# Graph Polynomials and Spectral Gaps 

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#### Abstract

These notes are mainly based on a paper by McKay [4] on the eigenvalues of large regular graphs, one by Sokal [5] on bounds of zeros of chromatic polynomials, and another by Sarnak [3], found in sections 4, 8 and 6 respectively. It turns out that the eigenvalues of large regular graphs follow with high probability a distribution depending only the degree and that for cubic graphs, for any $x \neq 3$ there exists a sequence of graphs such that their spectra avoid $x$ (i.e. $x$ can be gapped). In addition, the zeros of the chromatic polynomial of a regular graph is bounded linearly by its degree. Sections 1,2 and 7 give some background on algebraic graph theory, mostly taken from Biggs [1]. Section 3, on the process of generating random regular graphs, is due to Bollobás [2].

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## 1 The spectrum of a graph

Definition 1.1. The adjacency matrix of $G$ is the $n \times n$ matrix $\mathbf{A}=\mathbf{A}(G)$ where

$$
a_{i, j}= \begin{cases}1 & v_{i} \text { and } v_{j} \text { are adjacent } \\ 0 & \text { otherwise }\end{cases}
$$

By definition, $\mathbf{A}$ is real symmetric and has trace 0 . Since the ordering of rows and columns of $\mathbf{A}$ was arbitrary, we will be interested in permutation-invariant properties of $\mathbf{A}$, mainly its spectral properties.

Definition 1.2. The spectrum of $G$ is the set of eigenvalues of $\mathbf{A}(G)$ with their multiplicities. If $\lambda_{0}>\lambda_{1}>$ $\cdots>\lambda_{s-1}$ are eigenvalues with multiplicity $m\left(\lambda_{0}\right), \ldots, m\left(\lambda_{s-1}\right)$ respectively, then we write

$$
\text { Spec } G=\left(\begin{array}{cccc}
\lambda_{0} & \lambda_{1} & \cdots & \lambda_{s-1} \\
m\left(\lambda_{0}\right) & m\left(\lambda_{1}\right) & \cdots & m\left(\lambda_{s-1}\right)
\end{array}\right)
$$

These eigenvalues correspond precisely to the roots of the characteristic polynomial, with their multiplicity equal to the multiplicity of these roots.

[^0]Definition 1.3. The characteristic polynomial of $G$ is

$$
\chi_{G}(x)=\operatorname{det}(x \mathbf{I}-\mathbf{A})=x^{n}+c_{1} x^{n-1}+\cdots+c_{n}
$$

Note that $c_{i}$ is the sum of all $i \times i$ principal minors of $\mathbf{A}$, which is the determinant of a submatrix of $\mathbf{A}$ that takes $i$ rows and the same set of columns, since $c_{i}$ comes from expanding along the diagonal $n-i$ times, which always leaves such a matrix behind.

Proposition 1.4. The coefficients of $\chi_{G}$ satisfy

1. $c_{1}=0$;
2. $-c_{2}=|E|$;
3. $-c_{3}$ is twice the number of triangles of $G$.

Proof. Essentially follows from the fact that those $c_{i}$ are summed from principal minors.

1. All $1 \times 1$ principal minors of $\mathbf{A}$ are 0 ;
2. A $2 \times 2$ principal minor of $\mathbf{A}$ at $i, j$ is non-zero if and only if the submatrix is $\left[\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right]$, in which case the principal minor is -1 , and we have such a submatrix if and only if $v_{i}$ and $v_{j}$ are adjacent;
3. A $3 \times 3$ principal minor of $\mathbf{A}$ at $i, j, k$ is 0 unless all non-diagonal entries are 1 , which is true if and only if $v_{i}, v_{j}$ and $v_{k}$ form a triangle. In this case, the principal minor is $\left|\begin{array}{lll}0 & 1 & 1 \\ 1 & 0 & 1 \\ 1 & 1 & 0\end{array}\right|=2$.

This shows that algebraic properties of $\mathbf{A}(G)$ have profound implications on graph-theoretical properties of $G$ and vice-versa. Here we apply them to find Spec $K_{4}$.

Example 1.5. For $G=K_{4}$, note

$$
\mathbf{A}=\left[\begin{array}{llll}
0 & 1 & 1 & 1 \\
1 & 0 & 1 & 1 \\
1 & 1 & 0 & 1 \\
1 & 1 & 1 & 0
\end{array}\right] \Longrightarrow \mathbf{A}\left[\begin{array}{l}
1 \\
1 \\
1 \\
1
\end{array}\right]=\left[\begin{array}{l}
3 \\
3 \\
3 \\
3
\end{array}\right]
$$

so we see that 3 is an eigenvalue. We can use 1.4 to derive all but the last coefficient of $\chi$, which is simply $\operatorname{det} \mathbf{A}=-3$. Thus

$$
\begin{aligned}
& \chi_{K_{4}}(x)=x^{4}-6 x^{2}-8 x-3=(x-3)\left(x^{3}+3 x^{2}+3 x+1\right)=(x-3)(x+1)^{3} \\
\Longrightarrow & \operatorname{Spec} K_{4}=\left(\begin{array}{cc}
3 & -1 \\
1 & 3
\end{array}\right) .
\end{aligned}
$$

Later, in 2.6, we will see a general description of Spec $K_{n}$.

Another interesting algebraic object is the adjacency algebra:

Definition 1.6. The adjacency algebra of $G$ is the algebra generated by polynomials of $\mathbf{A}=\mathbf{A}(G)$, with the usual matrix addition and multiplication. We denote the adjacency algebra by $\mathcal{A}(G)$.

Since $\chi_{G}(G)=0, \operatorname{dim} \mathcal{A}(G) \leq n$. We can derive a lower bound for $\operatorname{dim} \mathcal{A}(G)$ via graph-theoretical properties of $G$. We start by associating powers of $\mathbf{A}$ to graph-theoretical properties.

Definition 1.7. A walk from $v_{i}$ to $v_{j}$ of length $\ell$ is a sequence of vertices $v_{i}=v_{0}, v_{1}, \ldots, v_{\ell}=v_{j}$ such that $v_{k-1}$ and $v_{k}$ are adjacent for all $1 \leq k \leq \ell$.

Lemma 1.8. The number of walks of length $\ell$ from $v_{i}$ to $v_{j}$ is precisely the entry $(i, j)$ in $\mathbf{A}^{\ell}$.

Proof. We prove this by induction. $\ell=0$ is true since $\mathbf{A}^{0}=\mathbf{I}$. Suppose this is true for $\ell=L$. Note that for each distinct walk from $v_{i}$ to $v_{j}$ of length $L+1$, there is a distinct walk of length $L$ from $v_{i}$ to some $v_{k}$ adjacent to $v_{j}$ and vice-versa. Thus the number of walks from $v_{i}$ to $v_{j}$ of length $L+1$ is precisely the sum of the number of walks of length $L$ from $v_{i}$ to each $v_{k}$ adjacent to $v_{j}$, i.e. $\left(\mathbf{A}^{L}\right)_{i k}$. And the $(i, j)$ entry of $\mathbf{A}^{L+1}$ is precisely

$$
\left(\mathbf{A}^{L+1}\right)_{i j}=\sum_{k=1}^{n}\left(\mathbf{A}^{L}\right)_{i k} a_{k j}=\sum_{\left(v_{k}, v_{j}\right) \in E}\left(\mathbf{A}^{L}\right)_{i k}
$$

Definition 1.9. We call $G$ connected if there exists a walk from each $v_{i}$ to each $v_{j}$. When there exists a walk from $v_{i}$ to $v_{j}$, the length of the shortest walk is called the distance from $v_{i}$ to $v_{j}$ on $G$, denoted $\partial\left(v_{i}, v_{j}\right)$. The maximum distance on a connected $G$ is called its diameter.

Proposition 1.10. Let $G$ be connected with diameter $d$. Then $\operatorname{dim} \mathcal{A}(G) \geq d+1$.

Proof. Let $x, y \in V$ be such that $\partial(x, y)=d$. Let $x=v_{0}, \ldots, v_{d}=y$ be a walk of length $d$. Then there are no walks from $x$ to $v_{\ell}$ that are shorter than $\ell$, so that entry in $\mathbf{A}^{k}$ is 0 for all $k<\ell$ so $\mathbf{A}^{\ell}$ is independent from $\left\{\mathbf{I}, \ldots, \mathbf{A}^{\ell-1}\right\}$ and this is true for all $0 \leq \ell \leq d$ so $\left\{\mathbf{I}, \ldots, \mathbf{A}^{d}\right\}$ are linearly independent, i.e. $\operatorname{dim} \mathcal{A}(G) \geq d+1$.

Since the dimension of $\mathcal{A}(G)$ corresponds directly to the degree of the minimal polynomial, which corresponds to the number of distinct eigenvalues of $G$, we have the following corollary:

Corollary 1.11. A connected graph with diameter $d$ has at least $d+1$ distinct eigenvalues.

Exercise 1.1. If $G_{i}$ denotes the induced subgraph of $V \backslash v_{i}$, then

$$
\chi_{G}^{\prime}=\sum_{i=1}^{n} \chi_{G_{i}} .
$$

Proof. Let $\mathbf{X}(x)$ be the diagonal matrix with $x_{i}(x)$ at the $i$-th row, where $x_{i}: x \mapsto x$ and let $\mathbf{X}_{i}(x)$ be $\mathbf{X}$ with the $i$-th row and column removed. Now we view $\chi_{G}$ as

$$
\chi_{G}(x)=\operatorname{det}(x \mathbf{I}-\mathbf{A}(G))=\operatorname{det}(\mathbf{X}(x)-\mathbf{A}(G))
$$

Then by chain rule ( $g_{i}$ denotes the part of $\operatorname{det}(\mathbf{X}(x)-\mathbf{A}(G))$ expanded at the $i$-th row that does not depend on $x_{i}$ )

$$
\begin{aligned}
\frac{\partial \chi_{G}}{\partial x}(x) & =\sum_{i=1}^{n} \frac{\partial \chi_{G}}{\partial x_{i}}(x) \cdot \frac{\partial x_{i}}{\partial x}(x)=\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} \operatorname{det}(\mathbf{X}(x)-\mathbf{A}(G)) \\
& =\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}}\left(x_{i} \operatorname{det}\left(\mathbf{X}_{i}(x)-\mathbf{A}\left(G_{i}\right)\right)-g_{i}\left(x_{1}, \ldots, x_{i-1}, x_{i+1}, \ldots, x_{n}\right)\right) \\
& =\sum_{i=1}^{n} \operatorname{det}\left(\mathbf{X}_{i}(x)-\mathbf{A}\left(G_{i}\right)\right)=\sum_{i=1}^{n} \chi_{G_{i}}
\end{aligned}
$$

Exercise 1.2. Let $g_{i j}(r)$ be the number of walks of length $r$ from $v_{i}$ to $v_{j}$ in $G$. Let $\mathbf{G}(z)$ be the matrix where

$$
\begin{equation*}
(\mathbf{G}(z))_{i j}=\sum_{r=0}^{\infty} g_{i j}(r) z^{r} \tag{1.1}
\end{equation*}
$$

under the assumption that (1.1) is absolutely convergent for all $i, j$. Then $\mathbf{G}(z)=(\mathbf{I}-z \mathbf{A})^{-1}$ (so in particular $\left.\frac{1}{z} \notin \operatorname{Spec} G\right)$.

Proof. We prove $\mathbf{G}(z)=\mathbf{G}(z) z \mathbf{A}+\mathbf{I}$. Indeed we have

$$
(\mathbf{G}(z) z \mathbf{A})_{i j}=z \sum_{k=1}^{n}(\mathbf{G}(z))_{i k} a_{k j}=z \sum_{k=1}^{n} \sum_{r=0}^{\infty} g_{i k}(r) z^{r} a_{k j}=\sum_{r=0}^{\infty} z^{r+1} \sum_{k=1}^{n} g_{i k}(r) a_{k j}
$$

Recall from the proof of 1.8 we have $\sum_{k=1}^{n} g_{i k}(r) a_{k j}=g_{i j}(r+1)$, thus

$$
(\mathbf{G}(z) z \mathbf{A})_{i j}=\sum_{r=0}^{\infty} z^{r+1} g_{i j}(r+1)=\sum_{r=1}^{\infty} z^{r} g_{i j}(r)
$$

If $i \neq j$ then $g_{i j}(0)=0$ so $(\mathbf{G}(z) z \mathbf{A})_{i j}=(\mathbf{G}(z))_{i j}$. If $i=j$ then $z^{0} g_{i j}(0)=1$ so $(\mathbf{G}(z) z \mathbf{A})_{i j}+1=(\mathbf{G}(z))_{i j}$. Thus $\mathbf{G}(z)=\mathbf{G}(z) z \mathbf{A}+\mathbf{I}$.

## 2 Regular graphs

We will see in this section that combinatorial regularity can have consequences on spectra of graphs. The class of graphs possessing the most simple kind of regularity are $k$-regular graphs, or graphs in which every vertex has degree $k$. This property has direct consequences on the graphs' eigenvalues.

Proposition 2.1. Let $G$ be $k$-regular. Then

1. $k \in \operatorname{Spec} G$;
2. The multiplicity of $k \in \operatorname{Spec} G$ is 1 if $G$ is connected;
3. For all $\lambda \in \operatorname{Spec} G,|\lambda| \leq k$.

Proof. 1. Take $\mathbf{x}=[1, \ldots, 1]^{\top} \in \mathbb{R}^{n}$, then $(\mathbf{A} \mathbf{x})_{i}=\sum_{j=1}^{n} a_{i j}=k$ so $\mathbf{A} \mathbf{x}=k \mathbf{x}$.
2. Suppose $\mathbf{x}=\left[x_{1}, \ldots, x_{n}\right]^{\top}$ satisfies $\mathbf{A x}=k \mathbf{x}$. Let $i$ be such that $x_{i}=\max _{1 \leq j \leq n}\left(x_{j}\right)$. Then

$$
(\mathbf{A} \mathbf{x})_{i}=\sum_{j=1}^{n} a_{i j} x_{j}=k x_{i}
$$

so $x_{i}=x_{j}$ for all $x_{j}$ adjacent to $x_{i}$. Since now those $x_{j}=\max _{1 \leq j \leq n}\left(x_{j}\right)$ also, their neighbours are also equal to $x_{i}$ and so forth. Since $G$ is connected, we get that $x_{i}=x_{j} \forall i, j$ so $\mathbf{x}=x_{1}[1, \ldots, 1]^{\top}$.
3. Suppose $\mathbf{x}=\left[x_{1}, \ldots, x_{n}\right]^{\top} \neq 0$ satisfies $\mathbf{A x}=\lambda \mathbf{x}$. Let $i$ be such that $\left|x_{i}\right|=\max _{1 \leq j \leq n}\left(\left|x_{j}\right|\right)$. Then

$$
\left|\lambda x_{i}\right|=\left|(\mathbf{A x})_{i}\right|=\left|\sum_{j=1}^{n} a_{i j} x_{j}\right| \leq \sum_{j=1}^{n} a_{i j}\left|x_{j}\right| \leq \sum_{j=1}^{n} a_{i j}\left|x_{i}\right|=k\left|x_{i}\right| \Longrightarrow|\lambda| \leq k
$$

In particular, 2.1 gives us a way to characterize regular connected graphs from its adjacency algebra. Let $\mathbf{J}$ denote the matrix whose entries are all 1.

Proposition 2.2. $G$ is regular and connected if and only if $\mathbf{J} \in \mathcal{A}(G)$.

Proof. Suppose $\mathbf{J} \in \mathcal{A}(G)$. Then for each $(i, j)$ there exists $r$ such that $\left(\mathbf{A}^{r}\right)_{i j} \neq 0$ so by $1.8 G$ is connected. Since $\mathbf{J}$ is a polynomial of $\mathbf{A}, \mathbf{J} \mathbf{A}=\mathbf{A} \mathbf{J}$ so we have

$$
(\mathbf{J A})_{i j}=\sum_{k=1}^{n} a_{k j}=(\mathbf{A J})_{i j}=\sum_{k=1}^{n} a_{i k}
$$

i.e. $\operatorname{deg}\left(v_{j}\right)=\operatorname{deg}\left(v_{i}\right)$ and this is for all $v_{i}, v_{j} \in V$ so $G$ is regular.

Conversely suppose $G$ is $k$-regular and connected. Let $p$ be the minimal polynomial of $G$. Then by 2.1 $p(x)=(x-k) q(x)$ for some polynomial $q$. Note $p(\mathbf{A})=0=(\mathbf{A}-k \mathbf{I}) q(\mathbf{A}) \Longrightarrow \mathbf{A} q(\mathbf{A})=k q(\mathbf{A})$. In particular, this means that every column vector of $q(\mathbf{A})$ is a $k$-eigenvector of $\mathbf{A}$. Since $A$ is connected, by 2.1 the columns of $q(\mathbf{A})$ are multiples of $[1, \ldots, 1]^{\top}$. Since $(q(\mathbf{A}))^{\top}=q(\mathbf{A}), q(\mathbf{A})$ is a multiple of $\mathbf{J}$.

Note that the polynomial $q$ used in 2.2 is simply $q(x)=\prod_{i=1}^{s-1}\left(x-\lambda_{i}\right)$, where $\lambda_{i}$ are the distinct eigenvalues of $G$ that are smaller than $k$. The following is an explicit construction of $\mathbf{J}$ from $q$ (i.e. from $\operatorname{Spec} G$ ) given a connected $k$-regular $G$.

