedral cone

$$
K:=\left\{\left[\begin{array}{c|c}
0 & -1 \\
\hline A & -\mathrm{b}
\end{array}\right]\binom{x}{\mu} \leqslant 0\right\}=\left\{\left.\binom{x}{\mu} \right\rvert\, A x-\mu \mathrm{b} \leqslant 0, \mu \geqslant 0\right\}
$$

By Proposition 5.2 .8 this cone is finitely generated, say by vectors $\binom{x_{1}}{\mu_{1}}, \ldots,\binom{x_{m}}{\mu_{m}}$, where we can suppose after rearranging and rescaling that there is $1 \leqslant k \leqslant m$ such that $\mu_{i}=1$ for all $i \leqslant k$ and $\mu_{i}=0$ for $i>k$.
Now by definition $P=\left\{\left.\binom{x}{\mu} \in K \right\rvert\, \mu=1\right\}$, so $x \in P$ if and only if $\binom{x}{1} \in K$, i.e., $\binom{x}{1}$ can be expanded as a positive combination of the generators of $K$ :

$$
\left\{\begin{array}{l}
x=\sum_{i \leqslant k} \lambda_{i} x_{i} \\
1=\lambda_{1}+\ldots+\lambda_{k}
\end{array}, \quad \lambda_{i} \geqslant 0 \text { for all } i\right.
$$

Recalling the definition of cone and convex hull, this is equivalent to

$$
x \in \operatorname{conv}\left\{x_{1}, \ldots, x_{k}\right\}+\operatorname{cone}\left\{x_{k+1}, \ldots, x_{m}\right\},
$$

and the claim follows with $Q=\operatorname{conv}\left\{x_{1}, \ldots, x_{k}\right\}$ and $C=\operatorname{cone}\left\{x_{k+1}, \ldots, x_{m}\right\}$, the latter being polyhedral by Proposition 5.2.8

Proof. Proof of Theorem 5.2.5 The claim follows from Proposition 5.2.9 by noting that the only bounded polyhedral cone is the cone generated by the single vector 0 , i.e., a single point.

### 5.3 Faces of polyhedra

Definition 5.3.1. A face of a polyhedron $P=\{A x \leqslant b\}$ is any subset of $P$ of the form

$$
F=\left\{x \in P \mid A^{\prime} x=b^{\prime}\right\}
$$

where $A^{\prime} x \leqslant b^{\prime}$ is a subsystem of (i.e., consists only of some of the inequalities of) $A x \leqslant b$.


Three faces of the polytope $P$ : a vertex $\left(F_{0}\right)$, a facet $\left(F_{1}\right)$ and the (improper) face $F_{2}=P$.

### 5.3.1 Supporting hyperplanes

An equivalent characterization of faces of polyhedra uses the notion of "supporting hyperplane". We call a hyperplane H a supporting hyperplane for a polyhedron P if " P touches H but lies wholly on one side of H ".

Definition 5.3.2. Let $\mathrm{P}=\{\mathrm{Ax} \leqslant \mathrm{b}\}$ be a polyhedron. An affine hyperplane $\mathrm{H}=\{\mathrm{c} \cdot \mathrm{x}=\delta\}$ is called a supporitng hyperplane for P if $\delta=\max \{\mathrm{c} \cdot x \mid \mathrm{x} \in \mathrm{P}\}$.

$\mathrm{H}_{1}$ and $\mathrm{H}_{2}$ are supporting hyperplanes of the polyhedron $P$.

Remark 5.3.3. Clearly every (nonempty) face F of a polyhedron P is of the form $F=P \cap H$ for some supporting hyperplane $H$. Indeed, suppose $F=\{x \in P \mid$ $\left.A^{\prime} x=b^{\prime}\right\}$ is nonempty. Then, if $c$ is the sum of all rows of $A^{\prime}$, clearly $F$ is the set of all points $x \in P$ attaining the maximum $\delta=\max \{c \cdot x \mid x \in P\}$.

Proving the converse to Remark 5.3 .3 - that every intersection of P with a supporting hyperplane is a face - needs more work, based on a corollary of Theorem 5.2.6 that bears considerable importance in its own right.

Lemma 5.3.4 ("Farkas' Lemma"). Let A be a matrix and b a vector. Then there exists $\mathrm{x} \geqslant 0$ such that $\mathrm{Ax}=\mathrm{b}$ if and only if $\mathrm{yb} \geqslant 0$ for each row vector y with $y A \geqslant 0$.

Proof. One direction is clear: if $x \geqslant 0, A x=b$ and $y \mathcal{A} \geqslant 0$ then clearly $y b=y A x \geqslant 0$.

For the other direction, let $a_{1}, \ldots, a_{m}$ be the columns of $A$ and argue by contraposition: if there is no $x \geqslant 0$ with $A x=b$, then $b \notin \operatorname{cone}\left\{a_{1}, \ldots, a_{m}\right\}$ and thus, by Theorem 5.2.6 there is $y$ with $y b<0$ and $y A \geqslant 0$ (the latter being equivalent to $y a_{i} \geqslant 0$ for all $i$ ).

Corollary 5.3.5. Let A be a matrix and b be a vector. Then the system $\mathrm{Ax} \leqslant \mathrm{b}$ has a solution x if and only if $\mathrm{yb} \geqslant 0$ for each row vector $\mathrm{y} \geqslant 0$ with $\mathrm{y} A=0$.

Proof. The system $A x \leqslant b$ has a solution if , and only if, the system

$$
[I|A|-A]\left[\begin{array}{l}
z_{0} \\
\left.\frac{z_{1}}{z_{2}}\right]
\end{array}\right]=\mathrm{b}
$$

has a nonnegative solution. (Indeed: given a solution $x$ to the first system, define vectors $x^{+}$and $x^{-}$as the positive and negative part ${ }^{1}$ of $x$, so that $x^{+} \geqslant$ $0, x^{-} \leqslant 0$ and $x=x^{+}+x^{-}$. then the (nonnegative) vector $\left[\frac{b-A x}{\frac{x^{+}}{-x^{-}}}\right]$solves the second. Conversely, if a nonnegative $\left[\begin{array}{l}z_{0} \\ \frac{z_{1}}{z_{2}}\end{array}\right]$ solves the second system, then $A\left(z_{1}-z_{2}\right)=b-z_{0} \leqslant 0$.) Now apply Farkas' Lemma to conclude that the latter system has a nonnegative solution if and only if $y b \geqslant 0$ for all $y$ with $y[I|A|-A]=[y I|y A|-y A] \geqslant 0$ - i.e, $y \geqslant 0$ (from the first block) and $y A=0$ (from the second and third block).

We are ready now to state and prove the following result, known as "duality theorem for linear programming".
Proposition 5.3.6. Let A be a matrix and $\mathrm{b}, \mathrm{c}$ be vectors. Then

$$
\max \{c \cdot x \mid A x \leqslant b\}=\min \{y b \mid y \geqslant 0, y A=c\}
$$

whenever both sets are nonempty.
Proof. If $A x \leqslant b, y \geqslant 0$ and $y A=c$, then $c \cdot x=y A x \leqslant y b$. This implies the inequality $\max \leqslant \min$, if the maxima and minima exist.

For the other inequality we need to show that there exist $x, y$ with $A x \leqslant b$, $y \geqslant 0, y A=c$ and $c \cdot x \geqslant y b$. This requirement we can translate in proving existence of a solution to

$$
\underbrace{\left[\begin{array}{cc}
0 & -1  \tag{5.6}\\
A & 0 \\
-c & b^{t} \\
0 & A^{t} \\
0 & -A^{t}
\end{array}\right]}_{=: \widetilde{A}}\left[\begin{array}{c}
x \\
y^{t}
\end{array}\right] \leqslant \underbrace{\left[\begin{array}{c}
0 \\
b \\
0 \\
c^{t} \\
-c^{t}
\end{array}\right]}_{\widetilde{b}}
$$

Now Corollary 5.3.5 implies that the solvability of (5.6) is equivalent to proving $w \widetilde{\mathrm{~b}} \geqslant 0$ for every $w \geqslant 0$ such that $w \widetilde{A}=0$. In order to check this last statement write $w=\left(w_{0}, w_{1}, w_{2}, w_{3}, w_{4}\right)$ according to the block-decomposition in (5.6). Then the statement's premise is that $w_{i} \geqslant 0$ for all $i$ and

$$
\begin{equation*}
\text { (1) } \quad w_{1} A-w_{2} \mathrm{c}=0, \quad \text { (2) } \quad w_{2} b^{t}+\left(w_{3}-w_{4}\right) A^{t} \geqslant 0 \tag{5.7}
\end{equation*}
$$

and we need to prove that

$$
\text { (3) } \quad w_{1} b+\left(w_{3}-w_{4}\right) c^{t} \geqslant 0
$$

Now we distinguish two cases:

$$
{ }^{1} \text { Explicitly: } x_{i}^{+}=\left\{\begin{array}{ll}
x_{i} & \text { if } x_{i}>0 \\
0 & \text { else. }
\end{array}, x_{i}^{-}= \begin{cases}x_{i} & \text { if } x_{i}<0 \\
0 & \text { else }\end{cases}\right.
$$

$w_{2}>0$. Here we write

$$
w_{2} w_{1} \mathrm{~b}=w_{2} \mathrm{~b}^{\mathrm{t}} w_{1}^{\mathrm{t}} \stackrel{(2)}{\geqslant}\left(w_{3}-w_{4}\right) \mathcal{A}^{\mathrm{t}} w_{1}^{\mathrm{t}} \stackrel{(1)}{=}\left(w_{3}-w_{4}\right) w_{2} \mathrm{c}=w_{2}\left(w_{3}-w_{4}\right) \mathrm{c},
$$

and (3) follows.
$w_{2}=0$ Since by assumption the sets over which maximum and minimum are taken in the proposition's claim are nonempty, we can choose elements $x_{0}, y_{0}$ with $A x_{0} \leqslant b$ and $y_{0} \geqslant 0, y_{0} A=c$. Then

$$
w_{1} b \geqslant w_{1} A x_{0} \stackrel{(1)}{=} 0 \stackrel{(2)}{\geqslant}\left(w_{4}-w_{3}\right) A^{t} y_{0}^{t}=\left(w_{4}-w_{3}\right) c^{t}
$$

and (3) follows, concluding the proof.

Corollary 5.3.7. Let P be a polyhedron. Then $\mathrm{F} \subsetneq \mathrm{P}$ is a face of P if and only if $\mathrm{F}=\mathrm{P} \cap \mathrm{H}$, where H is a supporting hyperplane of P .

Proof. Let $\mathrm{P}=\{\mathrm{Ax} \leqslant \mathrm{b}\}$ and let $\mathrm{F}=\mathrm{P} \cap \mathrm{H}$ with $\mathrm{H}=\{\mathrm{c} \cdot \mathrm{x}=\delta\}$ and $\delta:=\min \{\mathrm{c} \cdot \mathrm{x} \mid$ $x \in P\}$, for some $c \neq 0$. Now by Proposition 5.3.6 we have $\delta=\min \{y b \mid y \geqslant$ $0, y A=c\}$, and we can choose $y_{0}$ that attains this minimum. Now consider $A^{\prime} x \leqslant b^{\prime}$, the system obtained from the rows of $A x \leqslant b$ corresponding to positive entries of $y_{0}$. Then we have $F=\left\{x \in P \mid A^{\prime} x=b^{\prime}\right\}$ since for every $x \in P$ we have $A x \leqslant b$ and so $c \cdot x=\delta$ is equivalent to $y_{0} A x=y_{0} b$ and in turn (erasing trivial rows) to $A^{\prime} x=b^{\prime}$.

The other direction is proved in Remark 5.3.3

### 5.3.2 Facets

Definition 5.3.8. A facet of a polyhedron P is any maximal proper face of P .
Given a polyhedron $P=\{A x \leqslant b\}$ fix an enumeration $a_{1}, \ldots, a_{m}$ of the rows of $A$ (resp. the components of $b$ ), so that $P=\bigcap_{i=1, \ldots m}\left\{a_{i} \cdot x \leqslant b_{i}\right\}$. Such an inequality $a_{j} \cdot x \leqslant b_{j}$ is called an implicit equality if $P \subseteq\left\{a_{j} \cdot x=b_{j}\right\}$. We can then let $A^{=} x \leqslant b^{=}$denote the system of all implicit equalities and $A^{+} x \leqslant b^{+}$ the system consisting of all other inequalities, called effective inequalities.
Remark 5.3.9. Notice that any face $F=\left\{x \in P \mid A^{\prime} x=b^{\prime}\right\}$ can be written
 inequalities in $A^{\prime} x \leqslant \mathrm{~b}^{\prime}$. (in fact, any implicit equality is... implicit in the requirement " $x \in P$ ).
Remark 5.3.10. If a system has an effective inequality $a_{i} \cdot x \leqslant b_{i}$ then there must be a point $x \in P \backslash\left\{a_{i} \cdot x=b_{i}\right\}$ and in particular this point $x$ satisfies $A^{=} x=b^{=}$ and $A^{+} x<b^{+}$.

Recall that an inequality in a system of constraints is called redundant if its removal does not change the set of solutions.

Lemma 5.3.11. Suppose $\mathrm{a}_{\mathrm{i}} \cdot x \leqslant \mathrm{~b}_{\mathrm{i}}$ is an effective, non redundant inequality. Then $F:=\left\{x \in P \mid a_{i} \cdot x=b_{i}\right\}$ is a facet of $P$

Proof. Call $A^{\prime} x \leqslant b^{\prime}$ the system given by all effective inequalities other than $a_{i} \cdot x \leqslant b_{i}$. It is enough to prove that there is a point $x_{0}$ with

$$
A^{=} x_{0}=b^{=}, A^{\prime} x_{0}<b^{\prime} a_{i} \cdot x_{0}=b_{i}
$$

such an $x_{0}$ does not satisfy with equality any other row of $A^{\prime} x \leqslant b^{\prime}$, thus the only face bigger than $F$ is $P$ itself.

In order to find $x_{0}$ first recall Remark 5.3.10. which guarantees that we can find some $x_{1}$ with

$$
A^{=} x_{1}=b^{=} \text {and } A^{+} x_{1}<b^{+}
$$

and, since $a_{i} \cdot x \leqslant b_{i}$ is not redundant, we can also pick $x_{2}$ with

$$
A=x_{2}=b=\text { and } A^{\prime} x_{2}<b^{\prime}, a_{i} \cdot x_{2}>b_{i}
$$

Now we can pick $x_{0}:=t x_{1}+(1-t) x_{2}$ for some $t \in(0,1)$ such that $a_{i} \cdot x_{0}=b_{i}$.


Proposition 5.3.12. Every face of a polyhedron $P$ is an intersection of facets of $P$.
Proof. Let $F$ be a face of $P$. By definition, it can be written as $F=\left\{x \in P \mid A^{\prime} x=\right.$ $\left.b^{\prime}\right\}$ for some subsystem $A^{\prime} \leqslant b^{\prime}$ of $A x \leqslant b$, say consisting of the inequalities $a_{i}^{\prime} \cdot x \leqslant b_{i}^{\prime}$ for $i=1, \ldots, k$. By Remark 5.3 .9 we can assume that all $a_{i}^{\prime} \cdot x \leqslant b^{\prime}$ are effective, irredundant inequalities. Then, obviously

$$
\mathrm{F}=\bigcap_{i=1}^{\mathrm{k}}\left\{x \in \mathrm{P} \mid \mathrm{a}_{\mathrm{i}}^{\prime} \cdot x \leqslant \mathrm{~b}_{\mathrm{i}}^{\prime}\right\}
$$

and Lemma 5.3.11 says that every set on the right-hand side is a facet of $P$.

### 5.3.3 Vertices

Lemma 5.3.13. A polyhedron P is an affine subspace if and only if it has no nonempty faces except P itself.

Proof. If $P$ is an affine subspace then it has the form $P=\{M x=v\}$ for some matrix $M$ and some vector $v$. Its standard description via inequalities is $P=$ $\{A x \leqslant b\}$ with $A=\left[\begin{array}{c}M \\ -M\end{array}\right]$ and $b=\left[\begin{array}{c}v \\ -v\end{array}\right]$. Now obviously the set $P$ satisfies $A^{\prime} x=b^{\prime}$ for all subsystems of $A x \leqslant b$, so every nonempty face equals $P$.

On the other hand, if $P=\{A x \leqslant b\}$ has no nonempty faces except $P$, in particular for the "full" subsystem we have

$$
P=\{x \in P \mid A x=b\}=\{x \mid A x \leqslant b, A x=b\}=\{A x=b\}
$$

an affine subspace.
Corollary 5.3.14. The minimal nonempty faces of a polyhedron are affine subspaces.
Proof. From Definition 5.3.1 one gathers that every face F of P is a polyhedron in its own right, and every face $F^{\prime} \subseteq F$ is a face of $F$ if and only if it is a face of $P$. Thus the minimal nonempty faces of $P$ are polyhedra that do not have nontrivial faces themselves, and by Lemma 5.3.13 are affine subspaces.

Corollary 5.3.15. The minimal nonempty faces of a convex polytope P have dimension 0 (and are called "vertices" of P ). If $\mathrm{P}=\{\mathrm{A} x \leqslant \mathrm{~b}\}$, every vertex of P is the unique solution of some system $A^{\prime} x=b^{\prime}$, where the rows of $A^{\prime}$ are a maximal linearly independent set of the rows of $A$.

Proof. Since every convex polytope is bounded, the only affine subspaces it can contain are 0 -dimensional ones.

### 5.4 Unimodularity and integrality

Definition 5.4.1. Call a polytope $P \subseteq \mathbb{R}^{n}$ integral if all its vertices are points in $\mathbb{Z}^{n}$, rational if its vertices lie in $\mathbb{Q}^{n}$.

We are interested in finding conditions for a polytope to be integral - in particular, we are aiming at proving that Equation (5.2) defines the interior of an integral polytope. To this end, we introduce the notion of unimodularity of matrices.