Corollary 2.3. Let $G$ be connected and $k$-regular and let $q$ be as in the proof of 2.2. Then

$$
\mathbf{J}=\left(\frac{n}{q(k)}\right) q(\mathbf{A})
$$

Proof. From the proof of $2.2, q(\mathbf{A})=\alpha \mathbf{J}$ for some $\alpha$. Let $\mathbf{x}=[1, \ldots, 1]^{\top} \in \mathbb{R}^{n}$. Note that $\alpha \mathbf{J} \mathbf{x}=n \alpha \mathbf{x}$. However since $\mathbf{A x}=k \mathbf{x}$ we must have $q(\mathbf{A}) \mathbf{x}=q(k) \mathbf{x}\left(\right.$ since $\left.\mathbf{A}^{r} \mathbf{x}=k^{r} \mathbf{x}\right)$. Thus $\alpha=\frac{q(k)}{n}$ so $\mathbf{J}=\frac{1}{\alpha} q(\mathbf{A})=$ $\left(\frac{n}{q(k)}\right) q(\mathbf{A})$.

This shows that knowing $\operatorname{Spec} G$ can be particularly powerful. The following builds the linear algebraic theory for a special class of regular graphs for which we can compute Spec $G$ explicitly.

Definition 2.4. An $n \times n$ matrix $\mathbf{M}$ is circulant if $m_{i j}=m_{1,(j-i+1)_{n}}$ for $i>1$. A graph $G$ is circulant if $\mathbf{A}(G)$ is circulant.

In other words, $\mathbf{M}$ is circulant if the $i$-th row of $M$ is its first row slid right $i-1$ places. Let $\mathbf{W}$ be the $n \times n$ circulant matrix where the first row is $[0,1,0, \ldots, 0]$. Then if the first row of $\mathbf{M}$ is $\left[m_{1}, \ldots, m_{n}\right]$, we have

$$
\mathbf{M}=\sum_{j=1}^{n} m_{j} \mathbf{W}^{j-1}
$$

(since $\mathbf{W}^{j-1}$ is a circulant matrix where the first row has 1 at the $j$-th column and 0 elsewhere).
Suppose $\mathbf{x}=\left[x_{1}, \ldots, x_{n}\right]^{\top}$ is a non-zero eigenvector of $\mathbf{W}$, then note

$$
\mathbf{W} \mathbf{x}=\left[x_{2}, \ldots, x_{n}, x_{1}\right]^{\top}=\lambda \mathbf{x} \Longleftrightarrow \lambda x_{j}=x_{j+1} \forall 1 \leq j<n \wedge \lambda x_{n}=x_{1} \Longleftrightarrow \lambda^{n}=1 \wedge x_{j}=\lambda^{j-1} x_{1}
$$

In other words, the eigenvalues of $\mathbf{W}$ are precisely the $n$-th roots of 1 , i.e. $1, \omega, \ldots, \omega^{n-1}$ where $\omega=$ $\exp (2 \pi i / n)$. It follows that the eigenvalues of $\mathbf{M}$, which is a polynomial of $\mathbf{W}$, are

$$
\lambda_{r}=\sum_{j=1}^{n} m_{j} \omega^{(j-1) r}, 0 \leq r \leq n-1
$$

This gives us the following result:

Proposition 2.5. Suppose that $\left[0, a_{2}, \ldots, a_{n}\right]$ is the first row of $\mathbf{A}(G)$ for $G$ circulant, then the eigenvalues of $G$ are

$$
\lambda_{r}=\sum_{j=2}^{n} a_{j} \omega^{(j-1) r}, 0 \leq r \leq n-1
$$

We remark that these $n$ eigenvalues are not necessarily distinct. We will give three applications of 2.5 .

Example 2.6. $K_{n}$ is circulant, with the first row of $\mathbf{A}\left(K_{n}\right)$ being $[0,1, \ldots, 1]$. Note $\sum_{j=1}^{n} \omega^{(j-1) r}=0$ for $1 \leq r \leq n-1$, so in these cases

$$
\lambda_{r}=\sum_{j=2}^{n} a_{j} \omega^{(j-1) r}=-1
$$

For $r=0, \omega^{(j-1) r}=1 \forall 1 \leq j \leq n$ so

$$
\lambda_{0}=\sum_{j=2}^{n} a_{j} \omega^{(j-1) r}=n-1
$$

Thus

$$
\operatorname{Spec} K_{n}=\left(\begin{array}{cc}
n-1 & -1 \\
1 & n-1
\end{array}\right)
$$

Example 2.7. $C_{n}$ is circulant, with the first row of $\mathbf{A}\left(C_{n}\right)$ being $[0,1,0, \ldots, 0,1]$ (so $v_{i}$ is adjacent to $v_{i-1}$ and $v_{i+1}$ and $v_{1}$ is adjacent to $v_{n}$ ). Then

$$
\lambda_{r}=\sum_{j=2}^{n} a_{j} \omega^{(j-1) r}=\omega^{r}+\omega^{(n-1) r}=\cos (2 \pi r / n)+\cos (-2 \pi r / n)+i(\sin (2 \pi r / n)+\sin (-2 \pi r / n))=2 \cos (2 \pi r / n)
$$

However note that $2 \cos (2 \pi r / n)=2 \cos (2 \pi(n-r) / n)$ so if $n$ is odd then for all $1 \leq r \leq \frac{n-1}{2}$ there is a distinct $\frac{n+1}{2} \leq r^{\prime} \leq n-1$ such that $n-r=r^{\prime}$, thus

$$
\operatorname{Spec} C_{n}=\left(\begin{array}{cc}
2 & 2 \cos (2 \pi r / n) \forall 1 \leq r \leq \frac{n-1}{2} \\
1 & 2
\end{array}\right)
$$

On the other hand if $n$ is even then for all $1 \leq r \leq \frac{n}{2}-1$ there exists a distinct $\frac{n}{2}+1 \leq r^{\prime} \leq n-1$ such that $n-r=r^{\prime}$ and $2 \cos (2 \pi(n / 2) / n)=2 \cos \pi=-2$, thus

$$
\text { Spec } C_{n}=\left(\begin{array}{ccc}
2 & 2 \cos (2 \pi r / n) \forall 1 \leq r \leq \frac{n-2}{2} & -2 \\
1 & 2 & 1
\end{array}\right)
$$

Example 2.8. For $K_{2 s}$, separate $V$ into $V_{1}, V_{2}$ with $V_{1} \cup V_{2}=V$ and $\left|V_{1}\right|=\left|V_{2}\right|=s$ and pair $V_{1}$ and $V_{2}$ bijectively. Define $H_{s}$, called the hyperoctahedral graph, to be the subgraph of $K_{2 s}$ with the edges between those paired vertices removed. We now order $V\left(H_{s}\right)$ such that $\left(v_{i}, v_{i+s}\right) \notin E\left(H_{s}\right) \forall 1 \leq i \leq s$ (i.e. each $v_{i} \in V_{1}$ was paired to $\left.v_{i+s} \in V_{2}\right)$. Then we note that $\mathbf{A}\left(H_{s}\right)$ is circulant with the first row being all 1 s except at the first and the $s+1$-th positions. Again due to $\sum_{j=1}^{2 s} \omega^{(j-1) r}=0$ for $1 \leq r \leq 2 s-1$, in these cases

$$
\lambda_{r}=\sum_{j=2}^{2 s} a_{j} \omega^{(j-1) r}=-1-\omega^{s r}=-1-(-1)^{r}= \begin{cases}0 & r=2 k-1,1 \leq k \leq s \\ -2 & r=2 k, 1 \leq k \leq s-1\end{cases}
$$

When $r=0, \omega^{(j-1) r}=1 \forall 2 \leq j \leq 2 s$ so $\lambda_{0}=2 s-2$. Thus

$$
\text { Spec } H_{s}=\left(\begin{array}{ccc}
2 s-2 & 0 & -2 \\
1 & s & s-1
\end{array}\right)
$$

## 3 Random Regular Graphs

A commonly used method to generate a random graph on $n$ vertices ${ }^{2}$ is $G(n, 1 / 2)$, where we take a graph with $n$ vertices and each pair of vertices have independently a probability of $1 / 2$ to be adjacent. Under this model, all $2^{\binom{n}{2}}$ graphs have an equal chance of being chosen. It is also an effective algorithm for generating random graphs. Here, we define another algorithm for generating regular graphs because, as we will see later, a random graph is, with high probability, not regular.

Definition 3.1. Given $n>k$ natural numbers with $n k$ being even, we define a random $k$-regular graph on $n$ vertices generated by the Bollobás model as follows:

- Take $n k$ balls grouped into $n$ sets of $k$ balls;
- Uniformly randomly take a perfect matching on those balls. If two balls from the same set matched, or if there are two or more matches between balls from the same pair of sets, then we say the generation failed and take a new matching;
- For each set of balls, we define a vertex. Two vertices are adjacent if there is a matching between balls of their corresponding set.

Suppose for each matching, we label all sets and all balls. For each regular graph, we label all vertices and label each edge by a 2 -tuple corresponding to its two ends, such that the set of tuple values on edges out of any given vertex form exactly some $k$-element set (in Figure $1,\{a, b, c\}$ ). Then it is clear that there is a bijection between labeled graphs and labeled matchings.

graph

1a with 5 a
1 b with 2c
1c with 4 b
2 a with 6 a
2 b with 3 b
3 a with 6 b
3c with 4a
4c with 5b
5 c with 6 c
matching

graph with new labelling

1a with 2c
1 b with 5 a
1c with $4 b$ 2a with 6a 2 b with 3b 3 a with 6 b 3c with 4a 4c with 5b 5 c with 6 c
new matching

Figure 1: Two ways of labelling a 4-regular graph on 6 vertices and its corresponding matching.
The assignment of a labelling to matchings is clearly uniform by simply labelling the balls before matching. For a given vertex labelling (there are $n$ ! ways to vertex label a graph), a permutation in the way any vertex labels its incident edges will result in a different graph labelling, so there are exactly $n^{k!}$ ways to edge label a graph given a vertex labelling. Thus all $k$-regular graphs are equally likely to be obtained from a random matching, giving us the following:

[^1]Proposition 3.2. The Bollobás model generates all $k$-regular graphs on $n$ vertices with equal probability.

We note that the Bollobás model gives us an algorithm for generating random $k$-regular graphs. Now we prove that a random graph is not $k$-regular w.h.p.

## Proposition 3.3. $G(n, 1 / 2)$ is not $k$-regular w.h.p.

Proof. There are $2 \begin{gathered}\binom{n}{2}\end{gathered}$ graphs on $n$ labeled vertices. Given $n$ groups of $k$ balls, there are $\binom{n k}{n k / 2}$ ways to bipartition the balls and ( $n k / 2$ )! ways to make a matching. However, flipping any pair of matched balls will be counted as a new matching in this system, so each matching is actually counted $2^{n k / 2}$ times. Thus there are $\frac{(n k)!}{(n k / 2)!2^{n k / 2}}$ matchings. As noted above, each $k$-regular graph can be produced from $n^{k!}$ matchings, so there are $\frac{(n k)!}{(n k / 2)!2^{n k / 2} n^{k!}} k$-regular graphs on $n$ labeled vertices. Finally for $n \geq k / 2$

$$
\frac{(n k)!}{(n k / 2)!2^{n k / 2} n^{k!}} \leq\left(\frac{k}{2}\right)^{n k / 2} \cdot n^{n k / 2-k!} \leq n^{n k} \lll 2^{\binom{n}{2}}=(\sqrt{2})^{n^{2}-n}
$$

## 4 Distribution of eigenvalues of a large regular graph

This section will be notes on a paper by McKay 1981.

### 4.1 Introduction

Given $G$ (proper) $k$-regular graph, let $n(G)$ denote its number of vertices and $c_{r}(G)$ its number of cycles of length $r$. Let $F_{G}(x)$ denote the discrete CDF of its eigenvalues (i.e. each eigenvalue $\lambda$ is assigned $p(\lambda)=\frac{m(\lambda)}{n(G)}$ where $m(\lambda)$ is its multiplicity). Then by basic probability

1. $F_{G}(x)=0$ if $x<-k$;
2. $F_{G}(x)=1$ if $x \geq k$;
3. $F_{G}(x)$ is increasing and right-continuous on $\mathbb{R}$.

Now we state the main result of the McKay paper:

Theorem 4.1. Let $G_{1}, G_{2}, \ldots$ be a sequence of $k$-regular graphs satisfying the following:

- $n\left(G_{i}\right) \rightarrow \infty$ as $i \rightarrow \infty$;
- For each $r \geq 3$,

$$
\begin{equation*}
c_{r}\left(G_{i}\right) / n\left(G_{i}\right) \xrightarrow{i \rightarrow \infty} 0 . \tag{4.1}
\end{equation*}
$$

Then, $F_{G_{i}}(x) \rightarrow F(x)$ for $x \in \mathbb{R}$ as $i \rightarrow \infty$, where $F(x)$ is the function defined as follows:

1. $F(x)=0$ if $x \leq-2 \sqrt{k-1}$;
2. If $-2 \sqrt{k-1}<x<2 \sqrt{k-1}$,

$$
F(x)=\int_{-2 \sqrt{k-1}}^{x} \frac{k \sqrt{4(k-1)-t^{2}}}{2 \pi\left(k^{2}-t^{2}\right)} d t=\frac{1}{2}+\frac{k}{2 \pi}\left[\arcsin \frac{x}{2 \sqrt{k-1}}-\frac{k-1}{k} \arctan \frac{(k-2) x}{k \sqrt{4(k-1)-x^{2}}}\right]
$$

3. $F(x)=1$ if $x \geq 2 \sqrt{k-1}$.

Conversely, if $F_{G_{i}}(x)$ does not converge to $F(x)$ for some $x$, then (4.1) fails for some $r \geq 3$.

All graphs in this section take the assumptions of 4.1. We start with a short lemma or remark:
Lemma 4.2. Let $c\left(G_{i}\right)$ denote the number of connected components in $G_{i}$. Then $c\left(G_{i}\right) / n\left(G_{i}\right) \rightarrow 0$ as $i \rightarrow \infty$.

Proof. For contradiction suppose $\limsup _{i \rightarrow \infty} c\left(G_{i}\right) / n\left(G_{i}\right)=c>0$, so by taking a subsequence of $G_{i}$, we can WLOG assume that $c\left(G_{i}\right) / n\left(G_{i}\right) \rightarrow c$. Fix $\min (c / 2,1)>\epsilon>0$. Let $N \in \mathbb{N}$ be such that, for all $2(1 / c+\epsilon) \leq r \leq 3, c_{r}\left(G_{i}\right) / n\left(G_{i}\right)<\epsilon,\left|n\left(G_{i}\right) / c\left(G_{i}\right)-1 / c\right|<\epsilon$ and $\left|c\left(G_{i}\right) / n\left(G_{i}\right)-c\right|<\epsilon$ for all $i \geq N$.

Now fix $i \geq N$. Note that since each component in $G_{i}$ has on average at most $1 / c+\epsilon$ vertices, at least $c\left(G_{i}\right) / 2$ components have at most $2(1 / c+\epsilon)$ vertices. Note that each such component contains at least one cycle of length $2(1 / c+\epsilon) \leq r \leq 3$ (due to acyclic graphs being trees, which always have vertices of degree 1 or 0 so not $k$-regular). Thus each such cycle appears on average at least $c\left(G_{i}\right) /(4(1 / c+\epsilon))$ times, so there must exist some $2(1 / c+\epsilon) \leq r \leq 3$ such that

$$
c_{r}\left(G_{i}\right) \geq c\left(G_{i}\right) /(4(1 / c+\epsilon)) \Longrightarrow \epsilon>\frac{c_{r}\left(G_{i}\right)}{n\left(G_{i}\right)} \geq \frac{c\left(G_{i}\right)}{4(1 / c+\epsilon) n\left(G_{i}\right)}>\frac{c-\epsilon}{4(1 / c+\epsilon)}>\frac{c}{8(1 / c+1)}
$$

which clearly does not hold as $\epsilon \rightarrow 0$.

Corollary 4.3. For all $M \in \mathbb{N}$ and $c>0$, there exists $N \in \mathbb{N}$ such that for all $i \geq N$, components of $G_{i}$ with less than $M$ vertices have in total at most $n\left(G_{i}\right) / c$ vertices.

Proof. Components with less than $M$ vertices have in total at most $M c\left(G_{i}\right)$ vertices, and

$$
M c\left(G_{i}\right) \leq n\left(G_{i}\right) / c \Longleftrightarrow M c\left(G_{i}\right) / n\left(G_{i}\right) \leq 1 / c
$$

Thus we can simply apply 4.2 with $\epsilon \leq \frac{1}{c M}$.

Example 4.4. When $k=2, G_{i}$ are all cycles $C_{r(i)}$ where $r(i)=n\left(G_{i}\right) \rightarrow \infty$ and for $-2<x<2$,

$$
F(x)=\int_{-2}^{x} \frac{2 \sqrt{4-t^{2}}}{2 \pi\left(4-t^{2}\right)} d t=\int_{-2}^{x} \frac{1}{\pi \sqrt{4-t^{2}}} d t=\frac{1}{\pi}\left(\arcsin \frac{x}{2}+\frac{\pi}{2}\right)=\frac{1}{\pi} \arccos \frac{x}{2}
$$

If $X \sim \mathcal{U}(0,1)$ and $Y=2 \cos (\pi X)=2 \cos (2 \pi X)$, then $F(x)$ is exactly the CDF of $Y$. It is known that each component of $G_{i}$ must be a cycle. For each $r \geq 3$, let $F_{r}$ denote the distribution of eigenvalues of $C_{r}$. Then by 2.7 we see that $F_{r}(x) \rightarrow F(x)$ for all $x \in \mathbb{R}$. 4.3 allows us to assume that at least $(1-1 / c) n\left(G_{i}\right)$ eigenvalues are from cycles of size at least $M$, for large enough $i$, so it is clear that $F_{G_{i}}(x)$ is between $F_{M}(x)-1 / c$ and $\max \left(F(x), F_{M}(x)\right)+1 / c$ for each $x$, so with $M \rightarrow \infty$ and $c \rightarrow 0$ we have $F_{G_{i}}(x) \rightarrow F(x)$.