Let $A$ be a matrix. A minor of $A$ is any matrix obtained by deleting some rows and some columns from $A$.

Definition 5.4.2. An integral matrix $A$ is called totally unimodular if every square minor of $A$ has determinant 0,1 or -1 .

Proposition 5.4.3 (Franklin and Veblen, 1921). Let A be a matrix with entries from $\{0,+1,-1\}$ and such that each column contains exactly one +1 and exactly one -1 . Then A is totally unimodular.

Proof. Let $M$ be a square minor of $A$. If $M$ has a column with at most one nonzero entry, developing its determinant with respect to that column proves the claim for $M$. Otherwise every column of $M$ has exactly one +1 and one -1 : in this case the sum of all rows vanishes, thus $M$ does not have full rank and the determinant of $M$ is 0 .

Corollary 5.4.4. Incidence matrices of graphs are totally unimodular.
Proof. See Exercise V. 2 .
Our goal in this section is to prove the following characterization of totally unimodular matrices.

Theorem 5.4.5. For an integral matrix $A$, the following statements are equivalent.
(i) The matrix A is totally unimodular
(ii) For all integral vectors $b$, all vertices of the polyhedron $\{x \mid A x \leqslant b, x \geqslant 0\}$ are integral.
(iii) For all integral vectors $a, b, c, d$, the polyhedron $\{x \mid a \leqslant A x \leqslant b, c \leqslant x \leqslant d\}$ has only integral vertices.

In order to prove this theorem we start by considering a weaker notion: unimodularity.

Definition 5.4.6. An integer matrix $A$ is called unimodular if every maximal square minor has determinant 0,1 or -1 .

Remark 5.4.7. Notice that an $n \times m$-matrix $A$ is totally unimodular if and only if the block matrix $[\mathrm{I} \mid \mathrm{A}]$ is unimodular, where I is the $\mathrm{n} \times \mathrm{n}$-identity matrix (see Exercise V.4)
Remark 5.4.8. Obviously a square matrix $B$ is unimodular if and only if $\operatorname{det} B=$ $\pm 1$. Thus, an invertible square matrix $B$ is unimodular if and only if $B^{-1}$ has only integer entries (as can be ascertained by recalling that $B^{-1}=\frac{1}{\operatorname{det} B} \operatorname{Adj}(B)$ and that the adjugate of an integral matrix is again integral).

Proposition 5.4.9. Let $A$ be an integral matrix of full row rank. Then $A$ is unimodular if and only if for every integral vector $b$ the polyhedron $P_{b}:=\{x \mid A x=b, x \geqslant 0\}$ has only integral vertices.

Proof.

A unimodular $\Rightarrow \mathrm{P}_{\mathrm{b}}$ integral for all integral b . Suppose that $A$ is a unimodular $\mathrm{n} \times \mathrm{m}$-matrix and consider a vertex $v$ of P . Then, there are $m$ linearly independent constraints in the system

$$
\left[\frac{-\mathrm{I}}{\frac{A}{-A}}\right] x \leqslant\left[\frac{0}{\frac{\mathrm{~b}}{-\mathrm{b}}}\right]
$$

that are satisfied with equality at $v$. In particular, the set of columns $\left\{a_{i}\right\}_{I}$ of $A$ corresponding to nonzero components of $v$ are linearly independent (since, given a linear dependency $\sum_{i} \lambda_{i} a_{i}=0$, we could add to $v$ a (small enough) multiple of the vector with $\lambda_{i}$ in the $i$-th component for all $i \in I$ and 0 otherwise, and we would obtain a positive-dimensional space $W$ satisfying with equality the same $m$ independent constraints as $v$, contradicting the fact that $v$ is a vertex). Now since $\mathcal{A}$ is full-rowrank we can extend the columns $\left\{a_{i}\right\}_{I}$ to a column-basis of $A$, obtaining a maximal (and invertible) square minor $B$ with the property that $B^{-1} b$ equals $v$ at each of its nonzero components. Since $A$ is unimodular, $B^{-1}$ and hence $B^{-1} b$ is integral (see Remark 5.4.8). Since nonzero coordinates of $v$ equal some coordinates of $\mathrm{B}^{-1} \mathrm{~b}$, it follows that $v$ is integral.
$\mathrm{P}_{\mathrm{b}}$ integral for all integral $\mathrm{b} \Rightarrow A$ unimodular. Suppose that $\mathrm{P}_{\mathrm{b}}$ is integral for all integral $b$ and let $B$ be an invertible maximal minor of $A$. Consider $B^{-1}$ and, in order to prove $\operatorname{det} B= \pm 1$, we show that $B^{-1} c$ is integral for every integral vector $c$ (see Remark 5.4.8. Take $c$ integral and choose $y \in$ $\mathbb{Z}^{n}$ such that $z:=y+B^{-1} c \geqslant 0$. Then $b:=B z \in \mathbb{Z}^{n}$ is integral and we can extend $z$ by zero components to a vector $z^{\prime} \in \mathbb{Z}^{m}$ with $A z^{\prime}=\mathrm{Bz}=\mathrm{b}$. Now $z^{\prime}$ is a vertex of $\mathrm{P}_{\mathrm{b}}$ (as it is in the polyhedron and satisfies n linear independent constraints with equality) and by assumption it is integral. It follows that $z$ (obtained from $z^{\prime}$ by canceling zero components) and with it $\mathrm{B}^{-1} \mathrm{c}=z-\mathrm{y}$ are integral, as required.

Proof. Proof of Theorem 5.4 .5
(i) $\Leftrightarrow$ (ii) Recall that $A$ is totally unimodular if and only if $[\mathrm{I} \mid A]$ is unimodular (Remark 5.4.7. By Proposition 5.4.9, this is equivalent to all vertices of the polyhedron $\{z \mid[\mathrm{I} \mid A] z=\mathrm{b}, z \geqslant 0\}$ being integral, for every integral b. The latter statement is equivalent to saying that all vertices of $\{x \mid x \geqslant$ $0, A x \leqslant b\}$ (i.e., the projection of the former polyhedron on the last $m$ coordinates) are integral.
(i) $\Leftrightarrow$ (iii) Notice that an integer matrix $A$ is unimodular if and only if the matrix

$$
\widetilde{A}:=\left[\begin{array}{r}
\frac{I}{-I} \\
\frac{A}{-A}
\end{array}\right]
$$

is totally unimodular as well. Now by the equivalence "(i) $\Leftrightarrow$ (ii)" this means that

$$
\widetilde{P}_{f}:=\{z \mid z \geqslant 0, \widetilde{A} x \leqslant f\}
$$

has integral vertices for all $f \in \mathbb{Z}^{3 n}$. Now given $a, b, c, d \in \mathbb{N}$ let

$$
\mathrm{f}_{0}:=\left[\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{c}}  \tag{5.8}\\
\frac{\mathrm{~b}}{\mathrm{a}}
\end{array}\right]
$$

and notice that $\widetilde{P}_{f_{0}}=\{x \mid a \leqslant A x \leqslant b, c \leqslant x \leqslant d\}$. Since every $f \in \mathbb{Z}^{3 n}$ can be written as in Equation 5.8 for suitable $a, b, c, d$, the proof of equivalence is complete.

### 5.5 Triangulations

We will need to "dissect" a polytope into smaller polytopes.
Definition 5.5.1. A polyhedral (resp. polytopal) complex is a collection $\mathcal{K}$ of polyhedra (resp. polytopes) such that every face of every member of $\mathcal{K}$ is again a member of $\mathcal{K}$, and for any two members $\mathrm{P}, \mathrm{Q} \in \mathcal{K}$ the intersection $\mathrm{P} \cap \mathrm{Q}$ is a face of both P and Q . We call faces (resp. vertices, etc.) of $\mathcal{K}$ the faces (vertices, etc.) of members of $\mathcal{K}$. The support of the complex is $|\mathcal{K}|:=\bigcup_{K \in \mathcal{K}} \mathrm{~K}$.

We call a set $X$ of points affinely independent if $\operatorname{dim} \operatorname{conv} X=|X|-1$.


Remark 5.5.2. Notice that given $x_{0}, \ldots, x_{d} \in \mathbb{R}^{d}$, then $\operatorname{conv}\left\{x_{0}, \ldots, x_{d}\right\}=x_{0}+$ $\left\{\sum_{i=1, \ldots, \mathrm{~d}} \lambda_{i}\left(x_{i}-x_{0}\right) \mid \lambda_{i} \geqslant 0, \sum_{i>0} \lambda_{i} \leqslant 1\right\}$ and this has dimension $d$ if and only if the vectors $x_{i}-x_{0}$ are linearly independent in $\mathbb{R}^{d}$.

In view of the next definition, recall also that a simplex is any polytope of the form conv $X$ for an affinely independent set of points $X$. In particular, a simplex has the smallest possible number of vertices for a polytope of its dimension.

Definition 5.5.3. Let P be a polytope. We call subdivision of P any polytopal complex $\mathcal{K}$ with $|\mathcal{K}|=P$. A triangulation of $P$ is a subdivision $\mathcal{T}$ of $P$ such that every member of $\mathcal{T}$ is a simplex.