### 4.2 Existence and uniqueness of $F(x)$

We first prove a counting result related to Catalan numbers.

Lemma 4.5. Fix $n, k \in \mathbb{N}, k \leq n$. Assume we have $\mathbf{x}=\left(x_{1}, \ldots, x_{2 n}\right)$ where $x_{i} \in\{-1,1\}$ such that the following are satisfied:

- $\sum_{i=1}^{2 n} x_{i}=0$ and $\sum_{i=1}^{j} x_{i} \geq 0 \forall 1 \leq j \leq 2 n$ (note that these are the same as for Catalan numbers);
- $\left|\left\{1 \leq j \leq 2 n: \sum_{i=1}^{j} x_{i}=0\right\}\right|=k$.

Then there are

$$
\binom{2 n-k}{n} \frac{k}{2 n-k}
$$

possibilities for $\mathbf{x}$.

Proof. Let $N(n, k)$ denote the number of possibilities for $\mathbf{x}$ and let $P(n, k)=\binom{n+k-1}{k-1}$ denote the number of ways of putting $n$ indistinguishable balls into $k$ labeled baskets (equivalent to choosing $k-1$ positions out of $n+k-1$ to be cutoffs between baskets, with the rest being balls in each basket). Let $\mathbf{x}^{\prime}$ be defined from $\mathbf{x}$ according to the following:

- For each $1<j<2 n$ where $\sum_{i=1}^{j} x_{i}=0$, remove $x_{j}$ and $x_{j+1}$ (note that $x_{j}=-1$ and $x_{j+1}=1$ always);
- Remove $x_{1}$ and $x_{2 n}\left(x_{1}=1\right.$ and $x_{2 n}=-1$ always $)$.

Then $\mathbf{x}^{\prime}$ satisfies the assumptions with $n-k$ and some $\ell$ as its parameters $(1 \leq \ell \leq n-k)$. However, given $\mathbf{x}^{\prime}$, there are $P(k-1, \ell+1)$ possible $\mathbf{x}$ (since there are $k-1$ "zeros" to be put into $\ell+1$ potential positions). Thus we have

$$
N(n, k)=\sum_{\ell=1}^{n-k} N(n-k, \ell) P(k-1, \ell+1)=\sum_{\ell=1}^{n-k}\binom{\ell+k-1}{\ell} N(n-k, \ell) .
$$

We now prove by induction on $n-k$ that $N(n, k)=\binom{2 n-k}{n} \frac{k}{2 n-k}$. If $n=k$ then $N(n, k)=1$. Now assume the proposition is true for all $n, k$ such that $n-k<M(M \in \mathbb{N})$ and now assume $n-k=M$. By IH (note $n-k-\ell<M)$

$$
\begin{aligned}
N(n, k) & =\sum_{\ell=1}^{M}\binom{\ell+k-1}{\ell} N(M, \ell)=\sum_{\ell=1}^{M}\binom{\ell+k-1}{\ell}\binom{2 M-\ell}{M} \frac{\ell}{2 M-\ell} \\
& =\sum_{\ell=1}^{M} \frac{(\ell+k-1)!}{\ell!(k-1)!} \cdot \frac{(M+n-k-\ell)!}{(n-k)!(M-\ell)!} \cdot \frac{\ell}{M+n-k-\ell} \\
& =\frac{(2 n-k)!}{n!(n-k)!} \frac{k}{2 n-k} \sum_{\ell=1}^{M} \frac{(\ell+k-1)!(M-\ell+n-k-1)!n!}{(\ell-1)!k!(M-\ell)!(M+n-1)!}
\end{aligned}
$$

Note

$$
\begin{aligned}
\frac{(\ell+k-1)!(M-\ell+n-k-1)!M!}{(\ell-1)!k!(M-\ell)!(M+n-1)!} & =\frac{\binom{\ell+k-1}{\ell-1}(M-\ell+n-k-1)!n!(M-1)!}{(M-\ell)!(M+n-1)!(n-k-1)!} \\
& =\frac{\binom{\ell+k-1}{\ell-1}\binom{M-\ell+n-k-1}{M-\ell}}{\binom{M+n-1}{n}} \\
& =\frac{P(\ell-1, k+1) P(M-\ell, n-k)}{P(M-1, n+1)}
\end{aligned}
$$

which is the probability that when putting $M-1$ balls uniformly and independently randomly into $n+1$ baskets, there are $\ell-1$ balls inside $k+1$ specific baskets and the remaining $n-\ell$ balls inside the remaining $n-k$ baskets. Summing over $\ell$ from 1 to $M$, this covers all possible number of balls ending into those $k+1$ specific baskets (i.e. 0 to $M-1$ ), so this sums up to 1 over $1 \leq \ell \leq M$ and we are done.

Note that we can use the above result to count the number of closed walks in an acyclic regular graph:

Lemma 4.6. Suppose $G$ is $k$-regular. Let $v_{0} \in V(G)$ and suppose the subgraph of $G$ induced by vertices at distance at most $r / 2$ from $v_{0}$ is acyclic (i.e. a tree rooted at $v_{0}$ ). Then the number of closed walks of length $r$ in $X$ starting at $v_{0}$ is $\theta(r)$, where $\theta(r)=0$ if $r$ is odd, $\theta(0)=1$ and otherwise

$$
\begin{aligned}
\theta(2 s) & =\sum_{i=1}^{s}\binom{2 s-i}{s} \frac{i}{2 s-i} k^{i}(k-1)^{s-i} \\
& =k \sum_{i=0}^{s-1}\binom{2 s}{i} \frac{s-i}{s}(k-1)^{i} \\
& =\sum_{i=1}^{s}\binom{2 s}{i} \frac{2 s-2 i+1}{2 s-i+1}(k-1)^{i} .
\end{aligned}
$$

Proof. Let $\mathbf{v}=v_{0}, \ldots, v_{r}$ be a closed walk of length $r$. Corresponding to $\mathbf{v}$ we have a sequence of nonnegative integers $\mathbf{d}=d_{0}, \ldots, d_{r}$, where $d_{i}$ is the distance from $v_{0}$ to $v_{i}$ in $G$. Then $\left|d_{i}-d_{i-1}\right|=1$ for $r \geq i \geq 1$ and $d_{r}-d_{0}=\sum_{i=1}^{r}\left(d_{i}-d_{i-1}\right)=\left|\left\{i: d_{i}-d_{i-1}=1\right\}\right|-\left|\left\{i: d_{i}-d_{i-1}=-1\right\}\right|=0$ so $r$ is even so let $s=r / 2$ where $s>0$.

It is proven in 4.5 (to see that it applies here, note $d_{j}=\sum_{i=1}^{j} x_{i}$ ) that the number of possible $\mathbf{d}$ with $i$ $d_{j}$ (for $\left.1 \leq j \leq r\right)$ being $0(i \geq 0)$ (we say $\mathbf{d}$ has $i$ zeros) is

$$
\binom{2 s-i}{s} \frac{i}{2 s-i}
$$

Fix $\mathbf{d}$ with $i$ zeros. Note that this fixes a sequence of moving away or towards $v_{0}$, starting at $v_{0}$ for the walk. Every move towards $v_{0}$ is predetermined. Since we start at $v_{0} i$ times throughout the walk and each time has $k$ possibilities so those account for $k^{i}$ possibilities in total. All the $s-i$ other times we need to move away from $v_{0}$ so those account for $(k-1)^{s-i}$ possibilities, thus each $\mathbf{d}$ accounts for $k^{i}(k-1)^{s-i}$ closed walks, giving us the first expression. The other equalities are by algebra.

We then determine that $\theta(r)$ holds asymptotically without the acyclic assumption:

Lemma 4.7. For $r \geq 0, i \geq 1$ let $\phi_{r}\left(G_{i}\right)$ denote the total number of closed walks of length $r$ in $G_{i}$. Then for each $r, \phi_{r}\left(G_{i}\right) / n\left(G_{i}\right) \rightarrow \theta(r)$ as $i \rightarrow \infty$.

Proof. Let $n_{r}\left(G_{i}\right)$ denote the number of vertices in $G_{i}$ that satisfy the assumptions of 4.6 regarding $v_{0}$. Given $v_{0} \in V\left(G_{i}\right)$, call the subgraph used in $4.6 H_{r}\left(v_{0}\right) \subseteq G_{i}$. Then note $n_{r}\left(G_{i}\right) / n\left(G_{i}\right) \rightarrow 1$ as $i \rightarrow \infty$, since otherwise $n\left(G_{i}\right)-n_{r}\left(G_{i}\right)$ would not be $o\left(n\left(G_{i}\right)\right)$, and since the number of $v$ such that $V\left(H_{r}(v)\right) \cap V\left(H_{r}\left(v_{0}\right)\right) \neq$ $\varnothing$ is upper bounded ${ }^{3}$ regardless of $G_{i}$, the number of cycles that are in some $H_{r}(v)$, i.e. of size at most $k(k-1)^{r / 2-1}+1$ (upper bound of the size of $\left.H_{r}(v)\right)$, is not $o\left(n\left(G_{i}\right)\right.$ ), so there must be some $r_{0}$ at most $k(k-1)^{r / 2-1}+1$ such that $c_{r_{0}}\left(G_{i}\right)$ is not $o\left(n\left(G_{i}\right)\right)$, a contradiction to (4.1).

Note that for any $v \in V(G)$, the number of closed walks of length $r$ is trivially upper bounded by $k^{r-1}$, thus by squeeze theorem

$$
\frac{n_{r}\left(G_{i}\right) \theta(r)}{n\left(G_{i}\right)} \leq \frac{\phi_{r}\left(G_{i}\right)}{n\left(G_{i}\right)} \leq \frac{n_{r}\left(G_{i}\right) \theta(r)+\left(n\left(G_{i}\right)-n_{r}\left(G_{i}\right)\right) k^{r-1}}{n\left(G_{i}\right)} \Longrightarrow \frac{\phi_{r}\left(G_{i}\right)}{n\left(G_{i}\right)} \stackrel{i \rightarrow \infty}{\longrightarrow} \theta(r) .
$$

We can rewrite $\phi_{r} / n$ as a Lebesgue-Stieltjes integral with respect to $F_{G_{i}}(x)$ :

Lemma 4.8. For each $r \geq 0$,

$$
\int x^{r} d F_{G_{i}}(x) \xrightarrow{i \rightarrow \infty} \theta(r)
$$

Proof. By definition of $F_{G}$ and 1.8,

$$
\int x^{r} d F_{G_{i}}(x)=\sum_{\lambda \in \operatorname{Spec} G_{i}} \lambda^{r} p(\lambda)=\frac{1}{n\left(G_{i}\right)} \sum_{j=1}^{n\left(G_{i}\right)} \lambda_{j}^{r}=\frac{1}{n\left(G_{i}\right)} \operatorname{Tr}\left(A^{r}\right)=\frac{\phi_{r}\left(G_{i}\right)}{n\left(G_{i}\right)} \forall r \geq 0
$$

Now we arrive at the existence and uniqueness of $F(x)$.

Theorem 4.9. There is a unique $F(x)$ increasing and right-continuous such that

$$
\int x^{r} d F=\theta(r) \forall r \geq 0
$$

Furthermore, $F_{G_{i}}(x) \rightarrow F(x)$ as $i \rightarrow \infty$ for every $x$ where $F$ is continuous.

To prove this, we state 3 standard results in analysis. Let $I=[\alpha, \beta]$ and for each $M \geq 0$, define $\operatorname{RBV}(I, M)$ to be the set of all real $f(x)$ such that

1. $f(x)=0$ if $x<\alpha$ and $f(x)$ is constant if $x>\beta$;
2. $f$ is right-continuous;
3. the total variation of $f$ is at most $M$.
[^2]Lemma 4.10. If $f \in \operatorname{RBV}(I, M)$ and $\int x^{r} d f=0 \forall r \geq 0$, then $f(x)=0$ a.e.

Lemma 4.11 (Helley-Bray). Let $f_{1}, f_{2}, \ldots$ be a sequence in $\operatorname{RBV}(I, M)$ such that $f_{n}(x) \rightarrow f(x)$ as $i \rightarrow \infty$ for some $f \in \operatorname{RBV}(I, M)$ at every $x$ where $f$ is continuous. Then $\int x^{r} d f_{i} \rightarrow \int x^{r} d f$ as $i \rightarrow \infty$ for each $r \geq 0$.

Lemma 4.12 (Helley selection). Let $f_{1}, f_{2}, \ldots$ be a sequence in $\operatorname{RBV}(I, M)$. Then there exists a subsequence $f_{n_{1}}, f_{n_{2}}, \ldots$ and $f \in \operatorname{RBV}(I, M)$ such that $f_{n_{k}}(x) \rightarrow f(x)$ as $k \rightarrow \infty$ at every $x$ where $f$ is continuous.

We use the above to prove the following:

Theorem 4.13. Let $f_{1}, f_{2}, \ldots$ be a sequence in $\operatorname{RBV}(I, M)$ such that $\int x^{r} d f_{i} \rightarrow \mu_{r}$ as $i \rightarrow \infty$ for each $r \geq 0$, where $\mu_{r} \in \mathbb{R} \forall r \geq 0$. Then there exists a unique (up to being equal a.e.) $f \in \operatorname{RBV}(I, M)$ such that $\int x^{r} d f=\mu_{r} \forall r \geq 0$. Furthermore $f_{n}(x) \rightarrow f(x)$ wherever $f(x)$ is continuous.

Proof. Suppose $f, g \in \operatorname{RBV}(I, M)$ satisfy $\int x^{r} d f=\int x^{r} d g=\mu_{r} \forall r \geq 0$, then $f-g \in \operatorname{RBV}(I, 2 M)$ satisfies $\int x^{r} d(f-g)=\int x^{r} d f-\int x^{r} d g=0 \forall r \geq 0$ so by $4.10 f-g \equiv 0$ a.e.

We know by 4.12 that any subsequence of $f_{n}$ has itself a subsequence $\mathbf{x}$ that converge to some $f_{\mathbf{x}} \in$ $\operatorname{RBV}(I, M)$ wherever $f_{\mathbf{x}}$ is continuous. By 4.11, all such $f_{\mathbf{x}}$ satisfy $\int x^{r} d f_{\mathbf{x}}=\mu_{r}$ so there is a unique $f_{\mathbf{x}} \equiv f \forall \mathbf{x}$.

It suffices now to prove that $f_{n}(x) \rightarrow f(x)$ wherever $f(x)$ is continuous. Suppose there exists $x_{0}$ where $f$ is continuous, but $f_{n}\left(x_{0}\right)$ does not converge to $f\left(x_{0}\right)$. Then there must exist $\epsilon$ and a subsequence $f_{n_{k}}$ such that

$$
\left|f_{n}\left(x_{0}\right)-f_{n_{k}}\left(x_{0}\right)\right| \geq \epsilon \forall k \geq 1
$$

By $4.12 f_{n_{k}}$ has a subsequence $\mathbf{x}=f_{n_{k_{j}}}$ where $f_{n_{k}}(x) \rightarrow f_{\mathbf{x}}(x)=f(x)$ wherever $f(x)$ is continuous. However $f_{n_{k_{j}}}\left(x_{0}\right)$ cannot converge to $f\left(x_{0}\right)$ so $f$ is not continuous at $x_{0}$, a contradiction.

Then our result 4.9 is a direct corollary:
Proof of 4.9. 4.13 applies with $M=1$ (since $F_{G_{i}}$ are CDFs) and $I=[-k, k]$. It suffices to show that we have a unique, increasing and right-continuous $F$.

Let $X \subseteq \operatorname{RBV}(I, M)$ be the set of all $f$ obtainable from 4.13 with $F_{G_{i}}$. For any $f \in X$, since $f \in$ $\operatorname{RBV}(I, M), f$ is continuous a.e. Call the set on which $f$ is continuous $A_{f} \subseteq \mathbb{R}$ and let $A=\bigcup_{f \in X} A_{f}$. Then $F_{G_{i}}(x)$ converges in $A$. Since $F_{G_{i}}$ are increasing, $\left.f\right|_{A}$. Let $F$ be such that $F(x)=\lim _{i \rightarrow \infty} F_{G_{i}}(x) \forall x \in A$ and for every $x_{0} \notin A, F\left(x_{0}\right)=\inf _{x \in A, x \geq x_{0}} F(x)$. Then $F$ since $F$ agrees with any $f \in X$ on $A_{f}, F$ agrees with $f$ a.e. so $F \in X$. $F$ is unique and increasing, so the only discontinuity is of the first kind, so by its definition $F$ is right-continuous.

### 4.3 Derivations of $F(x)$

This section deals with the derivation of $F(x)$ from 4.9 and 4.6 and is quite technical so only an outline will be given. We start with an asymptotic expression for $\theta(2 s)$ :

Lemma 4.14. As $s \rightarrow \infty$,

$$
\theta(2 s) \sim \frac{4^{s} k(k-1)^{s+1}}{s(k-2)^{2} \sqrt{\pi s}}
$$

Outline of proof. By taking $z=1 /(k-1)$, the second form for $\theta(2 s)$ from 4.6 can be written as $\left(\right.$ let $\left.i^{\prime}=s-i\right)$

$$
\begin{aligned}
\theta(2 s) & =k \sum_{i=0}^{s-1}\binom{2 s}{i} \frac{s-i}{s}(k-1)^{i} \\
& =\frac{k(k-1)^{s-1}}{s} \sum_{i=0}^{s-1}\binom{2 s}{i}(s-i)(k-1)^{i-s+1} \\
& =\frac{k(k-1)^{s-1}}{s} \sum_{i^{\prime}=1}^{s}\binom{2 s}{s-i^{\prime}} i^{\prime} z^{i^{\prime}-1} \\
& =\frac{k(k-1)^{s-1}}{s}\binom{2 s}{s-1}\left(\sum_{i^{\prime}=1}^{s} i^{\prime} z^{i^{\prime}-1}-\varepsilon(s)\right)
\end{aligned}
$$

where

$$
\varepsilon(s)=\sum_{i=1}^{s}\left(1-\frac{\binom{2 s}{s-i}}{\binom{2 s}{s-1}}\right) i z^{i-1}
$$

It can be shown that $\varepsilon(s) \rightarrow 0$ as $s \rightarrow \infty$, and it can also be shown that

$$
\theta(2 s) \sim \frac{k(k-1)^{s-1}}{s}\binom{2 s}{s-1} \sum_{i^{\prime}=1}^{s} i^{\prime} z^{i^{\prime}-1} \sim \frac{4^{s} k(k-1)^{s+1}}{s(k-2)^{2} \sqrt{\pi s}}
$$

Lemma 4.15. Define $\omega=\sup \{|x| \mid 0<F(x)<1\}$. Then $\omega=2 \sqrt{k-1}$.

Outline of proof. For any $s$, note that $x^{2} \leq \omega^{2} \Longrightarrow x^{2 s+2} \leq \omega^{2} x^{2 s}$ wherever $\frac{d}{d x} F(x)>0$ by definition of $\omega$ so $\int x^{2 s+2} d F \leq \omega^{2} \int x^{2 s} d F$. Thus

$$
\limsup _{s \rightarrow \infty} \frac{\int x^{2 s+2} d F}{\int x^{2 s} d F} \leq \omega^{2}
$$

For any $0<\beta<\omega$, we can find (details omitted)

$$
\liminf _{s \rightarrow \infty} \frac{\int x^{2 s+2} d F}{\int x^{2 s} d F} \geq \beta^{2}
$$

Thus letting $\beta \rightarrow \omega$, we get $\lim _{s \rightarrow \infty} \frac{\int x^{2 s+2} d F}{\int x^{2 s} d F}=\lim _{s \rightarrow \infty} \frac{\theta(2(s+1))}{\theta(2 s)} \rightarrow \omega^{2}$ as $s \rightarrow \infty$. Thus by $4.14 \omega^{2}=$ $4(k-1) \Longrightarrow \omega=2 \sqrt{k-1}$.

Lemma 4.16. $F(x)$ is continuous at $x= \pm \omega$.

The proof is based on examining the jump at $x= \pm \omega$ and is omitted.
Finally some further analysis involving Tchebyshev polynomials was used to prove 4.1.

### 4.4 Application to random regular graphs

Let $k \geq 2$ and $n_{1}<n_{2}<\ldots$ be the sequence of number of vertices on which $k$-regular graphs are possible. ${ }^{4}$ For each $i$ define $R_{i}$ to be the set of all labeled $k$-regular graphs on $n_{i}$ vertices. Define $F_{i}(x)$ to be the average of $F_{G}(x)$ with $G$ taken over all $G \in R_{i}$, i.e. $F_{i}(x)=\frac{1}{n_{i}} \sum_{G \in R_{i}} F_{G}(x) . F_{i}$ can be seen as the expected eigenvalue distribution of $k$-regular graphs on $n_{i}$ vertices. We will state the following lemma without proof:

Lemma 4.17. For each $r \geq 3$ define $c_{r, i}$ to be the average number of $r$-cycles in one member of $R_{i}$. Then $c_{r, i} \rightarrow(k-1)^{r} / 2 r$ as $i \rightarrow \infty$.

For our purposes, we actually only need a weaker result, which will be proven:

## Lemma 4.18.

$$
c_{r, i} \xrightarrow{i \rightarrow \infty} M_{r}<\infty .
$$

Proof. Uniformly randomly take $G \in R_{i}$ and $v \in V(G)$. For the subgraph $H$ induced by the set of all vertices with distance at most $m=\lceil r / 2\rceil$ from $v$, if $H$ is acyclic then $v$ is not part of any $r$-cycle. There are at most $\binom{n_{i}-1}{k}\binom{n_{i}-1}{k-1} ~ k(k-1)^{m-1}$ possibilities for $H$ (the set of all $H$ is the subset of all trees where each child of $u$ is sampled from $V(G) \backslash\{u\}$ : there are $\binom{n_{i}-1}{k}$ possibilities for children of $v$ and $\binom{n_{i}-1}{k-1}$ possibilities for the children of $u \neq v$ ). Out of these, there are at least $\binom{n_{i}-1}{k}\binom{n_{i}-1-k(k-1)^{m-1}}{k-1}^{k(k-1)^{m-1}}$ possibilities where $H$ is acyclic (since there are never less than $n_{i}-1-k(k-1)^{m-1}$ options to chose from for each node).

Let $t=k(k-1)^{m-1}$, then the probability that $H$ is acyclic is at least $\prod_{j=0}^{k-2} \frac{n_{i}-1-t-i}{n_{i}-1-i} \geq\left(\frac{n_{i}-t-k+1}{n_{i}-k+1}\right)^{k-1}$. Thus the expected number of $v \in V(G)$ whose $H$ is acyclic is at least $n_{i}\left(\frac{n_{i}-t-k+1}{n_{i}-k+1}\right)^{k-1}$, so the expected number of $r$-cycles is at most ${ }^{5}$

$$
\frac{n_{i}}{r}\left(1-\left(\frac{n_{i}-t-k+1}{n_{i}-k+1}\right)^{k-1}\right)=\frac{n_{i}}{r}\left(1-\left(1-\frac{t}{n_{i}-k+1}\right)^{k-1}\right) \leq \frac{(k-1)^{2} n_{i} t}{r\left(n_{i}-k+1\right)}=O(1)
$$

## Theorem 4.19.

$$
F_{i}(x) \xrightarrow{i \rightarrow \infty} F(x)
$$

[^3]Proof. Define graph $Y_{i}$ to be the graph consisting of all graphs in $R_{i}$ put together in disjoint union. Then $F_{Y_{i}} \equiv F_{i}$. Now it suffices to show that $Y_{i}$ satisfies the assumptions in 4.1. Indeed since by 4.18 there are on average at most $M_{r}+\epsilon r$-cycles per member of $R_{i}$ (true for $i$ large enough),

$$
c_{r}\left(Y_{i}\right) \leq\left|R_{i}\right|\left(M_{r}+\epsilon\right)
$$

Since $n\left(Y_{i}\right)=\left|R_{i}\right| n_{i}$,

$$
\frac{c_{r}\left(Y_{i}\right)}{n\left(Y_{i}\right)} \leq \frac{M_{r}+\epsilon}{n_{i}} \xrightarrow{i \rightarrow \infty} 0
$$

## 5 Vertex replacement on regular graphs

Let $\mathcal{F}_{k}$ be the set of $k$-regular graphs ( $k \geq 2$ across this section). For a given $k$ we will define a map from $\mathcal{F}_{k}$ onto itself that will help us generate $k$-regular graphs, via vertex replacement.

Definition 5.1. Let $G$ be a graph. $D(G)=D_{G}$ (called the subdivision graph of $G$ ) is defined to be the graph created from subdividing every edge in $E(G)$. Formally, $V\left(D_{G}\right)=E(G) \cup V(G)$ and $E\left(D_{G}\right)=\{(e, v)$ : $e \in E(G), v \in V(G), e$ incident to $v\}$.

Remark. For $x \in V\left(D_{G}\right)$, if $x \in V(G)$ then $\operatorname{deg}_{D_{G}}(x)=\operatorname{deg}_{G}(x)$ and if $x \in E(G)$ then $\operatorname{deg}_{G}(x)=2$.

We start with a representation of the adjacency matrix of $D_{G}$.

Lemma 5.2. Let $G$ have $n$ vertices and $\ell$ edges. Let $\mathbf{X}_{n \times \ell}=\mathbf{X}(G)$ be the undirected incidence matrix of $G$ $\left(\mathbf{X}_{v e}=1\right.$ if and only if $e$ incident to $\left.v\right)$. Then

$$
\mathbf{A}\left(D_{G}\right)=\left(\begin{array}{cc}
0_{\ell \times \ell} & \mathbf{X}(G)^{\boldsymbol{\top}} \\
\mathbf{X}(G) & 0_{n \times n}
\end{array}\right)
$$

Proof. The first $\ell$ rows and columns of $\mathbf{A}\left(D_{G}\right)$ represent $E(G)$ and the following $n$ represent $V(G) . \mathbf{A}_{i j}=0$ if $i, j \in V(G)$ or $i, j \in E(G)$ and $\mathbf{A}_{i j}=1$ if and only if $\mathbf{X}(G)_{i j}=1$ (if $i \in V(G)$ and $\left.j \in E(G)\right)$ or $\mathbf{X}(G)_{j i}=1$ (if $i \in E(G)$ and $j \in V(G)$ ).

Definition 5.3. Let $G$ be a graph. $L(G)=L_{G}$ the line graph of $G$ is defined where $V\left(L_{G}\right)=E(G)$ and $(e, f) \in E\left(L_{G}\right)$ if and only if $e, f$ are incident to the same vertex in $G$.

Lemma 5.4. For each $v \in V(G)$, if it is incident to $e_{1}, \ldots, e_{d} \in E(G)$, then we can identify it with $H_{v} \subseteq L_{G}$, where $V\left(H_{v}\right)$ is the induced subgraph of $\left\{e_{1}, \ldots, e_{d}\right\}, H_{v}$ is isomorphic to $K_{d}$ and $H_{v} \cap H_{u} \neq \varnothing$ if and only if $(u, v) \in E(G)$, in which case $H_{v} \cap H_{u}=\{(u, v)\}$.

Proof. $H_{v}$ is isomorphic to $K_{d}$ since every $e, f \in V\left(H_{v}\right)$ are neighbours (they are all incident to $v$ ). The intersection property follows $V\left(H_{v} \cap H_{u}\right)=V\left(H_{v}\right) \cap V\left(H_{u}\right)$.

Lemma 5.5. For $e=(v, u) \in V\left(L_{G}\right)$ (where $v, u \in V(G)$ were neighbours in $\left.G\right)$, $\operatorname{deg}_{L_{G}}(e)=\operatorname{deg}_{G}(v)+$ $\operatorname{deg}_{G}(u)-2$.

Proof. $e$ is only neighbours with $f \in E(G)$ where $f \neq e$ incident to $v$ or to $u$, and there are $\operatorname{deg}_{G}(v)+$ $\operatorname{deg}_{G}(u)-2$ of those.

Proposition 5.6. Let $G \in \mathcal{F}_{k}$ be $k$-regular.

$$
\phi: \mathcal{F}_{k} \rightarrow \mathcal{F}_{k}, G \mapsto L(D(G))
$$

is a mapping where every $v \in V(G)$ can be identified to $H_{v} \subseteq \phi(G)$ whose vertex set are edges incident to $v$ in $G . H_{v}$ is isomorphic to $K_{k}$ and $H_{v}$ is joined by an edge to $H_{u}$ if and only if $(u, v) \in E(G)$.

Remark. $\phi$ indeed maps into $\mathcal{F}_{k}$ by 5.5, since every $e \in E\left(D_{G}\right)$ joins a vertex of degree $k$ (from $V(G)$ ) and one of degree $2\left(\right.$ from $\left.E(G) \subseteq V\left(D_{G}\right)\right)$. So $e \in V(\phi(G))$ would have degree $k+2-2=k$.

Proof. Since $V(G) \subseteq V(D(G))$ and the degrees stay unchanged, we use $H_{v}$ from 5.4 to get that they are isomorphic to $K_{k}$ and that they share a vertex with each of $H_{e}$, where $e \in E(G)$ was added by the subdivisions and $e$ is incident to $v$ and some $u$ neighbour of $v$. But this means $H_{e}$ also shares a vertex with $H_{u}$ (and $H_{u} \cap H_{v}=\varnothing$ since $u, v \in V(G)$ can never be neighbours in $D_{G}$ ), which means $H_{v}$ shares an edge with $H_{u}$.

By the degree of every $x \in V\left(H_{v}\right)$, since they have $k-1$ degrees inside $H_{v}$, they must have precisely one degree out of $H_{v}$, hence the only if.

Definition 5.7. $\phi(G)$ (from 5.6) is called the para-line graph of $G$.

In order to study the relationship between $\operatorname{Spec} G$ and $\operatorname{Spec} \phi(G)$ for $G \in \mathcal{F}_{k}$, we will use two well-known results from linear algebra, as well as one result that applies to graphs in general.

Lemma 5.8. If $\mathbf{M}$ is invertible and $\mathbf{T}$ is square,

$$
\operatorname{det}\left(\begin{array}{cc}
\mathbf{M} & \mathbf{X} \\
\mathbf{Y} & \mathbf{T}
\end{array}\right)=\operatorname{det} \mathbf{M} \operatorname{det}\left(\mathbf{T}-\mathbf{Y} \mathbf{M}^{-1} \mathbf{X}\right)
$$

Proof. Note (assume $\mathbf{M}_{\ell \times \ell}$ and $\mathbf{T}_{n \times n}$ )

$$
\left(\begin{array}{cc}
\mathbf{M}^{-1} & 0_{\ell \times n} \\
0_{n \times \ell} & \mathbf{I}_{n}
\end{array}\right)\left(\begin{array}{cc}
\mathbf{I}_{\ell} & -\mathbf{X} \\
0_{n \times \ell} & \mathbf{I}_{n}
\end{array}\right)=\left(\begin{array}{cc}
\mathbf{M}^{-1} & -\mathbf{M}^{-1} \mathbf{X} \\
0_{n \times \ell} & \mathbf{I}_{n}
\end{array}\right)
$$

is precisely the product of column operations that eliminate $\mathbf{M}$ and $\mathbf{X}$. Then

$$
\begin{aligned}
& \left(\begin{array}{cc}
\mathbf{M} & \mathbf{X} \\
\mathbf{Y} & \mathbf{T}
\end{array}\right)\left(\begin{array}{cc}
\mathbf{M}^{-1} & -\mathbf{M}^{-1} \mathbf{X} \\
0_{n \times \ell} & \mathbf{I}_{n}
\end{array}\right)=\left(\begin{array}{cc}
\mathbf{I}_{\ell} & 0_{\ell \times n} \\
\mathbf{Y} \mathbf{M}^{-1} & \mathbf{T}-\mathbf{Y M}^{-1} \mathbf{X}
\end{array}\right) \\
\Longrightarrow & \operatorname{det}\left(\begin{array}{cc}
\mathbf{M} & \mathbf{X} \\
\mathbf{Y} & \mathbf{T}
\end{array}\right)(\operatorname{det} \mathbf{M})^{-1}=\operatorname{det}\left(\mathbf{T}-\mathbf{Y} \mathbf{M}^{-1} \mathbf{X}\right) .
\end{aligned}
$$

Lemma 5.9. If $\mathbf{X}_{m \times n}, \mathbf{Y}_{n \times m}$ are matrices and $\chi_{\mathbf{A}}$ denotes the characteristic polynomial of a square matrix A, then

$$
\lambda^{n} \chi_{\mathbf{X} \mathbf{Y}}(\lambda)=\lambda^{m} \chi_{\mathbf{Y} \mathbf{X}}(\lambda)
$$

Proof. We start by the case where $\mathbf{X}, \mathbf{Y}$ are $n \times n$, and $\mathbf{X}$ is invertible. Then

$$
\mathbf{X}(\mathbf{Y X}) \mathbf{X}^{-1}=\mathbf{X Y} \Longrightarrow \operatorname{det}\left(\lambda \mathbf{I}_{n}-\mathbf{X Y}\right)=\operatorname{det}\left(\mathbf{X}\left(\lambda \mathbf{I}_{n}-\mathbf{Y X}\right) \mathbf{X}^{-1}\right)=\operatorname{det}\left(\lambda \mathbf{I}_{n}-\mathbf{Y X}\right)
$$

For $\mathbf{X}, \mathbf{Y}$ both singular, for fixed $\lambda \in \mathbb{R}$, if we consider

$$
\psi_{\lambda}: \mathbb{R}^{2 n^{2}} \rightarrow \mathbb{R},(\mathbf{X}, \mathbf{Y}) \mapsto \chi_{\mathbf{X Y}}(\lambda) ; \phi_{\lambda}: \mathbb{R}^{2 n^{2}} \rightarrow \mathbb{R},(\mathbf{X}, \mathbf{Y}) \mapsto \chi_{\mathbf{Y X}}(\lambda)
$$

then $\psi_{\lambda}=\phi_{\lambda}$ almost everywhere (wherever one of $\mathbf{X}, \mathbf{Y}$ is invertible), so since $\psi_{\lambda}$ and $\phi_{\lambda}$ are analytic, $\psi_{\lambda}=\phi_{\lambda}$ everywhere. Thus $\chi_{\mathbf{X Y}}=\chi_{\mathbf{Y X}}$ whenever $\mathbf{X}, \mathbf{Y}$ are square.