Our goal is to show that every polytope has a triangulation all of whose vertices are vertices of the polytope. For this, we set up some notation.

Throughout this section let $V$ be a finite subset of $\mathbb{R}^{d}$, let $P=$ conv $V$ be a (full-dimensional) polytope in $\mathbb{R}^{\mathrm{d}}$, and let $\mathrm{h}: \mathrm{V} \rightarrow \mathbb{R}_{>0}$ any assignment of a positive real number to every vertex. We think of $h$ as a height function and will "lift" the points of $V$ accordingly into $\mathbb{R}^{d+1}$ obtaining a set of points

$$
\mathrm{V}^{(\mathrm{h})}:=\{(v, \mathrm{~h}(v)) \mid v \in \mathrm{~V}\} \subset \mathbb{R}^{\mathrm{d}+1}
$$

We will throughout write a point in $\mathbb{R}^{d+1}=\mathbb{R}^{d} \times \mathbb{R}$ as a pair $(x, t)$ with $x \in \mathbb{R}^{\mathrm{d}}, \mathrm{t} \in \mathbb{R}$. Then, we consider

$$
\pi: \mathbb{R}^{\mathrm{d}+1} \rightarrow \mathbb{R}^{\mathrm{d}}, \quad(\mathrm{x}, \mathrm{t}) \mapsto \mathrm{x}
$$

the standard coordinate projection (for instance, $\pi\left(\mathrm{V}^{(\mathrm{h})}\right)=\mathrm{V}$ ).
Moreover, consider the polyhedron

$$
\ell:=\left\{(0, t) \in \mathbb{R}^{d+1} \mid t \geqslant 0\right\}
$$

i.e., the "vertical" coordinate halfline in $\mathbb{R}^{d+1}$ (notice that $\pi(\ell)=0$ ).

Definition 5.5.4. The epigraph of $P$ with respect to the height function $h$ is

$$
\mathrm{P}^{\uparrow h}:=\operatorname{conv} \mathrm{V}^{(\mathrm{h})}+\ell .
$$



By Proposition 5.2.9, $\mathrm{P}^{\uparrow h}$ is a polyhedron. Let us study its supporting hyperplanes.

Lemma 5.5.5. With the notation above, given a vector $(\mathrm{a}, \alpha) \in \mathbb{R}^{\mathrm{d}+1}$ and $\mathrm{b} \in \mathbb{R}$ define a hyperplane $\mathrm{H}^{=}$and a halfspace $\mathrm{H} \leqslant$ as

$$
\mathrm{H}^{\square}:=\left\{(\mathrm{x}, \mathrm{t}) \in \mathbb{R}^{\mathrm{d}+1} \mid\langle\mathrm{a}, \mathrm{x}\rangle+\alpha \mathrm{t} \square \mathrm{~b}\right\} \quad \text { for } \square \in\{=, \leqslant\} \text {. }
$$

If $\mathrm{P}^{\uparrow h} \subseteq \mathrm{H} \leqslant$, then $\alpha \leqslant 0$. Moreover, if $\mathrm{H}^{=}$is a supporting hyperplane for $\mathrm{P}^{\uparrow h}$, then the face $\mathrm{H}^{=} \cap \mathrm{P}^{\uparrow h}$ is bounded if and only if $\alpha<0$.


Proof. If $\mathrm{P}^{\uparrow h} \subseteq \mathrm{H} \leqslant$, then in particular for every $\mathrm{p} \in \mathrm{P}^{\uparrow h}$ the half-line $\mathrm{p}+\ell$ lies in $H \leqslant$, meaning $\langle a, p\rangle+\alpha t \leqslant b$ for all $t \in \mathbb{R}_{+}$. This is only possible if $\alpha \leqslant 0$.

Suppose now that $\mathrm{H}^{=}$is a supporting hyperplane and consider the face $F:=H^{=} \cap \mathrm{P}^{\uparrow h}$. Now $F$ is unbounded if and only if $F+\ell \subseteq F$. In particular, for every $p \in F$ we have $p+\left(0, \mathbb{R}_{\geqslant 0}\right) \subseteq H^{=}$, i.e., $\langle a, p\rangle+\alpha t=b$ for all $t \in \mathbb{R}_{\geqslant 0}$. This is only possible if $\alpha=0$.

Corollary 5.5.6. Every bounded face is a face of some bounded facet.
Proof. Let F be a bounded face. By Corollary 5.3.15 we know that F is an intersection of facets $F_{1}, \ldots, F_{k}$ : in particular, $F$ is a face of each $F_{i}$. Now assume that every $F_{i}$ is unbounded. By Lemma 5.5.5, this means that $F_{i}$ is contained in a vertical hyperplane, and with any $x \in F_{i}$ we have $x+\ell \subseteq F_{i}$. In particular, also $F$, the intersection of the $F_{i}$, contains the whole halfline $x+\ell$ for every $x \in F$. Since $F$ is non-empty, $F$ would be unbounded: this is a contradiction to all $F_{i}$ being unbounded.

Lemma 5.5.7. It is possible to choose the height function $h$ with the following property: if $\left\{v_{0}, \ldots, v_{\mathrm{d}}\right\} \subseteq \mathrm{V}$ is a maximal affinely independent set, then the $\left(v_{i}, \mathrm{~h}\left(v_{i}\right)\right)$ are affinely independent and span a hyperplane that does not contain any other point of $\mathrm{V}^{(\mathrm{h})}$.

Proof. Write $V=\left\{v_{1}, \ldots, v_{\mathrm{m}}\right\}$ and consider any $\mathrm{d}+1$-subset $\mathrm{I}=\left\{\mathfrak{i}_{0}, \ldots, \mathfrak{i}_{\mathrm{d}}\right\} \in$ $\binom{[m]}{d+1}$ such that $\left\{v_{i}\right\}_{i \in I}$ is affinely independent. In order for $\left\{\left(v_{i}, h_{i}\right)\right\}_{i \in I}$ to
determine a hyperplane in $\mathbb{R}^{\mathrm{d}+1}$ they must be affinely independent, and thus (by Remark 5.5.2) we must have that the determinant

$$
\delta(\mathrm{I}):=\left|\begin{array}{cccc}
v_{\mathfrak{i}_{1}}-v_{i_{0}} & \cdots & v_{i_{d}}-v_{i_{0}} & 0 \\
h_{i_{1}}-h_{i_{0}} & \cdots & h_{i_{d}}-h_{i_{0}} & 1
\end{array}\right|
$$

is non-zero (notice that the last column is never a nontrivial linear combination of the first $d$, since $v_{i_{1}}-v_{i_{0}}, \ldots, v_{i_{d}}-v_{i_{0}}$ are linearly independent by assumption). This is however automatically ensured since this determinant equals $\pm\left|v_{i_{1}}-v_{i_{0}} \cdots v_{i_{d}}-v_{i_{0}}\right| \neq 0$ because of affine independence of the $v_{i}$.

Now, if some $\left(v_{j}, h_{j}\right), j \notin I$, lies on the same hyperplane as $\left(v_{i_{0}}, h_{i_{0}}\right), \ldots,\left(v_{i_{d}}, h_{i_{d}}\right)$, then

$$
\gamma(\mathrm{I} \cup\{j\}):=\left|\begin{array}{cccc}
v_{i_{0}} & \cdots & v_{\mathfrak{i}_{\mathrm{d}}} & v_{\mathfrak{j}} \\
\mathrm{h}_{\mathrm{i}_{0}} & \cdots & h_{\mathfrak{i}_{\mathrm{i}}} & h_{\mathrm{j}} \\
1 & \cdots & 1 & 1
\end{array}\right|=0
$$

Thus, in order to satisfy the claim, we have to choose $h \in \mathbb{R}^{m}$ so that

$$
\prod_{\mathrm{I}, \mathrm{j} \notin \mathrm{I}} \gamma(\mathrm{I} \cup\{\mathrm{j}\}) \neq 0
$$

where the product is over all index sets I of maximal affinely independent points in V.

As the expression on the left-hand side is a non-constant polynomial, the set of "good" $h$ is nonempty.

Theorem 5.5.8. Every polyhedron $\mathrm{P}=\operatorname{conv}(\mathrm{V})$ has a triangulation all of whose vertices are in V .

Proof. Choose $h$ as in Lemma 5.5.7 and consider the associated $\mathrm{P}^{\uparrow h}$. Define

$$
\mathcal{T}:=\left\{\pi(F) \mid F \text { is a bounded face of } \mathrm{P}^{\uparrow h}\right\}
$$

Clearly the 0 -dimensional members of $\mathcal{T}$ are a subset of $V$. Moreover, for every bounded facet $F$ of $P^{\uparrow h}$, say supported on a hyperplane $H$ (say $F=H^{=} \cap P^{\uparrow h}$ and $\left.P^{\uparrow h} \subseteq H^{\leqslant}=\{\langle(a, \alpha),(x, t)\rangle \leqslant b\}\right)$ :
(a) $F$ is a simplex.