Given $n>m$, we write (note $\mathbf{X}_{n \times n}^{\prime}$ and $\mathbf{Y}_{n \times n}^{\prime}$ )

$$
\begin{aligned}
& \mathbf{X}^{\prime} \mathbf{Y}^{\prime}=\binom{\mathbf{X}}{0_{(n-m) \times n}}\left(\begin{array}{ll}
\mathbf{Y} & 0_{n \times(n-m)}
\end{array}\right)=\left(\begin{array}{cc}
\mathbf{X Y} & 0_{m \times(n-m)} \\
0_{(n-m) \times m} & 0_{(n-m) \times(n-m)}
\end{array}\right) \\
& \Longrightarrow \chi_{\mathbf{X}^{\prime} \mathbf{Y}^{\prime}}(\lambda)=\operatorname{det}\left(\begin{array}{cc}
\lambda \mathbf{I}_{m}-\mathbf{X Y Y} & 0_{m \times(n-m)} \\
0_{(n-m) \times m} & \lambda \mathbf{I}_{n-m}
\end{array}\right)=\lambda^{n-m} \chi_{\mathbf{X Y}}(\lambda)=\chi_{\mathbf{Y}^{\prime} \mathbf{X}^{\prime}}(\lambda) \text {. }
\end{aligned}
$$

and

$$
\mathbf{Y}^{\prime} \mathbf{X}^{\prime}=\left(\begin{array}{ll}
\mathbf{Y} & 0_{n \times(n-m)}
\end{array}\right)\binom{\mathbf{X}}{0_{(n-m) \times n}}=\mathbf{Y X} \Longrightarrow \lambda^{n-m} \chi_{\mathbf{X Y}}(\lambda)=\chi_{\mathbf{Y}^{\prime} \mathbf{X}^{\prime}}(\lambda)=\chi_{\mathbf{Y} \mathbf{X}}(\lambda)
$$

Lemma 5.10. Let $G$ be a graph with $n$ vertices, $\ell$ edges, degree matrix $\mathbf{D}$ and adjacency matrix $\mathbf{A}$. Then

$$
\chi_{L_{G}}(\lambda-2)=\lambda^{\ell-n} \operatorname{det}\left(\lambda \mathbf{I}_{n}-\mathbf{A}-\mathbf{D}\right)
$$

Proof. Note $\mathbf{A}\left(L_{G}\right)+2 \mathbf{I}_{\ell}=\mathbf{X}(G)^{\top} \mathbf{X}(G)$. Indeed for $i \neq j, \mathbf{A}\left(L_{G}\right)_{i j}=1$ if and only if $i$ and $j$ are incident to the same vertex, which is precisely when the $i$-th and $j$-th columns of $\mathbf{X}(G)$ have a 1 at the same position. For $i=j$, the $i$-th column of $\mathbf{X}(G)$ always has precisely two 1 s , so $\left(\mathbf{X}(G)^{\top} \mathbf{X}(G)\right)_{i i}=2$. This implies

$$
\chi_{L_{G}}(\lambda-2)=\operatorname{det}\left(\lambda \mathbf{I}_{\ell}-2 \mathbf{I}_{\ell}-\mathbf{A}\left(L_{G}\right)\right)=\chi \mathbf{x}(G)^{\top} \mathbf{X}(G)(\lambda)
$$

On the other hand $\mathbf{X}(G) \mathbf{X}(G)^{\boldsymbol{\top}}=\mathbf{A}+\mathbf{D}$. Indeed for $i \neq j, \mathbf{A}_{i j}=1$ if and only if $i$ and $j$ are incident to the same edge, which is when $i$ is adjacent to $j$. If $i=j$ then $(\mathbf{X}(G) \mathbf{X}(G))_{i i}=\operatorname{deg}(i)=\mathbf{D}_{i i}$. The result then follows from 5.9:

$$
\chi_{L_{G}}(\lambda-2)=\chi_{\mathbf{X}(G)^{\top} \mathbf{X}(G)}(\lambda)=\lambda^{\ell-n} \chi_{\mathbf{x}(G) \mathbf{X}(G)^{\top}}(\lambda)=\lambda^{\ell-n} \chi_{\mathbf{A}+\mathbf{D}}(\lambda)=\lambda^{\ell-n} \operatorname{det}\left(\lambda \mathbf{I}_{n}-\mathbf{A}-\mathbf{D}\right)
$$

Now we are ready to state a result due to Cvetković to study the eigenvalues of $\phi(G)$.

Theorem 5.11. Let $G \in \mathcal{F}_{k}$ have $n$ vertices and $\ell=n k / 2$ edges. Then

$$
\chi_{\phi(G)}(\lambda)=(\lambda(\lambda+2))^{\ell-n} \chi_{G}\left(\lambda^{2}+(2-k) \lambda-k\right)
$$

Proof. Recall

$$
\mathbf{A}\left(D_{G}\right)=\left(\begin{array}{cc}
0_{\ell \times \ell} & \mathbf{X}(G)^{\boldsymbol{\top}} \\
\mathbf{X}(G) & 0_{n \times n}
\end{array}\right), \mathbf{D}\left(D_{G}\right)=\left(\begin{array}{cc}
2 \mathbf{I}_{\ell} & 0_{\ell \times n} \\
0_{n \times \ell} & k \mathbf{I}_{n}
\end{array}\right)
$$

By 5.10

$$
\chi_{\phi(G)}(\lambda)=(\lambda+2)^{\ell-n} \operatorname{det}\left((\lambda+2) \mathbf{I}_{\ell+n}-\mathbf{A}\left(D_{G}\right)-\mathbf{D}\left(D_{G}\right)\right)
$$

Then by 5.8

$$
\begin{aligned}
(\lambda+2) \mathbf{I}_{\ell+n}-\mathbf{A}\left(D_{G}\right)-\mathbf{D}\left(D_{G}\right) & =\left(\begin{array}{cc}
(\lambda+2-2) \mathbf{I}_{\ell} & -\mathbf{X}(G)^{\top} \\
-\mathbf{X}(G) & (\lambda+2-k) \mathbf{I}_{n}
\end{array}\right)=\left(\begin{array}{cc}
\lambda \mathbf{I}_{\ell} & -\mathbf{X}(G)^{\top} \\
-\mathbf{X}(G) & (\lambda-k+2) \mathbf{I}_{n}
\end{array}\right) \\
\Longrightarrow \operatorname{det}\left(\begin{array}{cc}
\lambda \mathbf{I}_{\ell} & -\mathbf{X}(G)^{\top} \\
-\mathbf{X}(G) & (\lambda-k+2) \mathbf{I}_{n}
\end{array}\right) & =\lambda^{\ell} \operatorname{det}\left((\lambda-k+2) \mathbf{I}_{n}-\lambda^{-1} \mathbf{X}(G) \mathbf{X}(G)^{\top}\right) \\
& =\lambda^{\ell-n} \operatorname{det}\left((\lambda(\lambda-k+2)-k) \mathbf{I}_{n}-\mathbf{A}(G)\right)=\lambda^{\ell-n} \chi_{G}(\lambda(\lambda-k+2)-k) \\
\Longrightarrow \chi_{\phi(G)}(\lambda) & =(\lambda(\lambda+2))^{\ell-n} \chi_{G}(\lambda(\lambda-k+2)-k)
\end{aligned}
$$

Corollary 5.12. Let $G \in \mathcal{F}_{k}$ have $n$ vertices and $\ell=n k / 2$ edges and let $f(\lambda)=\lambda^{2}+(2-k) \lambda-k$. Then

$$
\operatorname{Spec} \phi(G)=\{0\}^{(\ell-n)} \cup\{-2\}^{(\ell-n)} \cup f^{-1}(\operatorname{Spec} G)
$$

where $\operatorname{Spec} G$ and $\operatorname{Spec} \phi(G)$ are understood to be multisets (and $\{x\}^{(\ell-n)}$ denotes the multiset with $x$ having multiplicity $\ell-n$ ).

Proof. Everything follows 5.11: if $\lambda=0$ or -2 then $\chi_{\phi(G)}(\lambda)=0$; If $\lambda \in f^{-1}(\operatorname{Spec} G)$ then $f(\lambda) \in \operatorname{Spec} G$ so $\chi_{G}(f(\lambda))=0$ so $\chi_{\phi(G)}(\lambda)=0$. The multiplicities follow from the multiplicity of $\lambda$ as root of $\chi_{\phi(G)}$.

Now we use the notation

$$
f_{k}(\lambda)=\lambda^{2}+(2-k) \lambda-k
$$

Denote

$$
\Gamma_{k}=\bigcap_{j=0}^{\infty} f_{k}^{-j}([-k, k])
$$

We are looking for a set $A_{k} \subseteq[-k, k]$ such that

$$
\begin{equation*}
\operatorname{Spec} G \subseteq A_{k} \Longrightarrow \operatorname{Spec} \phi(G)=f_{k}^{-1}(\operatorname{Spec} G) \cup\{0,2\} \subseteq A_{k} \tag{5.1}
\end{equation*}
$$

Note that if $f_{k}^{-j}\left(A_{k} \cup\{0,-2\}\right) \subseteq A_{k}, \forall j \geq 0$ then $A_{k}$ would satisfy (5.1). We will proceed to show that $A_{k}=\Gamma_{k} \cup \bigcup_{j=0}^{\infty} f_{k}^{-j}(\{0\})$ has the desired property.

Lemma 5.13. For $x \notin[-k, k], f_{k}(x) \notin[-k, k]$. In particular, $x \notin f_{k}^{-1}[-k, k] \supseteq \Gamma_{k}$ and $f_{k}^{j}(x) \neq 0, \forall j \geq 0$.

Proof. Assume $x>k$ i.e. $x=k+\epsilon$. Then

$$
f_{k}(x)=(k+\epsilon)^{2}+(2-k)(k+\epsilon)-k=(k+\epsilon)(\epsilon+2)-k>(k+\epsilon)(\epsilon+2)-k-\epsilon=(k+\epsilon)(\epsilon+1)>k .
$$

Now assume $x<-k$ i.e. $x=-k-\epsilon$. Then (recall $k \geq 2$ so $2 k-3 \geq 1$ )

$$
f_{k}(x)=(-k-\epsilon)^{2}+(k-2)(k+\epsilon)-k>(k+\epsilon)(k+\epsilon+k-2)-k-\epsilon=(k+\epsilon)(2 k+\epsilon-3)>k .
$$

Theorem 5.14. Let $A_{k}=\Gamma_{k} \cup \bigcup_{j=0}^{\infty} f_{k}^{-j}(\{0\})$. Then $f_{k}^{-j}\left(A_{k} \cup\{0,-2\}\right) \subseteq A_{k} \subseteq[-k, k], \forall j \geq 0$. In other words, $A_{k}$ satisfies (5.1)

Proof. By 5.13, $\Gamma_{k} \subseteq[-k, k]$ and $\bigcup_{j=0}^{\infty} f_{k}^{-j}(\{0\}) \subseteq[-k, k]$ so $A_{k} \subseteq[-k, k]$. Note

$$
f_{k}^{-j}\left(A_{k} \cup\{0,-2\}\right)=f_{k}^{-j}\left(\Gamma_{k} \cup\{-2\}\right) \cup \bigcup_{i=0}^{\infty} f_{k}^{-(i+j)}(\{0\}), \forall j \geq 0
$$

Clearly $\bigcup_{i=0}^{\infty} f_{k}^{-(i+j)}(\{0\}) \subseteq A_{k}$ by definition. Furthermore, $f_{k}(-2)=(-2)^{2}+(2-k)(-2)-k=k$ and $f_{k}^{2}(-2)=f_{k}(k)=k^{2}+(2-k) k-k=k$ so $f_{k}^{j}(-2) \in[-k, k]$ for all $j \geq 0$ so $-2 \in f_{k}^{-j}([-k, k]), \forall j \in \mathbb{N}$ i.e. $-2 \in \Gamma_{k}$ so $\Gamma_{k} \cup\{-2\}=\Gamma_{k}$. Thus it would be sufficient if we can show that $f_{k}^{-j}\left(\Gamma_{k}\right) \subseteq \Gamma_{k}, \forall j \geq 0$.

Via induction this is equivalent to saying that $f_{k}^{-1}\left(\Gamma_{k}\right) \subseteq \Gamma_{k}$. If $x \in f_{k}^{-1}\left(\Gamma_{k}\right)$ i.e. $f_{k}(x) \in \Gamma_{k}$ so $f_{k}(x) \in$ $f_{k}^{-j}([-k, k]) \Longrightarrow x \in f_{k}^{-(j+1)}([-k, k]), \forall j \geq 0$. Now we only need $x \in[-k, k]$, but this follows 5.13: $x \notin[-k, k] \Longrightarrow f_{k}(x) \notin[-k, k] \supseteq \Gamma_{k}$.

What we showed is that suppose we are given $G \in \mathcal{F}_{k}$ such that $\operatorname{Spec} G \subseteq A_{k}$, define $\left\{G_{j}=\phi^{j}(G): j \geq 0\right\}$. Then $\operatorname{Spec} G_{j} \subseteq A_{k}, \forall j \geq 0$.

Example 5.15. Recall that

$$
\operatorname{Spec} K_{k+1}=\left(\begin{array}{cc}
k & -1 \\
1 & k
\end{array}\right)
$$

Note $f_{k}(-1)=1-2+k-k=-1$ is a fixed point so $-1 \in \Gamma_{k} \subseteq A_{k}$. So let $G=K_{k+1}$. Then $\left\{G_{j}=\phi^{j}(G)\right.$ : $j \geq 0\}$ has $\operatorname{Spec} G_{j} \subseteq A_{k}, \forall j \geq 0$.

### 5.1 Spectral gaps

One major application of studying spectra of graphs is the study of mixing times of random walks.

Definition 5.16. Let $G \in \mathcal{F}_{k}$. Then $\mathbf{M}_{G}=\frac{1}{k} \mathbf{A}_{G}$ is the random walk matrix on $G$.

The idea of random walk matrices is that $\mathbf{M}_{G}(i, j)$ is the probability of walking to node $j$ in 1 step, if one starts at node $i$ (namely $\frac{1}{k}$ for $G \in \mathcal{F}_{k}$ ).

Proposition 5.17. Let $\mathbf{M}$ be a random walk matrix on $G$. Let $p_{i, j}^{(m)}$ denote the probability of going to node $j$ in $m \in \mathbb{N}$ steps, starting at node $i$. Then $p_{i, j}^{(m)}=\left(\mathbf{M}^{m}\right)_{i, j}$.

Proof. We do induction on $m . m=1$ follows the definition.

## 6 Gap sets for the spectra of cubic graphs

This section is an outline of a paper by Kollár and Sarnak.
Throughout this section, we will only study cubic graphs, i.e. 3-regular connected, simple graphs (unless noted otherwise), and the set of all such graphs will be denoted $\mathbf{X}$ and its planar subset, $\mathbf{X}_{\text {Planar }}$. Recall that the spectrum of $G \in \mathbf{X}, \operatorname{Spec} G$, is a subset of $[-3,3]$ and $3 \in \operatorname{Spec} G$.

Definition 6.1. A closed $K \subseteq[-3,3]$ is called spectral if there exists infinitely many $\left\{G_{n}\right\}_{n \in \mathbb{N}} \subseteq \mathbf{X}$ such that $\operatorname{Spec} G_{n} \subseteq K \forall n \in \mathbb{N}$. $A=[-3,3] \backslash K$ is then called a gap set and each $x \in A$ is said to be gapped by $\left\{G_{n}\right\}$.

Note that trivially $\overline{\bigcup_{n \in \mathbb{N}} \operatorname{Spec} G_{n}}$ is spectral. Since any closed superset of a spectral set is spectral, we are interested in minimal spectral sets (or equivalently maximal gap sets) in particular, i.e. spectral sets for which no proper subset is spectral.

Gap sets are always open in $[-3,3]$ and 3 is never gapped. As such any $x \in[-3,3)$ that is gapped must also have a neighbourhood around it within $[-3,3)$ in which every point is gapped. The following is one of the main results of the paper:

Theorem 6.2. Every point in $[-3,3)$ is planar gapped.

The proof of 6.2 will, in particular, make use of the triangle map $\mathcal{T}$, defined as follows:

Definition 6.3. The triangle map $\mathcal{T}: \mathbf{X} \rightarrow \mathbf{X}$ is defined by replacing every vertex in $G$ by a triangle and joining them according to adjacency in $G$.

Explicitly, we first take a subdivision of $G$ by adding a vertex on every edge, then take the edge dual of that graph, such that every vertex that was in $V(G)$ becomes $K_{3}$ and every added vertex becomes $K_{2}$ joining the $K_{3}$ together.

It turns out that the spectrum for $G$ and $\mathcal{T}(G)$ are linked by the function

$$
f(x)=x^{2}-x-3
$$

In fact, it can be shown that $\operatorname{Spec} G$ relates to $\operatorname{Spec} \mathcal{T}(G)$ via the following:

$$
\begin{equation*}
\operatorname{Spec} \mathcal{T}(G)=f^{-1}(\operatorname{Spec} G) \cup\{0\} \cup\{-2\} \tag{6.1}
\end{equation*}
$$

with the multiplicities preserved in $f^{-1}(\operatorname{Spec} G)$ and with 0 and 2 having multiplicity $|V(G)| / 2$ (on top of the multiplicity of 0 and 2 from $\left.f^{-1}(\operatorname{Spec} G)\right)$. This makes the iterative properties of $f$ important. Consider

$$
\begin{equation*}
\Lambda=\bigcap_{m=0}^{\infty} f^{-m}([-3,3]) \tag{6.2}
\end{equation*}
$$

We will start by proving some general facts about sets invariant under some function. Let $g: \mathbb{R} \rightarrow \mathbb{R}$ and $C$ non-empty bounded subset of $\mathbb{R}$ be such that for $x \notin C, g^{n}(x) \notin C$ and is unbounded as $n \rightarrow \infty$. Let

$$
K=\bigcap_{m=0}^{\infty} g^{-m}(C)
$$

Proposition 6.4. $K$ is the unique largest bounded set closed under $g$ (meaning $g(K) \subseteq K$ ). Furthermore if $C \subseteq g(\mathbb{R})$ and is closed, then $K$ is non-empty.

Proof. Clearly no bounded $g$-closed set can contain any $x \notin C$. As a consequence, no bounded $g$-closed set can contain any $x$ where $g^{n}(x) \notin C$. Thus suppose $K^{\prime}$ is some bounded set where $g\left(K^{\prime}\right) \subseteq K^{\prime}$, then $K^{\prime} \subseteq K$. Now it suffices to prove that $K$ is indeed $g$-closed, i.e. $g(x) \in K$ for $x \in K$, but this is clear from the definition of $K$ (since $\left.g^{n}(g(x))=g^{n+1}(x) \in C\right)$. Since $g^{m}(x) \notin C \Longrightarrow g^{m+1} \notin C, g^{-m-1}(C) \subseteq g^{-m}(C)$ which means that if $C \subseteq g(\mathbb{R})$, then $g^{-m+1}(C) \subseteq g(\mathbb{R})$ so $g^{-m}(C)$ is non-empty for $m \in \mathbb{N}$ and compact, so $K \neq \varnothing$.

Corollary 6.5. $K$ is the unique largest bounded set invariant under $g$ (meaning $g(K)=K$ )

Proof. Since any $g$-invariant set is necessarily $g$-closed, by 6.4 it suffices to prove that $K$ itself is $g$-invariant. Note that we already have $g(K) \subseteq K \Longleftrightarrow K \subseteq g^{-1}(K)$. For any $x \in g^{-1}(K)$, note $g(x) \in K \subseteq g^{-1}(K)$, so $g^{-1}(K)$ is $g$-closed, so $g^{-1}(K) \subseteq K$ and we are done.

Now we take $g=f, C=[-3,3]$ and $K=\Lambda$. Then it suffices to show that $x \notin[-3,3]$ has $f(x) \notin[-3,3]$ and $f^{n}(x) \rightarrow \infty$. Note that for $\epsilon>0$, we have

$$
f(3+\epsilon)=(3+\epsilon)^{2}-6-\epsilon=3+5 \epsilon+\epsilon^{2}
$$

so taking $5 \epsilon+\epsilon^{2}$ as the new $\epsilon$, since $5\left(5^{n} \epsilon+\epsilon^{2}\right)>5^{n+1} \epsilon$ we inductively have $f^{n}(3+\epsilon)>3+5^{n} \epsilon \xrightarrow{n \rightarrow \infty} \infty$. For any $x<-3$, note $f(x)>x^{2}$ so $f^{n}(x)>x^{2 n} \rightarrow \infty$ (and obviously $f(x)>3$ ). Thus 6.4 and 6.5 apply to $f$ and $\Lambda$.

In fact, taking $C=[-2,3]$ would produce the same $\Lambda$, since

$$
f(-2-\epsilon)=(-2-\epsilon)^{2}-1+\epsilon=3+5 \epsilon+\epsilon^{2}
$$

so $f(-2-\epsilon)>3$ so by the above, we have $f^{n}(-2-\epsilon) \rightarrow \infty$ and everything applies. Now we show some properties of $\Lambda$.

Proposition 6.6. Let $C_{n}=f^{-n}([-2,3])$ be the $n$-th step in the construction of $\Lambda(n \geq 0)$. Then $C_{n}$ is made of $2^{n}$ disjoint (non-singleton) intervals. Furthermore, each endpoint $x$ in $C_{n}$ is also an endpoint in $C_{n+1}$.

Proof. Note that since the minimum of $f$ is $f(1 / 2)=-3-1 / 4$ which is smaller than all our intervals, each interval has two intervals as pre-images (and disjoint sets have disjoint pre-images) so $C_{n}$ consists of precisely $2^{n}$ intervals.

Any interval $[a, b]$ in $C_{n}, n \geq 1$ (meaning $[a, b] \subseteq C_{n}$ and $a-\epsilon, b+\epsilon \notin C_{n}$ for all $\epsilon$ small enough) must have $f([a, b])$ being an interval in $C_{n-1}$, since $f([a, b]) \subseteq C_{n-1}$ and $f(a) \pm \epsilon \in C_{n-1} \Longrightarrow a \pm \delta \in C_{n}$ for $\delta$ small enough (analogous for $f(b)$ ), and $f([a, b])$ is an actual interval since $f$ is monotonic on $[a, b]$.

Thus if $x$ were an endpoint of some interval in $C_{n}$, then $f^{m}(x)$ is an endpoint in $C_{n-m}$ so $x \in \Lambda$, since $f^{n}(x) \in\{-2,3\}$ so $f^{m}(x)=3$ for $m>n$. This means $x \in C_{n+1}$, which forces $x$ to be an endpoint, since $C_{n+1} \subseteq C_{n}$.

## Corollary 6.7. $\Lambda$ is uncountable.

Proof. 6.6 showed that, since $C_{n+1} \subseteq C_{n}$ and has twice as many intervals, with one of the endpoints of each interval matching with endpoints of intervals in $C_{n}$ and all the intervals in $C_{n+1}$ being subsets of those in $C_{n}$, each interval in $C_{n}$ must have been split into two intervals in $C_{n+1}$. The proof of uncountability then is exactly the same as the one for the Cantor set.

## 7 The chromatic polynomial

In this section, we deal with general graphs rather than simple graphs. By an $r$-color partition of $G$, we mean a partitioning of $V(G)$ into $\left\{V_{1}, \ldots, V_{r}\right\}$ such that no $v \in V_{i}$ and $u \in V_{j}$ are adjacent for $i \neq j$. In other words, an $r$-color partition can define an $r$-coloring of $G$.

Definition 7.1. Given $G$ (general) graph on $n$ vertices and given $u \in \mathbb{C}$, for each $r \in \mathbb{N}$ let $m_{r}(G)$ denote the number of $r$-color partitions of $G$ and let $u_{(r)}$ denote $u(u-1)(u-2) \cdots(u-r+1)$. We define the chromatic polynomial of $G$ to be

$$
C_{G}(u)=\sum_{r=1}^{n} m_{r}(G) u_{(r)}
$$

Proposition 7.2. If $s \in \mathbb{N}$, then $C_{G}(s)$ is the number of colorings of $G$ using at most $s$ colors.

Proof. We claim that $m_{r}(G) s_{(r)}$ is precisely the number of $s$-colorings of $G$ after an $r$-color partition (i.e. using precisely $r$ colors). If $r>s$ then there are no $s$-colorings allowed, and indeed $s_{(r)}=0$ if and only if $s \in \mathbb{N}$ and $s<r$. If $r \leq s$ then $s_{(r)}=\binom{s}{r} r$ ! and indeed, for each $r$-color partition, we have $\binom{s}{r}$ ways of choosing $r$ colors out of $s$ options, and then $r$ ! ways of permuting $r$ colors across the color classes.

Arguably the simplest graph to compute the chromatic polynomial on is the complete graph.

Example 7.3. Since $K_{n}$ has no color partition using less than $n$ color classes, $m_{1}\left(K_{n}\right)=\cdots=m_{n-1}\left(K_{n}\right)=0$ and $m_{n}\left(K_{n}\right)=1$, thus

$$
C_{K_{n}}(u)=u(u-1) \cdots(u-n+1)
$$

and we can observe that this is indeed the number of ways to assign $n$ colors out of $u$ options to color $K_{n}$.

Note that the chromatic polynomial is indeed a polynomial since each $u_{(r)}$ is a polynomial. Its degree is no more than $n$. Since $m_{n}(G)=1, C_{G}$ is monic.

For $G$ disconnected with components $G_{1}$ and $G_{2}$, for each way to color $G_{1}$ in $s \in \mathbb{N}$ colors or less paired with each way to $s$-color $G_{2}$, we have a distinct way to color $G$ and vice versa. Thus we have, by 7.2 ,

$$
C_{G}(s)=C_{G_{1}}(s) C_{G_{2}}(s)
$$

Since a polynomial is uniquely determined by its values on $\mathbb{N}$, we have

$$
C_{G}=C_{G_{1}} C_{G_{2}}
$$

Since $u$ is a factor in $u_{(r)}$ for all $r \geq 1$, we have $C_{G}(0)=0$ for any $G$ (i.e. the coefficient of $1=u^{0}$ is 0 ). If $G$ has $c$ components, then the coefficients of $u^{0}, u^{1}, \ldots, u^{c-1}$ are all 0 due to $C_{G}=\prod_{i=1}^{c} C_{G_{i}}$ where each $C_{G_{i}}$ has coefficient 0 on $u^{0}$. Finally, if $E(G) \neq \varnothing$ then $m_{1}(G)=0$ and $C_{G}(1)=0$ and $u-1$ is a factor of $C_{G}$.

Note that the problem of finding zeros of $C_{G}$ encompasses the problem of finding the chromatic number of $G, \chi(G)$, since the smallest natural number that is not a zero of $C_{G}$ is $\chi(G)$.

Now we discuss some techniques of calculating chromatic polynomials.

Definition 7.4. Suppose $e \in E(G)$ is not a loop. We define $G^{(e)}$ to be the graph with $V\left(G^{(e)}\right)=V(G)$ and $E\left(G^{(e)}\right)=E(G) \backslash e$ and say that $G^{(e)}$ is obtained from $G$ by deleting $e$. We define $G_{(e)}$ to be the graph with $e$ removed and with the two vertices incident to $e$ identified as one vertex and say that $G_{(e)}$ is obtained from $G$ by contracting $e$.

Note that $G^{(e)}$ has one edge fewer than $G$ and $G_{(e)}$ has one edge and one vertex fewer than $G$. Thus the following proposition provides a method to compute the chromatic polynomial by repeated reduction to smaller graphs. This is known as the deletion-contraction method.

## Proposition 7.5.

$$
C_{G}=C_{G^{(e)}}-C_{G_{(e)}} .
$$

Proof. We show $C_{G^{(e)}}=C_{G}+C_{G_{(e)}}$. Indeed for each way of coloring $G^{(e)}$, either the vertices incident to $e$ in $G$ are not colored the same, thus corresponding to a coloring of $G$, or they are, thus corresponding to a coloring of $G_{(e)}$. Thus $C_{G^{(e)}}(s)=C_{G}(s)+C_{G_{(e)}}(s)$ on $s \in \mathbb{N}$ and the general equality follows.

We use 7.5 to obtain the general formula for $C_{T}$ when $T$ is a tree:

Corollary 7.6. If $T$ is a tree with $n$ vertices, then

$$
C_{T}(u)=u(u-1)^{n-1}
$$

Proof. We use induction on $n$ to prove $C_{T}(u)=u(u-1)^{n-1}$. $n=1$ is trivial since there are $s$ ways to $s$-color one vertex. Assume it is true for $n=N$ and let $T$ be a tree with $N+1$ vertices. It is known that $T$ has at least one vertex of degree 1 (in fact, at least two vertices) so call it $v$. Let $e$ be the edge attaching $v$ to the rest of $T$. Let $T^{\prime}=\langle V(T) \backslash\{v\}\rangle$ and note that $T^{\prime}$ is a tree with $N$ vertices, so by $\mathrm{IH} C_{T^{\prime}}(u)=u(u-1)^{N-1}$. Then note $T_{(e)}=T^{\prime}$ so

$$
C_{T_{(e)}}(u)=C_{T^{\prime}}(u)=u(u-1)^{N-1} .
$$

$T^{(e)}$ has disconnected components $\langle\{v\}\rangle$ and $T^{\prime}$ so

$$
C_{T^{(e)}}(u)=C_{\langle\{v\}\rangle}(u) C_{T^{\prime}}(u)=u\left(u(u-1)^{N-1}\right) .
$$

Thus by 7.5

$$
C_{T}=C_{T^{(e)}}-C_{T_{(e)}}=u(u-1)^{N-1}(u-1)=u(u-1)^{N} .
$$

We can apply 7.5 and 7.6 together to obtain the general formula for $C_{C_{n}}$.

Example 7.7. $C_{C_{1}} \equiv 0$. When $n \geq 2$, let $e \in E\left(C_{n}\right)$, then $C_{n,(e)}=C_{n-1}$ and $C_{n}^{(e)}=P_{n}$ so

$$
\begin{aligned}
C_{C_{n}}(u) & =C_{C_{n}^{(e)}}(u)-C_{C_{n,(e)}}(u)=C_{P_{n}}(u)-C_{C_{n-1}}(u)=u(u-1)^{n-1}-C_{C_{n-1}}(u) \\
& =(u-1)^{n}-\left(C_{C_{n-1}}(u)-(u-1)^{n-1}\right) .
\end{aligned}
$$

We note that $C_{C_{n}}(u)-(u-1)^{n}=-\left(C_{C_{n-1}}(u)-(u-1)^{n-1}\right) \forall n \geq 2$ so we know

$$
C_{C_{n}}(u)=(u-1)^{n}+(-1)^{n} f(u)
$$

for some $f(u)$ independent of $n$. Since the base case is $C_{C_{1}}=0=u-1+(-1)(u-1)$, we know $f(u)=u-1$ so

$$
C_{C_{n}}(u)=(u-1)^{n}+(-1)^{n}(u-1)
$$

Now we describe two additional methods of calculating chromatic polynomials. The first is based on a join operation for graphs:

Definition 7.8. Let $G_{1}, G_{2}$ be graphs (with $V\left(G_{1}\right) \cap V\left(G_{2}\right)=\varnothing$ ). We define their join $G_{1}+G_{2}$ to be the graph with

$$
V\left(G_{1}+G_{2}\right)=V\left(G_{1}\right) \cup V\left(G_{2}\right)
$$

and

$$
E\left(G_{1}+G_{2}\right)=E\left(G_{1}\right) \cup E\left(G_{2}\right) \cup\left\{(u, v): u \in V\left(G_{1}\right), v \in V\left(G_{2}\right)\right\}
$$

In other words, $G_{1}+G_{2}$ contains $G_{1}$ and $G_{2}$ as disjoint subgraphs and edges joining all vertices between them, and nothing else.

Proposition 7.9. Suppose $G=G_{1}+G_{2}$, then

$$
m_{r}(G)=\sum_{i+j=r} m_{i}\left(G_{1}\right) m_{j}\left(G_{2}\right)
$$

Proof. For each $r$-color partition of $G$, the colors used on $G_{1}$ and $G_{2}$ must be disjoint, so $G_{1}$ must have an $i$-color partition and $G_{2}$ a $j$-color partition such that $i+j=r$ and conversely any such pair indeed gives an $r$-color partition of $G$.

This directly give us the following formula for computing chromatic polynomials of joins:
Corollary 7.10. Define

$$
\sum_{r=1}^{n} m(r) u_{(r)} \circ \sum_{s=1}^{\ell} m(s) u_{(s)}=\sum_{r=1}^{n} \sum_{s=1}^{\ell} m(r) m(s) u_{(r+s)}=\sum_{t=1}^{m+n} \sum_{r+s=t} m(r) m(s) u_{(t)}
$$

Then

$$
C_{G_{1}+G_{2}}=C_{G_{1}} \circ C_{G_{2}}
$$

Here use 7.10 to derive $C_{K_{3,3}}$ :

Example 7.11. Note that $K_{3,3}=N_{3}+N_{3}$, where $N_{k}$ defines a graph with $k$ vertices and no edges. Since there are 3 ways to partition a set of 3 elements into 2 classes, note

$$
C_{N_{3}}(u)=u_{(1)}+3 u_{(2)}+u_{(3)} .
$$

Thus

$$
\begin{aligned}
C_{K_{3,3}}(u) & =C_{N_{3}}(u) \circ C_{N_{3}}(u)=\left(u_{(1)}+3 u_{(2)}+u_{(3)}\right) \circ\left(u_{(1)}+3 u_{(2)}+u_{(3)}\right) \\
& =u_{(2)}+3 u_{(3)}+u_{(4)}+3 u_{(3)}+9 u_{(4)}+3 u_{(5)}+u_{(4)}+3 u_{(5)}+u_{(6)} \\
& =u_{(2)}+6 u_{(3)}+11 u_{(4)}+6 u_{(5)}+u_{(6)} .
\end{aligned}
$$

Theoretically, we can apply the same argument to all $K_{n, m}$ and algorithmically find its chromatic polynomial. However note that there is no known closed-form expression for the number of ways to partition $k$ elements into $r$ classes, so we cannot obtain a closed-form $C_{K_{n, m}}$ this way for all $n, m \in \mathbb{N}$.

Another application of 7.10 relates the chromatic polynomial of $G$ with that of $N_{1}+G$, called the cone of $G$ and denoted $c(G)$, and $N_{2}+G$, called the suspension of $G$ and denoted $s(G)$.

Proposition 7.12. The chromatic polynomial of $a$ cone and a suspension are given by

$$
\begin{gather*}
C_{c(G)}(u)=u C_{G}(u-1)  \tag{7.1}\\
C_{s(G)}(u)=u(u-1) C_{G}(u-2)+u C_{G}(u-1) \tag{7.2}
\end{gather*}
$$

Proof. Note $C_{N_{1}}=u_{(1)}$ and $C_{N_{2}}=u_{(1)}+u_{(2)}$. (7.1) is due to 7.10 and the fact that $u_{(r+1)}=u(u-1)_{(r)}$, so

$$
C_{c(G)}(u)=u_{(1)} \circ C_{G}(u)=u C_{G}(u-1)
$$

(7.2) is similar, using $u_{(r+2)}=u(u-1)(u-2)_{(r)}$, so

$$
C_{s(G)}(u)=\left(u_{(1)}+u_{(2)}\right) \circ C_{s(G)}(u)=u(u-1) C_{G}(u-2)+u C_{G}(u-1) .
$$

The second technique applies to graphs described as follows:

Definition 7.13. The (general) graph $G$ is quasi-separable if there is $K \subseteq V(G)$ such that the induced subgraph $\langle K\rangle$ is complete but $\langle V(G) \backslash K\rangle$ is disconnected. $G$ is separable if $|K| \leq 1$; in this case either $K=\varnothing$ so $G$ is in fact disconnected, or $|K|=1$, in which case we call $v \in K$ a cut-vertex.

It follows that by taking each of the disconnected components and unioning each with $K$, we can get $V_{1}, V_{2} \subseteq V(G)$ such that $V_{1} \cup V_{2}=V(G), V_{1} \cap V_{2}=K$ and there are no edges between $V_{1} \backslash K$ and $V_{2} \backslash K$. We call $\left(V_{1}, V_{2}\right)$ a quasi-separation, or a separation if $|K| \leq 1$.