Proof. By the choice of $h$ we know that every bounded facet $F$ of $P^{\uparrow h}$ is supported on a hyperplane that contains at least $d+1$ affinely independent vertices (since $\operatorname{dim}(F)=d$ ), but no hyperplane can contain more than $d+1$ points from $V^{(h)}$ : thus $F$ is a simplex.
(b) The restriction $\pi: F \rightarrow \pi(F)$ is bijective, with inverse $\iota: x \mapsto\left(x, \frac{b-\langle a, x\rangle}{\alpha}\right)$.

Proof. Since F is bounded, by Lemma $5.5 .5 \alpha \neq 0$. Thus the inverse is well-defined and proves bijectivity.
(c) For every $X \subseteq F, \operatorname{conv}(\pi(X))=\pi(\operatorname{conv}(X))$

Proof. Given $\left(q_{0}, t_{0}\right), \ldots,\left(q_{k}, t_{k}\right) \in F$ and $\lambda_{1}, \ldots, \lambda_{k} \in \mathbb{R}$, we have $\left(q_{0}, t_{0}\right)=$ $\sum_{i} \lambda_{i}\left(q_{i}, t_{i}\right)$ if and only if $q_{0}=\sum_{i} \lambda_{i} q_{i}$ (for the "if" direction: recall the inverse map of item (b) and notice that $t_{i}=\left(b-\left\langle a, q_{i}\right\rangle\right) / \alpha$ for all $i$, and thus $\left.\alpha \sum_{i \geqslant 1} \lambda_{i} t_{i}=b-\left\langle a, \sum_{i \geqslant 1} \lambda_{i} q_{i}\right\rangle=b-\left\langle a, q_{0}\right\rangle=\alpha t_{0}\right)$
(d) $\pi(\mathrm{F})$ is a simplex with set of vertices $\pi\left(\mathrm{V}_{\mathrm{F}}\right)$, where $\mathrm{V}_{\mathrm{F}}$ is the set of vertices of $F$.

Proof. Since $F$ is a simplex of dimension $d,\left|V_{F}\right|=d+1$. Then, by item (c) $\pi(\mathrm{F})=\operatorname{conv} \pi\left(\mathrm{V}_{\mathrm{F}}\right)$ has at most $\mathrm{d}+1$ vertices. Moreover, linearity of $\pi$ implies that $\operatorname{dim} \pi(\mathrm{F})=\operatorname{dim}(\mathrm{F})=\mathrm{d}$ and, in particular, $\pi(\mathrm{F})$ has at least $d+1$ vertices. The claim follows.
(e) The faces of $\pi(F)$ correspond bijectively to projections of faces of $F$.

Proof. Faces of simplices correspond bijectively to (convex hulls of) subsets of vertices. Per items (a) and (d) both F and $\pi(\mathrm{F})$ are simplices with vertices, respectively, $V_{F}$ and $\pi\left(V_{F}\right)$. Thus the bijection $2^{V_{F}} \rightarrow 2^{\pi\left(V_{F}\right)}$, $\mathrm{I} \mapsto \pi(\mathrm{I})$, induces a bijection between the faces of F and $\pi(\mathrm{F})$.

## Moreover,

(f) $\mathcal{T}$ consists of all faces of projections of bounded facets.

Proof. By Corollary 5.5.6 every bounded face of $\mathrm{P}^{\uparrow h}$ is a face of a bounded facet. Thus, by item (e), $\mathcal{T}$ is the set of all faces of some $\pi(F), F$ a bounded facet of $F$.

We now claim that $\mathcal{T}$ is a triangulation of $P$. For this we need to check three things.
(1) $|\mathcal{T}|=P$.

The inclusion $\subseteq$ follows trivially from $\mathrm{P}=\pi\left(\mathrm{P}^{\uparrow h}\right)$.
In order to prove the inclusion $\supseteq$ is is enough to prove that the open interior of P is contained in $\pi\left(\mathrm{P}^{\uparrow h}\right)$, as $\pi$ is continuous. Let then $p$ be a point in the interior of $P$ and let $\epsilon$ be such that the $\epsilon$-ball $B_{\epsilon}(p)$ is fully contained in P. Moreover, let

$$
\ell_{p}:=(p+\ell) \cap p^{\uparrow h} \text {,i.e., } \ell_{p}=x+\ell
$$

for some $x \in \mathrm{P}^{\uparrow h}$. Then $x$ is on the boundary of $\mathrm{P}^{\uparrow h}$ and therefore it is in the intersection of $\mathrm{P}^{\uparrow h}$ with a supporting hyperplane. In the notation of Lemma 5.5.5, there is a vector $(a, \alpha)$ (which we can suppose to be of norm 1) and a number $b$ such that $\mathrm{H}^{=}$is a supporting hyperplane of $\mathrm{P}^{\uparrow h}$ and

$$
x \in \mathrm{~F}:=\mathrm{P}^{\uparrow h} \cap \mathrm{H} \leqslant .
$$

Now, if $\alpha=0$ then $\pi\left(\mathrm{H}^{=}\right)$is a hyperplane in $\mathbb{R}^{\mathrm{d}}$ that contains p and bounds the halfspace $\pi\left(\mathrm{H}^{\leqslant}\right)$. Now $\mathrm{P}=\pi\left(\mathrm{P}^{\uparrow h}\right) \subseteq \pi\left(\mathrm{H}^{\leqslant}\right)$, but $\mathrm{p}+\frac{\epsilon}{2} \mathrm{a} \notin$ $\pi\left(\mathrm{H}^{\leqslant}\right)$, a contradiction to $\mathrm{B}_{\epsilon}(\mathrm{p}) \subseteq \mathrm{P}$.

This means that $\alpha>0$ and, by Lemma 5.5.5, the face $F$ is bounded. Therefore $p=\pi(x) \in \pi(F) \subseteq|T|$ as was to be shown.
(2) Every member of $\mathcal{T}$ is a simplex.

Item (f) shows that every member of $\mathcal{T}$ is a face of $\pi(F)$ for some bounded facet $F$. Since every face of a simplex is a simplex, item (d) ensures that every member of $\mathcal{T}$ is a simplex.
(3) $\mathcal{T}$ is a polytopal complex.

Since every face of a bounded face is bounded, item (f) shows that $\mathcal{T}$ contains all faces of each of its members. We are left to show that the intersection of any two $T_{1}, T_{2} \in \mathcal{T}$ is a face of both. Let $F_{1}$ and $F_{2}$ be bounded faces of $P^{\uparrow h}$ such that $T_{i}=\pi\left(F_{i}\right)$ for $i=1$,2. By item (e) it is enough to prove

$$
\pi\left(\mathrm{F}_{1} \cap \mathrm{~F}_{2}\right)=\mathrm{T}_{1} \cap \mathrm{~T}_{2}
$$

The inclusion $\subseteq$ is trivial.
Furthermore, for any $p \in T_{1} \cap T_{2} \backslash \pi\left(F_{1} \cap F_{2}\right)$ the halfline $p+\ell \supseteq$ intersects both $\mathrm{F}_{i} \mathrm{~s}$. Therefore there are $t_{1}, t_{2}$ such that $\left(p, t_{1}\right) \in$
 $F_{1}$ and $\left(p, t_{2}\right) \in F_{2}$, and $t_{1} \neq t_{2}$ because $p$ is not in $\pi\left(\mathrm{F}_{1} \cap \mathrm{~F}_{2}\right)$. Without loss of generality we can assume $t_{1}<t_{2}$. Consider a supporting hyperplane $H=\{(a, \alpha) \cdot x=\delta\}$ to $F_{2}$. Then $\left(p, t_{2}\right) \in F_{2}$ implies $\langle a, p\rangle+\alpha t_{2}=\delta$. Since $F_{2}$ is bounded, we have $\alpha<0$ and so $\langle a, p\rangle+\alpha t_{1}>\delta$, contradicting the fact that, since $H$ is a supporting hyperplane, we should have $\mathrm{P} \subseteq \mathrm{H} \leqslant$. Thus $\mathrm{T}_{1} \cap \mathrm{~T}_{2} \subseteq$ $\pi\left(F_{1} \cap F_{2}\right)$ and we are done.

Definition 5.5.9. We say that a polyhedron P (resp. a pointed cone) can be triangulated using no new vertices if there exists a triangulation $\mathcal{T}$ of $P$ such that the vertices of $\mathcal{T}$ are vertices of $P$.

Analogously, we say that a pointed cone C can be triangulated using no new generators if there exists a triangulation $\mathcal{T}$ of $\mathbb{C}$ such that the generators of $\mathfrak{T}$ are generators of $C$.

Corollary 5.5.10. Every polyhedron P can be triangulated using no new vertices.

### 5.6 Sources and Bibliography

The approach of Section 5.1] is from the original paper by Kochol [5]. Sections 5.2 to 5.4 follow more or less closely the treatment by Schrijver [6]. Section 5.5 is based on Chapter 5 of Beck and Sanyal [2] and Chapter 3 of Beck and Robins [1]

## Exercises V

V. 1 This exercise should give a geometric intuition for the "Farkas' type" statements we have seen. Let
(a) Consider the statement of Farkas' lemma. Let

$$
A:=\left[\begin{array}{cc}
1 & 2 \\
-1 & -1
\end{array}\right], \quad b_{1}:=\left[\begin{array}{c}
8 \\
-2
\end{array}\right], \quad b_{2}:=\left[\begin{array}{c}
8 \\
-5
\end{array}\right] .
$$

For $i=1,2$ do the following: determine the solutions of $A x=b_{i}$; then draw the set of points $y$ with $y \mathcal{A} \geqslant 0$ and check whether $y b_{i} \geqslant 0$ for all such $y$.
(b) In order to illustrate Corollary 5.3.5, consider

$$
A:=\left[\begin{array}{cc}
1 & 1 \\
1 & -2 \\
-2 & 1
\end{array}\right], \quad b_{1}:=\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right], \quad b_{2}:=\left[\begin{array}{c}
-1 \\
0 \\
0
\end{array}\right] .
$$

For $i=1,2$ determine the set of solutions of the system $A x \leqslant b_{i}$, consider the set of all $y \in \mathbb{R}^{3}$ with $y \geqslant 0$ and $y A=0$ and verify whether $y b \geqslant 0$ for all such $y$.
(c) Check the claim of Proposition 5.3.6 using the data $A$ from (b) above and $\mathrm{b}:=\left[\begin{array}{l}1 \\ 2 \\ 0\end{array}\right], \mathrm{c}:=\left[\begin{array}{l}1 \\ 1\end{array}\right]$. In particular, draw a picture of $A x \leqslant b$ and of the points where $c \cdot x$ is maximum; then draw the set $y A=c$, $y \geqslant 0$ and compute the minimum on the right-hand side.
V. 2 Recall the definition of the incidence matrix of a graph and prove Corollary 5.4.4.
V. $3(+)$ Let P be a polyhedron and consider the set $\mathscr{F}(\mathrm{P})$ of all faces of P , partially ordered by inclusion.
(a) Let $x_{0}, \ldots, x_{d} \in \mathbb{R}^{d}$ be chosen so that these points do not lie all on the same (affine) hyperplane. Then $P:=\operatorname{conv}\left\{x_{0} l d o t s, x_{d}\right\}$ is called a simplex. What is the poset of faces of a simplex? (You can start by concrete examples with $d=2,3$ ).
(b) Prove that for every polyhedron P , the poset $\mathscr{F}(\mathrm{P})$ is ranked.
(c) ++ Prove that for every convex polytope $P$ the poset $\mathscr{F}(P)$ is a ranked lattice.
V. 4 Prove that an integer $n \times m$ matrix $A$ is totally unimodular if and only if the block matrix $[1 \mid A$ ] is unimodular (here 1 denotes the $n \times n$ identity matrix).
V. 5 ([1, Exercise 3.3]) Let $e_{1}, \ldots, e_{d}$ be the standard basis vectors of $\mathbb{R}^{d}$. For every permutation $\pi$ of the set $[n]$ consider the polytope

$$
\Delta_{\pi}:=\operatorname{conv}\left\{0, e_{\pi(1)}, e_{\pi(1)}+e_{\pi(2)}, \ldots, e_{\pi(1)}+e_{\pi(2)}+\ldots+e_{\pi(\mathrm{d})}\right\}
$$

Prove that the polytopes $\Delta_{\pi}$ and all their faces define a triangulation of the unit cube $[0,1]^{\mathrm{d}}$. Moreover, prove that any two such simplices can be obtained from one another by a congruence (a concatenation of reflections, rotations and translations).

## Chapter 6

## Rational generating functions

The last ingredient we need towards Ehrhart's theorem is rational generating functions as "representations" of number sequences.

### 6.1 Basics

Fix a natural number $d$ and a field $\mathbb{K}$.
The main idea is that any function $f: \mathbb{Z}^{d} \rightarrow \mathbb{K}$ can be represented by the formal expression in the variables $z_{1}, \ldots, z_{\mathrm{d}}$ :

$$
\sum_{m \in \mathbb{Z}^{\mathrm{d}}} f(m) z^{m}, \quad \text { where we write } z^{m}:=z_{1}^{m_{1}} \cdots z_{d}^{m_{\mathrm{d}}}
$$

Every such expression is called formal Laurent series. The word "formal" indicates that we are not concerned about analytic convergence. We treat such "sum" as abstract formal objects.

It is common to call

$$
\mathbb{C}\left[\left[z_{1}, \ldots, z_{d}\right]\right]:=\left\{\sum_{m \in \mathbb{N}} f(m) z^{m} \mid f: \mathbb{N}^{d} \rightarrow \mathbb{C}\right\}
$$

the set of formal power series in the variables $z_{1}, \ldots, z_{d}$. The degree of $\sum_{m \in \mathbb{Z}^{d}} f(m) z^{m}$ is defined as $\max \left\{\sum_{i} m_{i} \mid f(m) \neq 0\right\} \in \mathbb{N} \cup\{\infty\}$. Thus an element of $\mathbb{C}\left[\left[z_{1}, \ldots, z_{d}\right]\right]$ of finite degree is a polynomial, i.e., an element of $\mathbb{C}\left[z_{1}, \ldots, z_{d}\right]$.

### 6.1.1 Algebraic structure

The following addition and multiplication rules define a ring structure on the set of formal power series (and of formal Laurent series, by changing the index
set from $\mathbb{N}$ to $\mathbb{Z}$ ):

$$
\begin{gathered}
\sum_{\mathfrak{m} \in \mathbb{N}} f(m) z^{m}+\sum_{m \in \mathbb{N}} g(m) z^{m}:=\sum_{m \in \mathbb{N}}(f(m)+g(m)) z^{m} \\
\sum_{m \in \mathbb{N}} f(m) z^{\mathfrak{m}} \cdot \sum_{m \in \mathbb{N}} g(m) z^{\mathfrak{m}}:=\sum_{m \in \mathbb{N}}\left(\sum_{k \in \mathbb{Z}} f(k) g(m-k)\right) z^{m} .
\end{gathered}
$$

Example 6.1.1. The following identity holds in the ring $\mathbb{C}[[z]]$ :

$$
\frac{1}{1-z}=\sum_{k \in \mathbb{N}} z^{k}
$$

Thus, the function $f: \mathbb{N} \rightarrow \mathbb{C}$ constant equal to 1 can be represented by the fraction on the left-hand side. We call "rational" any formal power series that equals a rational expression in the ring of formal power series.

The gist now is: whether a function $\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ agrees with a polynomial can be ascertained by looking at the rationality of the associated formal power series. In the next section we will make this precise.

Before that, let us collect some considerations on the case $\mathbb{K}=\mathbb{C}$

### 6.2 The main theorem

Theorem 6.2.1. Fix $\alpha_{1}, \ldots, \alpha_{d} \in \mathbb{C}$ with $\mathrm{d} \geqslant 1$ and $\alpha_{\mathrm{d}} \neq 0$. Then, for every function $\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ the following are equivalent:
(i)

$$
\sum_{n \in \mathbb{N}} f(n) z^{n}=\frac{p(z)}{q(z)}
$$

where $\mathrm{q}(z)=1+\alpha_{1} z+\ldots+\alpha_{\mathrm{d}} z^{\mathrm{d}}$ and $\mathrm{p}(z)$ is any polynomial in $z$ of degree less than d.
(ii) For all $\mathrm{n} \in \mathbb{N}$ :

$$
\sum_{i=0}^{d} \alpha_{i} f(n+d-i)=0
$$

(iii) For all $\mathrm{n} \in \mathbb{N}$,

$$
f(n)=\sum_{i=1}^{k} p_{i}(n) \gamma_{i}^{n}
$$

where the $\gamma_{i}$ are distinct, nonzero complex numbers with $1+\alpha_{1} z+\ldots+$ $\alpha_{\mathrm{d}} z^{\mathrm{d}}=\left(1-\gamma_{1} z\right)^{\mathrm{d}_{1}} \cdots\left(1-\gamma_{\mathrm{k}} z\right)^{\mathrm{d}_{\mathrm{k}}}$ and each $\mathrm{p}_{\mathrm{i}}$ is a polynomial of degree less than $\mathrm{d}_{\mathrm{i}}$

Proof. Fix $q(z):=1+\alpha_{i} z+\ldots+\alpha_{d} z^{\mathrm{d}}=\left(1-\gamma_{1} z\right)^{\mathrm{d}_{1}} \cdots\left(1-\gamma_{\mathrm{k}} z\right)^{\mathrm{d}_{\mathrm{k}}}$, all $\gamma_{\mathrm{i}}$ nonzero and distinct. Consider the following complex vectorspaces:

- $V_{1}:=\{\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ such that (i) holds $\}$.

Notice that $\operatorname{dim} V_{1}=d$, as in (i) we can choose the $d$ coefficients of $p(z)$ freely.

- $\mathrm{V}_{2}:=\{\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ such that (ii) holds $\}$.

Here we have $\operatorname{dim} V_{2}=d$ because in (ii) the choice of the $f(n+i), i=$ $0, \ldots, d-1$ is free and determines $f$ completely.

- $\mathrm{V}_{3}:=\{\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ such that (iii) holds $\}$.