The following reduces the chromatic polynomial of a graph with quasi-separation $\left(V_{1}, V_{2}\right)$ to the chromatic polynomials of $\left\langle V_{1}\right\rangle,\left\langle V_{2}\right\rangle$ and $\left\langle V_{1} \cap V_{2}\right\rangle$ :

Proposition 7.14. If $G$ has quasi-separation $\left(V_{1}, V_{2}\right)$, then

$$
C_{G}=\frac{C_{\left\langle V_{1}\right\rangle} C_{\left\langle V_{2}\right\rangle}}{C_{\left\langle V_{1} \cap V_{2}\right\rangle}}
$$

with the convention that $C_{\langle\varnothing\rangle}(u)=1$.

Proof. If $V_{1} \cap V_{2}=\varnothing$, then $G$ is disjoint and the result follows a previous remark about 7.2. Suppose that $\langle K\rangle=\left\langle V_{1} \cap V_{2}\right\rangle \cong K_{t}$ with $t \geq 1$. Let $s \in \mathbb{N}$.

Due to the fact that all colorings on $\left\langle K_{t}\right\rangle$ are isomorphic (i.e. one can bijectively map each color to some other color to go from one coloring to another), each coloring on $\langle K\rangle$ can be extended in $C_{G}(s) / s_{(t)}$ ways to a coloring in $G$ (because there exists a bijection between extensions starting from different colorings of $\langle K\rangle$ ). Since $K_{t} \subseteq\left\langle V_{1}\right\rangle,\left\langle V_{2}\right\rangle$, the same argument applies to $C_{\left\langle V_{1}\right\rangle}(s) / s_{(t)}$ and $C_{\left\langle V_{2}\right\rangle}(s) / s_{(t)}$.

Because $\left\langle V_{1} \backslash K\right\rangle$ and $\left\langle V_{2} \backslash K\right\rangle$ have no edges between them, extensions for $\left\langle V_{1}\right\rangle$ and $\left\langle V_{2}\right\rangle$ starting from the same coloring on $K_{t}$ are independent, hence

$$
\begin{aligned}
& \frac{C_{G}(s)}{s_{(t)}}=\frac{C_{\left\langle V_{1}\right\rangle}(s)}{s_{(t)}} \frac{C_{\left\langle V_{1}\right\rangle}(s)}{s_{(t)}} \\
\Longrightarrow C_{G}(s) & =\frac{C_{\left\langle V_{1}\right\rangle}(s) C_{\left\langle V_{1}\right\rangle}(s)}{s_{(t)}}=\frac{C_{\left\langle V_{1}\right\rangle}(s) C_{\left\langle V_{1}\right\rangle}(s)}{C_{\left\langle V_{1} \cap V_{2}\right\rangle}(s)} \\
\Longrightarrow C_{G}= & \frac{C_{\left\langle V_{1}\right\rangle} C_{\left\langle V_{2}\right\rangle}}{C_{\left\langle V_{1} \cap V_{2}\right\rangle}} .
\end{aligned}
$$

7.14 is often useful for finding chromatic polynomials of small graphs for which the subgraphs induced by the quasi-separation are well-known. The following is an example:


Figure 2: G separable with $V_{1}=\{1,2,3\}$ and $V_{2}=\{2,3,4\}$

Example 7.15. Consider the graph in Figure 2. Under the given separation and using the fact that $\left\langle V_{1}\right\rangle \cong$ $\left\langle V_{2}\right\rangle \cong K_{3}$ while $\left\langle V_{1} \cap V_{2}\right\rangle \cong K_{2}$, we get

$$
C_{G}(u)=\frac{C_{\left\langle V_{1}\right\rangle}(u) C_{\left\langle V_{2}\right\rangle}(u)}{C_{\left\langle V_{1} \cap V_{2}\right\rangle}(u)}=\frac{(u(u-1)(u-2))^{2}}{u(u-1)}=u(u-1)(u-2)^{2} .
$$

Note that there are other ways to compute $C_{G}$. For example, $G$ can also be viewed as $c P_{3}$, in other words $N_{1}+P_{3}$ where $V\left(N_{1}\right)=\{3\}$ and $V\left(P_{3}\right)=\{1,2,4\}$, in which case

$$
C_{G}(u)=u C_{P_{3}}(u-1)=u(u-1)(u-2)^{2}
$$

## 8 Bounds on the complex zeros of chromatic polynomials

This section consists of notes on a paper by Sokal 2000.

### 8.1 Introduction

We have previously introduced (in 7) the chromatic polynomial for (general) graph $G$. For each $e \in E$ let $v_{e} \in \mathbb{C}$. Then we define (let $s \in \mathbb{N}, x_{e}, y_{e} \in V$ be the two endpoints of $e, \mathbf{1}$ denote the indicator function, and product over the empty set understood to give 1)

$$
\begin{equation*}
Z_{G}\left(s,\left\{v_{e}\right\}\right)=\sum_{\sigma: V \rightarrow\{1, \ldots, s\}} \prod_{e \in E}\left[1+v_{e} \mathbf{1}\left\{\sigma\left(x_{e}\right)=\sigma\left(y_{e}\right)\right\}\right] . \tag{8.1}
\end{equation*}
$$

We will show below (8.1) that (8.1) is the restriction of a polynomial to $\mathbb{N} \times \mathbb{C}^{|E|}$. If one takes $v_{e}=-1 \forall e \in E$, then

$$
\prod_{e \in E}\left[1-\mathbf{1}\left\{\sigma\left(x_{e}\right)=\sigma\left(y_{e}\right)\right\}\right]=\mathbf{1}\left\{\sigma\left(x_{e}\right) \neq \sigma\left(y_{e}\right) \forall e \in E\right\}=\mathbf{1}\{\sigma \text { is an } s \text {-coloring of } V\}
$$

so $Z_{G}(s,-1, \ldots,-1)=C_{G}(s)$ is precisely the chromatic polynomial. If one takes $v_{e}=v \forall e \in E$ then the two-variable $Z_{G}(s, v)$ is called, up to some trivial transformation (see 8.2), the dichromatic polynomial or the Tutte polynomial.

Proposition 8.1. $Z_{G}\left(s,\left\{v_{e}\right\}\right)$ is a polynomial in its arguments with 1 as the only non-zero coefficient.

Proof. First we expand the multiplication in (8.1) ( $E^{\prime}$ is taken over all subsets of $\left.E\right)$ :

$$
\begin{aligned}
Z_{G}\left(s,\left\{v_{e}\right\}\right) & =\sum_{\sigma} \prod_{e \in E}\left[1+v_{e} \mathbf{1}\left\{\sigma\left(x_{e}\right)=\sigma\left(y_{e}\right)\right\}\right] \\
& =\sum_{\sigma} \sum_{E^{\prime} \subseteq E} \prod_{e \in E^{\prime}} v_{e} \mathbf{1}\left\{\sigma\left(x_{e}\right)=\sigma\left(y_{e}\right)\right\} \\
& =\sum_{\sigma} \sum_{E^{\prime} \subseteq E} \mathbf{1}\left\{\sigma\left(x_{e}\right)=\sigma\left(y_{e}\right) \forall e \in E^{\prime}\right\} \prod_{e \in E^{\prime}} v_{e} \\
& =\sum_{E^{\prime} \subseteq E}\left|\left\{\sigma: \sigma\left(x_{e}\right)=\sigma\left(y_{e}\right) \forall e \in E^{\prime}\right\}\right| \prod_{e \in E^{\prime}} v_{e}
\end{aligned}
$$

Fix $E^{\prime} \subseteq E$. For any $\sigma$, note that $\sigma\left(x_{e}\right)=\sigma\left(y_{e}\right) \forall e \in E^{\prime}$ if and only if all components of the subgraph induced by $E^{\prime}$ are colored in the same color. Since each component would have $s$ choices of colors, let $c\left(E^{\prime}\right)$ denote the number of components in the said subgraph, then there are $s^{c\left(E^{\prime}\right)}$ possibilities of $\sigma$ for a given $E^{\prime} \subseteq E$. Thus

$$
\begin{equation*}
Z_{G}\left(s,\left\{v_{e}\right\}\right)=\sum_{E^{\prime} \subseteq E} s^{c\left(E^{\prime}\right)} \prod_{e \in E^{\prime}} v_{e} \tag{8.2}
\end{equation*}
$$

Since (8.1) was only defined for $s \in \mathbb{N}$, we take (8.2) to be the definition of $Z_{G}\left(s,\left\{v_{e}\right\}\right)$ for $s,\left\{v_{e}\right\} \in \mathbb{C}$. We now show the relationship between the two-variable $Z_{G}(s, v)$ (i.e. $v_{e}=v \forall e \in E$ ) and the conventional Tutte polynomial $T_{G}(x, y)$.

Corollary 8.2. Define

$$
T_{G}(x, y)=\sum_{E^{\prime} \subseteq E}(x-1)^{c\left(E^{\prime}\right)-c(E)}(y-1)^{\left|E^{\prime}\right|+c\left(E^{\prime}\right)-|V|}
$$

Then

$$
T_{G}(x, y)=(x-1)^{-c(E)}(y-1)^{-|V|} Z_{G}((x-1)(y-1), y-1) .
$$

Proof.

$$
Z_{G}((x-1)(y-1), y-1)=\sum_{E^{\prime} \subseteq E}(x-1)^{c\left(E^{\prime}\right)}(y-1)^{\left|E^{\prime}\right|+c\left(E^{\prime}\right)} .
$$

In the context of statistical mechanics, (8.1) is known as the partition function of the s-state Potts model ${ }^{6}$. In this model, each site $x \in V$ can exist in any of the $s \in \mathbb{N}$ different 'states' (e.g. spins). The energy of $e \in E$ is 0 if the states of the end points are unequal and $-J_{e}$ if they are equal, and we sum over all $e \in E$ to get $H(\sigma)$, the energy of a given configuration $\sigma$. The Boltzmann weight of $\sigma$ is then $e^{-\beta H(\sigma)}$, where $\beta \geq 0$ is the inverse temperature ${ }^{7}$. The partition function is the sum of Boltzmann weights over all $\sigma$. It is easy to see that by taking $v_{e}=e^{\beta J_{e}}-1 \forall e \in E$, the partition function is equivalent to (8.1). An interaction $e \in E$ is called ferromagnetic if $J_{e}>0 \Longleftrightarrow v_{e}>0$, antiferromagnetic if $-\infty \leq J_{e}<0 \Longleftrightarrow-1 \leq v_{e}<0$, and otherwise the ends of $e$ are non-interacting (i.e. $J_{e}=0 \Longleftrightarrow v_{e}=0$, equivalent to $e$ being removed).

We remark that the partition function (provided it is non-zero) serves as a normalizing constant for the Boltzmann weights, such that

$$
f_{G, s,\left\{v_{e}\right\}}(\sigma)=\frac{e^{-\beta H(\sigma)}}{Z_{G}\left(s,\left\{v_{e}\right\}\right)}
$$

is a probability distribution (called an ensemble in statistical mechanics) over possible configurations $\sigma$. In fact, complex zeros of $Z_{G}$ are often particular interest, as explained below.

Often, one would want to study phase transitions, i.e. points where some physical quantities (e.g. energy) depend nonanalytically (or even discontinuously) on the parameters of the system (e.g. temperature or the magnetic field, i.e. $J_{e}$ ). Such points are impossible in (8.1) for any finite $G$, but they could arise in an infinite-volume limit, i.e. some limit $G_{\infty}$ of increasing graphs $G_{i}$ with $n\left(G_{i}\right) \rightarrow \infty$. Regular lattices such as the $d$-dimensional integer lattice are typical examples. It can then be shown (under modest assumptions on $G_{i}$ ) that the (limiting) free energy per unit volume

$$
f_{G_{\infty}}(s, v)=\lim _{i \rightarrow \infty}\left|n\left(G_{i}\right)\right|^{-1} \log Z_{G_{i}}(s, v)
$$

[^4]exists for all nondegenerate physical ${ }^{8}$ values of the parameters, namely either

1. $s \in \mathbb{N}$ and $-1<v<\infty$ (using (8.1)), or
2. $s>0$ and $0 \leq v<\infty$ (using (8.2)).

The limit $f_{G_{\infty}}$ is in general continuous in $v$, but can fail to be real-analytic in $v$, because complex singularities of $\log Z_{G_{i}}$, i.e. complex zeros of $Z_{G_{i}}$, can approach the real axis as $i \rightarrow \infty$. Therefore the real limits of such zeros are precisely the points of interest when studying phase transitions, so theorems on the location of the zeros of the partition function are very important. ${ }^{9}$

The purpose of the paper was to give an upper bound on complex zeros of $Z_{G}$ for $v_{e}$ in the 'complex antiferromagnetic regime' $A=\{v \in \mathbb{C}:|1+v| \leq 1\}$, i.e. $\left|1+v_{e}\right| \leq 1$. This bound can be valid for infinite families of $G$, if some local conditions on $v_{e}$ hold. A corollary is an upper bound on the zeros of the chromatic polynomial based on the maximum degree of $G$ : for each $r \geq 0$, there exists $C(r)$ such that for any loopless $G$ with maximum degree $r$, the zeros of $C_{G}(u)$ lie in $|u|<C(r)$. It is also proven that $C(r)$ grows at most linearly in $r$. Note that $C(r)$ grows at least linearly because $r+1$ is a root of $C_{K_{r+1}}$. An explicit bound of $C(r) \leq 7.963907 r$ was given. Finally, with one vertex of degree exceeding $r$, the zeros of $C_{G}(u)$ are bounded by $|u|<C(r)+1$. It is known that the roots of $C_{G}$ are unbounded when there are two such vertices.

These notes will discuss the bound on the zeros of $Z_{G}$ for a fixed $G$ under $\left|1+v_{e}\right| \leq 1$ as well as the corollary on $C_{G}$. We will also show that $C(r)$ grows at most linearly, but with a worse bound than the one given.

### 8.2 Transformation of the Potts-model partition function to a polymer gas

Let $G=(V, E)$ have complex edge weights $\left\{v_{e}\right\}_{e \in E}$. If $e$ is a loop then by (8.1) its presence simply multiplies $Z_{G}$ by $\left(1+v_{e}\right)$, thus WLOG we assume $G$ is loopless. ${ }^{10}$ For each $E^{\prime} \subseteq E$, if we decompose ( $V, E^{\prime}$ ) into connected components $S_{1}, \ldots, S_{N}$, note that $c\left(E^{\prime}\right)=\sum_{i=1}^{N}\left(1+\left|S_{i}\right|-\left|S_{i}\right|\right)=|V|-\sum_{i=1}^{N}\left(\left|S_{i}\right|-1\right)$.

Proposition 8.3. Let $G=(V, E)$ be loopless with edge weights $\left\{v_{e}\right\}_{e \in E}$. Then

$$
\begin{equation*}
Z_{G}\left(s,\left\{v_{e}\right\}\right)=s^{|V|} Z_{\text {polymer }, G}\left(s,\left\{v_{e}\right\}\right), \tag{8.3}
\end{equation*}
$$

where

$$
\begin{equation*}
Z_{\text {polymer }, G}\left(s,\left\{v_{e}\right\}\right)=\sum_{N=0}^{\infty} \frac{1}{N!} \sum_{\bigsqcup_{i=1}^{N} S_{i} \subseteq V} \prod_{i=1}^{N} w\left(S_{i}\right) \tag{8.4}
\end{equation*}
$$

and where

$$
w(S)= \begin{cases}s^{-(|S|-1)} \sum_{E^{\prime} \subseteq E_{S},\left(S, E^{\prime}\right) \text { connected }} \prod_{e \in E^{\prime}} v_{e} & |S| \geq 2  \tag{8.5}\\ 0 & |S| \leq 1\end{cases}
$$

In particular note $w(S)=0$ if $\left(S, E_{S}\right)$ is not connected, i.e. if $S$ is not the subset of a connected component.