Here, for every $i$ we can choose the $d_{i}$ coefficients of $p_{i}$ freely, and this choice determines $f$ completely. Thus, $\operatorname{dim} V_{3}=d_{1}+\ldots+d_{k}=d$.

- $V_{4}:=\left\{\begin{array}{l}f: \mathbb{N} \rightarrow \mathbb{C} \text { such that } \\ \sum_{n \in \mathbb{N}} f(n) z^{n}=\sum_{i=1}^{k} \sum_{j=1}^{d_{i}} \beta_{i j}\left(1-\gamma_{i} z\right)^{-j} \quad(\sharp) \\ \text { for some } \beta_{i j} \in \mathbb{C} \text { and where the } \gamma_{i}, d_{i} \text { are as in (iii) }\end{array}\right\}$.

We claim that $\operatorname{dim} V_{4}=d$. This is not as evident as in the other cases.
First, notice that $V_{4}$ is spanned over $\mathbb{C}$ by the rational functions

$$
\mathrm{R}_{\mathfrak{i j}}(z):=\frac{1}{\left(1-\gamma_{\mathfrak{i}} z\right)^{\mathfrak{j}}}, \quad \mathfrak{i}=1, \ldots, k, \mathfrak{j}=1, \ldots, \mathrm{~d}_{\mathfrak{i}} .
$$

Since there are $\sum_{i} d_{i}=d$ such $R_{i j} s$, we have $\operatorname{dim} V_{4} \leqslant d$.
Now it is enough to show that the $R_{i j}$ are linearly independent. By way of contradiction suppose that there is a linear dependency $\sum c_{i j} R_{i j}(z)=$ 0 with complex coefficients $c_{i j}$ not all equal to zero. Let $i_{0}$ such that $c_{i_{0} j} \neq 0$ for some $j$, and let $j_{0}$ be the largest such index. Now multiplying the linear dependency by $\left(1-\gamma_{i_{0}} z\right)^{j_{0}}$ and setting $z=\frac{1}{\gamma_{i_{0}}}\left\{^{1}\right.$ ) gives the equality $\mathrm{c}_{\mathfrak{i}_{0} \mathrm{j}_{0}}=0$, a contradiction. Thus the $\mathrm{R}_{\mathfrak{i j}}$ are linearly independent and $\operatorname{dim} V_{4}=\mathrm{d}$.

Now it is enough to show that $V_{1}=V_{2}=V_{3}=V_{4}$. We now do so, in three steps.
$V_{3}=V_{4}$ Since $\operatorname{dim} V_{3}=\operatorname{dim} V_{4}$ it is enough to prove $V_{4} \subseteq V_{3}$. To this end, recall that

$$
\frac{1}{\left(1-\gamma_{i} z\right)^{j}}=\left(\sum_{n \in \mathbb{N}} \gamma_{i}^{n} z^{n}\right)^{j}=\sum_{n \in \mathbb{N}}\binom{j+n-1}{j-1} \gamma_{i}^{n} z^{n} .
$$

[^3]Notice that $\binom{j+n-1}{j-1}$, as a function of $n$, is a polynomial of degree $j-1$. Now let $f \in V_{4}$. For all $n, f(n)$ equals the coefficient of $z^{n}$ in $(\sharp)$, which we can compute to be

$$
\begin{equation*}
\sum_{i=1}^{k}\left(\sum_{j=1}^{d_{i}} \beta_{i j}\binom{j+n-1}{j-1}\right) \gamma_{i}^{n} \tag{b}
\end{equation*}
$$

where the expression inside the big parenthesis is a polynomial in $n$ of degree at most $d_{i}-1$. Therefore $f \in V_{3}$ and the proof is complete.
$V_{4}=V_{1}$ The right-hand side of ( $\#$ ) equals

$$
\begin{equation*}
\frac{\sum_{i=1}^{k}\left(\prod_{h \neq i}\left(1-\gamma_{h} z\right)^{d_{h}} \sum_{j=1}^{d_{i}} \beta_{i j}\left(1-\gamma_{i} z\right)^{d_{i}-j}\right)}{\prod_{i=1}^{k}\left(1-\gamma_{i} z\right)^{d_{i}}} \tag{bb}
\end{equation*}
$$

where the numerator is a polynomial of degree strictly smaller than $\sum_{i} d_{i}=d$. Thus, $V_{4} \subseteq V_{1}$. Since both spaces have dimension $d$, equality follows.
$V_{1}=V_{2}$ If $f \in V_{1}$, then

$$
q(z) \sum_{n \in \mathbb{N}} f(n) z^{n}=p(z)
$$

Now equating the coefficient of $z^{\mathrm{d}}$ on both sides of the last equality gives the relation in (ii).

Call a function $\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ polynomial if there is a polynomial $p \in \mathbb{C}[z]$ with $f(n)=p(n)$ for all $n \in \mathbb{N}$.

Corollary 6.2.2. Let $\mathrm{f}: \mathbb{N} \rightarrow \mathbb{C}$ and $\mathrm{d} \in \mathbb{N}$. Then f is polynomial of degree at most d if and only if

$$
\begin{equation*}
\sum_{n \in \mathbb{N}} f(n) z^{n}=\frac{p(z)}{(1-z)^{d+1}} \tag{দ̆}
\end{equation*}
$$

for some $\mathfrak{p}(\mathrm{x}) \in \mathbb{C}[\mathrm{x}]$ of degree at most d .
Moreover, f is polynomial of degree exactly d if and only if $\mathfrak{p}(1) \neq 0$.
Proof. Equation ( $(\mathrm{f})$ is an instance of Theorem 6.2.1(i), with $\mathrm{k}=1$ and $\gamma_{1}=1$. From the equivalence of (i) and (iii) in Theorem 6.2.1] we have then that ( $t$ ) is equivalent to

$$
\mathrm{f}(\mathrm{n})=\mathrm{p}_{1}(\mathrm{n})
$$

for a polynomial $p_{1} \in \mathbb{C}[x]$ of degree at most $d$.

For the second claim notice that following (b) in the proof of Theorem 6.2.1 our fatisfies

$$
f(n)=\sum_{j=1}^{d+1} \beta_{1 j}\binom{j+n-1}{j-1}
$$

for some numbers $\beta_{1 j}$, so $f$ is polynomial of degree $d$ if and only if $\beta_{1(d+1)} \neq$ 0 . Now, the numbers $\beta_{i j}$ are related to $p(z)$ via the expression in (bb). In our case, this means

$$
p(z)=\sum_{j=1}^{d+1} \beta_{i j}(1-z)^{d+1-j}
$$

and from here $p(1) \neq 1$ if and only if $\beta_{1(d+1)} \neq 0$. The claim follows.

## Exercises VI

VI. 1 Give a nonrecursive form for the general term of the following recursively defined integer sequences.

- $a_{0}=2, a_{1}=3, a_{n}=3 a_{n-1}-2 a_{n-2}$ for $n \geqslant 2$
- $a_{0}=0, a_{1}=2, a_{n}=4 a_{n-1}-4 a_{n-2}$ for $n \geqslant 2$
- $a_{0}=5, a_{1}=12, a_{n}=4 a_{n-1}-3 a_{n-2}-2^{n-2}$ for $n \geqslant 2$


## Chapter 7

## Ehrhart theory

> NOTE: This chapter follows very closely chapter 3 of Beck-Robins "Computing the continuum discretely". I omit proofs and examples that can be found in that chapter.

Our goal in this chapter is to determine the lattice-point enumerator function of a given convex polytope $P$ in $\mathbb{R}^{d}$, defined for every integer $t$ as

$$
\mathrm{L}_{\mathrm{P}}(\mathrm{t}):=\#\left(\mathrm{tP} \cap \mathbb{Z}^{\mathrm{d}}\right)=\#\left(\mathrm{P} \cap \frac{1}{\mathrm{t}} \mathbb{Z}^{\mathrm{d}}\right)
$$

i.e., as the number of integer points contained in the $t$-th dilate of $P$.

This chapter is devoted to proving
Theorem 7.0.1. For every integral d-dimensional polytope $P, L_{P}(t)$ is polynomial of degree d.

The preceding chapter should suggest to us to consider the Ehrhart series in one variable $z$ :

$$
\operatorname{Ehr}_{P}(z):=\sum_{t \geqslant 0} L_{P}(t) z^{t}
$$

### 7.1 Formal Laurent series from integer points

Definition 7.1.1. Let $P$ be any polyhedron in $\mathbb{R}^{d}$. Define

$$
\sigma_{\mathrm{P}}\left(z_{1}, \ldots, z_{\mathrm{d}}\right):=\sum_{\mathrm{m} \in \mathrm{P} \cap \mathbb{Z}^{\mathrm{d}}} z^{\mathrm{m}}
$$

We call this the integer-point transform of P .
Example 7.1.2. See Examples 3.3 and 3.4 in the book by Beck and Robins [1]

Recall that a simplicial cone is a polyhedral cone $C$ generated by $\operatorname{dim}(C)$ vectors.

Theorem 7.1.3. Let $m_{1}, \ldots, m_{d} \in \mathbb{Z}^{\mathrm{d}}$ be such that

$$
C:=\operatorname{cone}\left\{m_{1}, \ldots, m_{d}\right\}
$$

is a simplicial d-dimensional cone. Let

$$
\Pi:=\left\{\sum_{i=1}^{d} \lambda_{i} m_{i} \mid 0 \leqslant \lambda_{i}<1\right\}
$$

be the "fundamental parallelepiped" of C . Then, for every $v \in \mathbb{R}^{n}$ we have

$$
\sigma_{v+\mathrm{C}}\left(z_{1}, \ldots, z_{\mathrm{d}}\right)=\frac{\sigma_{v+\Pi}\left(z_{1}, \ldots, z_{\mathrm{d}}\right)}{\left(1-z^{\mathrm{m}_{1}}\right) \cdots\left(1-z^{\mathrm{m}_{\mathrm{d}}}\right)}
$$

Proof. The proof follows the template of the previous example - see [1].

### 7.2 Coning over polytopes

Definition 7.2.1. Let $P \subseteq \mathbb{R}^{d}$ be a convex polytope with vertices $x_{1}, \ldots, x_{k}$. We call standard cone over $P$ the polyhedron

$$
C(P):=\operatorname{cone}\left\{\left(x_{1}, 1\right), \ldots\left(x_{k}, 1\right)\right\} \subseteq \mathbb{R}^{\mathrm{d}+1}
$$

Now consider the ("vertical") vector $e_{d+1}:=(0, \ldots, 0,1) \in \mathbb{R}^{d+1}$. Given an integer $t$ we let

$$
\mathrm{H}^{(\mathrm{t})}:=\left\{\mathrm{y} \in \mathbb{R}^{\mathrm{d}+1} \mid\left\langle e_{\mathrm{d}+1}, \mathrm{y}\right\rangle=\mathrm{t}\right\}
$$

denote the "horizontal" plane at "height" $t$ and define

$$
\mathrm{P}^{(\mathrm{t})}:=\mathrm{C}(\mathrm{P}) \cap \mathrm{H}^{(\mathrm{t})} .
$$

Remark 7.2.2. Notice that, under the canonical identification of $H^{(t)}$ with $\mathbb{R}^{d}$, we have $P^{(1)}=P$ and, in general, $P^{(t)}=t P$, the $t$-th dilate of $P$. Thus, for the integer point counting function we have

$$
L_{P}(t)=L_{P^{(t)}}(1)
$$

and, since every integer point of $C(P)$ is obviously an integer point of exactly one $\mathrm{P}^{(\mathrm{t})}$,

$$
\sigma_{C(P)}(1, \ldots, 1, z)=1+\sum_{t \geqslant 1} L_{P(t)}(1) z^{t}=\sum_{t \geqslant 0} L_{P}(t) z^{t}=\operatorname{Ehr}_{P}(z)
$$

## Proof of Theorem 7.0.1

By Corollary 6.2.2 it is enough to prove that the power series satisfies

$$
\operatorname{Ehr}_{\mathrm{p}}(z)=\frac{p(z)}{(1-z)^{\mathrm{d}+1}}
$$

for a polynomial $p$ of degree at most $d$ and $p(1) \neq 0$.
Moreover, by Theorem 5.5 .8 it is enough to prove this for P an integral simplex (since any integral $P$ can be triangulated with no new vertices, and in particular the triangulation consists of integral simplices).

Then, let $P$ be an integral simplex - say with vertices $x_{1}, \ldots, x_{d+1} \in \mathbb{R}^{d}$. Then $C(P)$ is a simplicial $d$-cone generated by

$$
\mathfrak{m}_{1}:=\left(x_{1}, 1\right), \ldots, m_{d+1}=\left(x_{d+1}, 1\right) \in \mathbb{Z}^{\mathrm{d}+1}
$$

and now from Remark 7.2.2 and Theorem 7.1.3 that

$$
\operatorname{Ehr}_{\mathrm{P}}(z)=\sigma_{\mathrm{C}(\mathrm{P})}\left(1, \ldots 1, z_{\mathrm{d}+1}\right)=\frac{\sigma_{\Pi}\left(1, \ldots, 1, z_{\mathrm{d}+1}\right)}{(1-z)^{\mathrm{d}+1}}
$$

where $\Pi=\left\{\lambda_{1} m_{1}+\ldots+\lambda_{d+1} m_{d+1} \mid 0 \leqslant \lambda_{i}<1\right\}$ is the half-open fundamental parallepiped.

Now $\sigma_{\Pi}\left(1, \ldots 1, z_{\mathrm{d}+1}\right)$ is a polynomial in $z_{\mathrm{d}+1}$ (since $\Pi$ is bounded), and its value at $z_{\mathrm{d}+1}=1$ is not zero (in fact, $\sigma_{\Pi}(1, \ldots 1)$ is the total number of integer points in $\Pi$, and $\Pi$ contains at least the origin). We only have to show that the degree of $\sigma_{\Pi}\left(1, \ldots 1, z_{d+1}\right)$ is at most $d$. For this, notice that the maximum $z_{\mathrm{d}+1}$-degree in $\sigma_{\Pi}\left(z_{1}, \ldots, z_{\mathrm{d}+1}\right)$ is the maximum value of the last coordinate over all integer points in $\Pi$. Now the last coordinate of a generic point $\lambda_{1} m_{1}+\ldots+\lambda_{d+1} m_{d+1}$ in the parallelepiped is $\lambda_{1}+\ldots+\lambda_{d+1}$. Thus, every point in $\Pi$ has last coordinate strictly less than $d+1$, and so every integer point in $\Pi$ has last coordinate at most d .

## Exercises VII

VII. 1 For each of the following simplices compute the Ehrhart series and the Ehrhart polynomial.

- $\operatorname{conv}\{(0,0),(1,0),(0,1)\}$
- $\operatorname{conv}\{(0,0,0),(1,0,0),(0,2,0),(0,0,3)\}$
- $\operatorname{conv}\{(0,0,0,0),(1,0,0,0),(0,2,0,0),(0,0,3,0),(0,0,0,4)\}$
VII. 2 Compute the Ehrhart polynomial of the unit square $[0,1]^{2}$ in $\mathbb{R}^{2}$ using the theory developed in this chapter, and compare the result with the "common sense" computation.
VII. 3 Let G be the graph depicted below, oriented so that "every edge is oriented clockwise"


Compute the number of positive k-flows on G using Ehrhart theory, following the steps below:

- Consider the incidence matrix $A$ of $G$ and compute the vertices of the polytope $\mathrm{P}=\{A x=0, \mathbf{0} \leqslant x \leqslant \mathbf{1}\}$ (this is the polytope $\overline{\mathrm{U}}$ of the beginning of chapter V ).
(Hint: this polytope lives in the 2-dimensional space $W=\operatorname{ker}(A) \subseteq$ $\mathbb{R}^{4}$, hence the candidates for being vertices are the intersections of the lines $W \cap\left\{x_{i}=0\right\}, W \cap\left\{x_{i}=1\right\}, i=1, \ldots, 4$.)
- Compute the integer point enumerator $L_{Q}(t)$ for $Q$ ranging over all faces of $P$.
- Give a polynomial expression for the function $I_{G}(k)$, counting the number of integer points in $U$.

Bonus question: how many flippings of G support positive k-flows? Can you compute the "full" integer flow polynomial of G ?

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[^0]:    ${ }^{1}$ There are several such total orderings - for instance, we can consider the lexicographic ordering, which is defined by setting $B<B^{\prime}$ if and only if max $<\left(B \backslash B^{\prime} \cup B^{\prime} \backslash B\right) \in B^{\prime}$. (The name derives from the following alternative definition of this total order: for every broken circuit $B$ consider the ordered list $\lambda_{1}(\mathrm{~B})>\lambda_{2}(\mathrm{~B})>\ldots$ of its elements in decreasing order, and set $B<B^{\prime}$ if $\lambda_{i}(B)<\lambda_{i}\left(B^{\prime}\right)$ for the smallest $i$ where the lists differ.)

[^1]:    ${ }^{1}$ Unique by Remark 3.2 .2

[^2]:    ${ }^{1}$ https://en.wikipedia.org/wiki/Harold_Scott_MacDonald_Coxeter
    ${ }^{2}$ A reflection
    ${ }^{3}$ See for instance Bourbaki

[^3]:    ${ }^{1}$ Formal power series are formal objects, and one has to make sense of "setting $z$ to...". In this case, the subring $\mathcal{R}$ of $\mathbb{C}[[z]]$ generated by the $R_{i j}$ with $\mathfrak{i} \neq \mathfrak{i}_{0}$ and $\left(1-\gamma_{i_{0}} z\right)$ is in the image of the ring of rational functions $\mathbb{C}(z)$, and the map $\mathbb{C}(z) \rightarrow \mathbb{C}[[z]]$ is injective. Now every function in $\mathcal{R}$ is defined for $\frac{1}{\gamma_{i_{0}}}$ and thus the map $\mathcal{R} \rightarrow \mathbb{C}$ obtained by formally replacing $z$ with $\frac{1}{\gamma_{i_{0}}}$ is well-defined, so that an identity in $\mathcal{R}$ goes to an identity in $\mathbb{C}$.