[^5]Proof. Clearly we only need to consider $\bigsqcup_{i=1}^{N} S_{i}=S \subseteq V$ where each $\left(S_{i}, E_{S_{i}}\right)$ is connected and where each $\left|S_{i}\right|>1$, due to which $E_{S_{i}} \neq \varnothing$. In that case

$$
\prod_{i=1}^{N} w\left(S_{i}\right)=\prod_{i=1}^{N} s^{-\left(\left|S_{i}\right|-1\right)} \sum_{\substack{E^{\prime} \subseteq E_{S_{i}},\left(S_{i}, E^{\prime}\right) \text { connected }}} \prod_{e \in E^{\prime}} v_{e}=\sum_{\substack{E^{\prime}=\bigsqcup_{i=1}^{N} E_{i}^{\prime}: \\\left(E_{i}^{\prime} \subseteq E_{S_{i}},\left(S_{i}, E_{i}^{\prime}\right)\right. \text { connected }}} s^{-\sum_{i=1}^{N}\left(\left|S_{i}\right|-1\right)} \prod_{e \in E^{\prime}} v_{e}
$$

Note that since $|S|-\sum_{i=1}^{N}\left(\left|S_{i}\right|-1\right)=N$ which is the number of components in $\left(S, E^{\prime}\right)$, and there are precisely $|V|-|S|$ more components in $\left(V, E^{\prime}\right)$ than in $\left(S, E^{\prime}\right)$, we have $c\left(E^{\prime}\right)=|V|-\sum_{i=1}^{N}\left(\left|S_{i}\right|-1\right)$. Thus (8.3) is

$$
\begin{equation*}
s^{|V|} Z_{\text {polymer }, G}\left(s,\left\{v_{e}\right\}\right)=\sum_{N=0}^{\infty} \frac{1}{N!} \sum_{\substack{\bigsqcup_{i=1}^{N}}} \sum_{\substack{S_{i}=S \subseteq V \\ E^{\prime}=\bigsqcup_{i=1}^{N} \\ E_{i}^{\prime} \subseteq E_{S_{i}},\left(S_{i}^{\prime}, E_{i}^{\prime}\right) \text { connected }}} s^{c\left(E^{\prime}\right)} \prod_{e \in E^{\prime}} v_{e} \tag{8.6}
\end{equation*}
$$

For each $E^{\prime} \subseteq E$ such that $\left(V, E^{\prime}\right)$ has $N$ components $S_{1}, \ldots, S_{N}$ where each $\left|S_{i}\right|>1$, there are precisely $N$ ! possible ways to permute those $S_{i}$ such that $S_{1}, \ldots, S_{N}$ are still the connected components of ( $S, E^{\prime}$ ), so such $E^{\prime}$ appears throughout the summation in (8.6) exactly $N$ ! times. By summing across all $N$, we taking into account every possible $E^{\prime} \subseteq E$ so

$$
s^{|V|} Z_{\text {polymer, } G}\left(s,\left\{v_{e}\right\}\right)=\sum_{E^{\prime} \subseteq E} s^{c\left(E^{\prime}\right)} \prod_{e \in E^{\prime}} v_{e}=Z_{G}\left(s,\left\{v_{e}\right\}\right)
$$

The 'polymer model' (8.4) has the form of a grand-canonical gas (8.4)

$$
\begin{equation*}
Z_{\text {polymer }, G}\left(s,\left\{v_{e}\right\}\right)=\sum_{N=0}^{\infty} \frac{1}{N!} \sum_{S_{1}, \ldots, S_{N} \subseteq V} \prod_{i=1}^{N} w\left(S_{i}\right) \prod_{1 \leq i, j \leq N} 1\left\{S_{i} \cap S_{j}=\varnothing\right\} \tag{8.7}
\end{equation*}
$$

with single-particle state space $\mathscr{P}_{*}(V)$ (non-empty subsets of $V$, or, equivalently, due to how $w(S)$ was defined, $S \subseteq V$ with $|S| \geq 2$ and connected $\left(S, E_{S}\right)$ ), fugacities $w(S)$ and two-particle Boltzmann factor $1\left\{S_{i} \cap S_{j}=\varnothing\right\}$.

If we define an exponential generating function (EGF) in $M$ variables w.r.t. sequence $\left\{a\left(n_{1}, \ldots, n_{M}\right)\right\}_{n_{i} \in \mathbb{N}}$ as

$$
\begin{equation*}
G\left(\left\{a_{n}\right\}, x_{1}, \ldots, x_{M}\right)=\sum_{N=0}^{\infty} \frac{1}{N!} \sum_{\sum_{m=1}^{M} n_{m}=N} a\left(n_{1}, \ldots, n_{M}\right) \prod_{m=1}^{M} x_{m}^{n_{m}} \tag{8.8}
\end{equation*}
$$

then (8.7) is precisely (8.8) with $M=\left|\mathscr{P}_{*}(V)\right|$ with each $m$ associated to some $S \in \mathscr{P}_{*}(V), x_{S}=w(S) \forall S \in$ $\mathscr{P}_{*}(V)$ and $a\left(\left\{n_{S}\right\}_{S \in \mathscr{P}_{*}(V)}\right)$ is an indicator function for the following: $n_{S}=1$ or $n_{S}=0$ for each $S$, and $S \cap S^{\prime}=\varnothing$ for each $S, S^{\prime}$ where $n_{S}, n_{S^{\prime}}=1$. (8.7) is thus the generating function for independent sets in the intersection graph constructed from $\mathscr{P}_{*}(V)$ with variables $w(S)$.

We remark that since $-(|S|-1) \leq-1$ when $w(S) \neq 0, w(S)$ decreases in $s$. Thus if the sum over $E^{\prime}$ can be controlled, one can expect an exponential decay of $w(S)$ in $|S|$ provided $s$ is large enough so that the exponential decay overshadows the increase in the sum over $E^{\prime}$. Thus as we show that $Z_{\text {polymer, } G}$ does not vanish for small enough $w(S)$, we have reason to believe that the same is true for large enough $s$. This essentially outlines the following sections, but in the opposite order.

### 8.3 Dobrushin and Kotechý-Preiss conditions for the nonvanishing of $Z$

Definition 8.4. A grand-canonical gas is defined by a single-particle state space $X$ (here finite), a fugacity vector $\mathbf{w}=\left\{w_{x}\right\}_{x \in X} \in \mathbb{C}^{|X|}$ and a symmetric two-particle Boltzmann factor $W: X \times X \rightarrow \mathbb{C}$. The (grand) partition function $z(\mathbf{w}, W)$ is then defined to be the sum over ways of placing $N \geq 0$ particles on sites $x_{1}, \ldots, x_{N} \in X$, with each configuration given a Boltzmann weight, which is the product over all $w_{x_{i}}$ and $W\left(x_{i}, x_{j}\right)($ for $i \neq j)$ :

$$
\begin{equation*}
Z(\mathbf{w}, W)=\sum_{N=0}^{\infty} \frac{1}{N!} \sum_{x_{1}, \ldots, x_{N} \in X} \prod_{i=1}^{N} w_{x_{i}} \prod_{1 \leq i<j \leq N} W\left(x_{i}, x_{j}\right) \tag{8.9}
\end{equation*}
$$

Under weak assumptions on $W$ (e.g. $|W(x, y)| \leq 1$ like in (8.7)) $Z(\mathbf{w}, W)$ is analytic in $\mathbf{w}$. We want to find a sufficient condition for $Z(\mathbf{w}, W)$ to be nonvanishing in a polydisc $D_{R}=\left\{\mathbf{w}:\left|w_{x}\right|<R_{x}\right\}$. This would imply in particular that $\log Z(\mathbf{w}, W)$ is analytic in $D_{R}$.

Note that if $W$ were hard-core self-repulsive, i.e. $W(x, x)=0 \forall x \in X$, then we only need to consider summation over distinct $x_{i}$, i.e. (8.9) would be equivalent to summing over subsets:

$$
Z(\mathbf{w}, W)=\sum_{X^{\prime} \subseteq X} \prod_{x \in X^{\prime}} w_{x} \prod_{\{x, y\} \subseteq X^{\prime}} W(x, y) .
$$

Under this assumption, we introduce the notation where, for each $\Lambda \subseteq X$, we write ${ }^{11}$

$$
\begin{equation*}
Z_{\Lambda}(\mathbf{w}, W)=\sum_{X^{\prime} \subseteq \Lambda} \prod_{x \in X^{\prime}} w_{x} \prod_{\{x, y\} \in X^{\prime}} W(x, y) \tag{8.10}
\end{equation*}
$$

Now we give an extension of a theorem due to Dobrushin:

Theorem 8.5 (Extension to Dobrushin). Let $X$ be finite and let $W$ satisfy

1. $0 \leq W(x, y) \leq 1 \forall x, y \in X$;
2. $W(x, x)=0 \forall x \in X)$.

Suppose for each $x \in X$, there exists $R_{x} \geq 0$ and $0 \leq K_{x}<1 / R_{x}$ satisfying

$$
\begin{equation*}
K_{x} \geq \prod_{y \in X: y \neq x} \frac{1-W(x, y) K_{y} R_{y}}{1-K_{y} R_{y}} \tag{8.11}
\end{equation*}
$$

Then, for each $\Lambda \subseteq X, Z_{\Lambda}(\mathbf{w}, W)$ is nonvanishing in the closed polydisc $\overline{D_{R}}=\left\{\mathbf{w} \in \mathbb{C}^{|X|}:\left|w_{x}\right| \leq R_{x}\right\}$ and satisfies in $\overline{D_{R}}$

$$
\left|\frac{\partial \log Z_{\Lambda}(\mathbf{w}, W)}{\partial w_{x}}\right| \leq \begin{cases}\frac{K_{x}}{1-K_{x}\left|w_{x}\right|} & x \in \Lambda  \tag{8.12}\\ 0 & x \in X \backslash \Lambda\end{cases}
$$

Moreover, given $\mathbf{w}, \mathbf{w}^{\prime} \in \bar{D}_{R}$ such that $\mathbf{w}_{x}^{\prime} / \mathbf{w}_{x} \in[0, \infty] \forall x \in X$, by integration, one can get ${ }^{12}$

$$
\begin{equation*}
\left|\log \frac{Z_{\Lambda}\left(\mathbf{w}^{\prime}, W\right)}{Z_{\Lambda}(\mathbf{w}, W)}\right| \leq \sum_{x \in \Lambda}\left|\log \frac{1-K_{x}\left|w_{x}^{\prime}\right|}{1-K_{x}\left|w_{x}\right|}\right| \tag{8.13}
\end{equation*}
$$

[^6]We remark that since $W(x, y) K_{y} R_{y} \leq K_{y} R_{y}$, by (8.11) we have $K_{x} \geq 1$ and thus $R_{x}<1$.
Proof. On any $\Lambda$, (8.12) implies (8.13) because let $\mathbf{w}(t):[0,1] \rightarrow \mathbb{C}^{|X|}, t \mapsto \mathbf{w}(1-t)+\mathbf{w}^{\prime} t$, then $\mathbf{w}(0)=\mathbf{w}$ and $\mathbf{w}(1)=\mathbf{w}^{\prime}$ and note $\left|\frac{\partial w_{x}}{t}\right|=\left|w_{x}^{\prime}-w_{x}\right|=\left|\left|w_{x}^{\prime}\right|-\left|w_{x}\right|\right|=\left|\frac{\partial\left|w_{x}\right|}{\partial t}\right|$, thus

$$
\begin{aligned}
& \left|\frac{\partial}{\partial t} \log Z_{\Lambda}(\mathbf{w}(t), W)\right|=\left|\sum_{x \in X} \frac{\partial \log Z_{\Lambda}(\mathbf{w}, W)}{\partial w_{x}} \frac{\partial w_{x}}{\partial t}\right| \leq \sum_{x \in X}\left|\frac{\partial \log Z_{\Lambda}(\mathbf{w}, W)}{\partial w_{x}} \frac{\partial w_{x}}{\partial t}\right| \\
& \leq \sum_{x \in X} \frac{K_{x}}{1-K_{x}\left|w_{x}\right|}\left|\frac{\partial\left|w_{x}\right|}{\partial t}\right|=\sum_{x \in X}\left|\frac{\partial}{\partial t} \log \left(1-K_{x}\left|w_{x}(t)\right|\right)\right| \\
\Longrightarrow & \left|\int_{0}^{1} \frac{\partial}{\partial t} \log Z_{\Lambda}(\mathbf{w}(t), W) d t\right| \leq \sum_{x \in X} \int_{0}^{1}\left|\frac{\partial}{\partial t} \log \left(1-K_{x}\left|w_{x}(t)\right|\right)\right| d t \\
\Longrightarrow & \left|\log \frac{Z_{\Lambda}\left(\mathbf{w}^{\prime}, W\right)}{Z_{\Lambda}(\mathbf{w}, W)}\right| \leq \sum_{x \in \Lambda}\left|\log \frac{1-K_{x}\left|w_{x}^{\prime}\right|}{1-K_{x}\left|w_{x}\right|}\right|
\end{aligned}
$$

The rest of the proof is by induction on the cardinality $\Lambda$. The claim is vacuously true for $\Lambda=\varnothing$. Assume (8.12) (and thus (8.13)) hold for all sets with cardinality less than $n \in \mathbb{N}$ and let $|\Lambda|=n$. Let $X \in \Lambda$ and let $\Lambda^{\prime}=\Lambda \backslash\{x\}$. Then from (8.10), by breaking down the summation over $X^{\prime} \subseteq \Lambda$ into $X^{\prime}$ containing $x$ and those that do not, we get

$$
\begin{equation*}
Z_{\Lambda}(\mathbf{w}, W)=w_{x} Z_{\Lambda^{\prime}}(\tilde{\mathbf{w}}, W)+Z_{\Lambda^{\prime}}(\mathbf{w}, W) \tag{8.14}
\end{equation*}
$$

where $\tilde{w}_{y}=W(x, y) w_{y}$, so

$$
w_{x} Z_{\Lambda^{\prime}}(\tilde{\mathbf{w}}, W)=\sum_{X^{\prime}: x \notin X^{\prime}} w_{x} \prod_{y \in X^{\prime}} w_{y} W(x, y) \prod_{\left\{y, y^{\prime}\right\} \in X^{\prime}} W\left(y, y^{\prime}\right)=\sum_{X^{\prime}: x \in X^{\prime}} \prod_{y \in X^{\prime}} w_{y} \prod_{\left\{y, y^{\prime} \in X^{\prime}\right\}} W\left(y, y^{\prime}\right)
$$

We note $\tilde{\mathbf{w}} \in \overline{D_{R}}$ since $|W(x, y)| \leq 1$. From (8.14) we directly get

$$
\frac{\partial}{\partial w_{x}} \log Z_{\Lambda}(\mathbf{w}, W)=\frac{Z_{\Lambda^{\prime}}(\tilde{\mathbf{w}}, W)}{w_{x} Z_{\Lambda^{\prime}}(\tilde{\mathbf{w}}, W)+Z_{\Lambda^{\prime}}(\mathbf{w}, W)}=\frac{k(w)}{k(w) w_{x}+1}
$$

where $k(w)=\frac{Z_{\Lambda^{\prime}}(\tilde{\mathbf{w}}, W)}{Z_{\Lambda^{\prime}}(\mathbf{w}, W)}$. Finally since $\tilde{w}_{y} / w_{y}=W(x, y) \geq 0$, by IH we can apply (8.13) to bound $|k(w)|$ with

$$
\begin{equation*}
|k(w)| \leq \prod_{y \in \Lambda^{\prime}} \frac{1-K_{y}\left|\tilde{w}_{y}\right|}{1-K_{y}\left|w_{y}\right|}=\prod_{y \in \Lambda^{\prime}} \frac{1-K_{y} W(x, y)\left|w_{y}\right|}{1-K_{y}\left|w_{y}\right|} \leq \prod_{y \in X \backslash x} \frac{1-K_{y} W(x, y)\left|w_{y}\right|}{1-K_{y}\left|w_{y}\right|} \tag{8.15}
\end{equation*}
$$

Note

$$
\frac{1-K_{y} W(x, y)\left|w_{y}\right|}{1-K_{y}\left|w_{y}\right|} \leq \frac{1-K_{y} W(x, y) R_{y}}{1-K_{y} R_{y}} \Longleftrightarrow-K_{y} R_{y} \leq-K_{y}\left|w_{y}\right|
$$

which is true so by $(8.11)|k(w)| \leq K_{x}$. Note $\frac{|k(w)|}{\left|k(w) w_{x}+1\right|} \leq \frac{|k(w)|}{1-|k(w)|\left|w_{x}\right|}$ and similarly to the above

$$
\frac{|k(w)|}{1-|k(w)|\left|w_{x}\right|} \leq \frac{K_{x}}{1-K_{x}\left|w_{x}\right|} \Longleftrightarrow|k(w)| \leq K_{x}
$$

so (8.12) is proven and we are done.

[^7]Now we restrict to $W$ being a hard-core interaction, i.e. $W(x, y) \in\{0,1\} \forall x, y \in X$ (while still being hard-core self-repulsive). If $W(x, y)=0$ (resp. 1 ), we say $x$ and $y$ are incompatible (resp. compatible) and write $x \nsim y$ (resp. $x \sim y$ ). In particular $x \nsim x$. The hypothesis of $8.5,(8.11)$, is then equivalent to

$$
\begin{aligned}
& K_{x} \geq \prod_{y \sim x} \frac{1-W(x, y) K_{y} R_{y}}{1-K_{y} R_{y}} \prod_{y \neq x: y \nsim x} \frac{1-W(x, y) K_{y} R_{y}}{1-K_{y} R_{y}}=\prod_{y \neq x: y \nsim x} \frac{1}{1-K_{y} R_{y}} \\
\Longleftrightarrow & \prod_{y \nsim x}\left(1-K_{y} R_{y}\right) \frac{K_{x} R_{x}}{1-K_{x} R_{x}}=\prod_{y \nsim x}\left(1-K_{y} R_{y}\right)\left(\frac{1}{1-K_{x} R_{x}}-1\right) \geq R_{x}
\end{aligned}
$$

Thus if $c_{x}=-\log \left(1-K_{x} R_{x}\right) \geq 0 \forall x \in X$ then (8.11) is equivalent to finding $c_{x} \geq 0$ such that

$$
\exp \left[-\sum_{y \nsim x} c_{y}\right]\left(e^{c_{x}}-1\right) \geq R_{x} \forall x \in X
$$

This is the Dobrushin condition. Since $e^{c_{x}}-1 \geq c_{x}$, a stronger and more convenient to check condition is the Kotechy-Preiss condition:

$$
\begin{equation*}
\exp \left[-\sum_{y \nsim x} c_{y}\right] c_{x} \geq R_{x} \forall x \in X \tag{8.16}
\end{equation*}
$$

Now we consider the special case where $X$ can be partitioned into $X=\bigsqcup_{n=1}^{\infty} X_{n}$ and where there are suitable $\left\{A_{n}\right\}_{n \in \mathbb{N}}$ such that

$$
\begin{equation*}
\sum_{y \in X_{n}: y \nsim x} R_{y} \leq A_{n} m \forall n \in \mathbb{N} \forall x \in X_{m} \tag{8.17}
\end{equation*}
$$

(8.17) typically arises when $X$ is a set of subsets of some $V$, and $x \nsim y \Longleftrightarrow x \cap y \neq \varnothing$. Then we take $X_{n}$ to be $\{x \in X:|x|=n\}$. Then if we showed that we can find $\left\{A_{n}\right\}$ such that

$$
\begin{equation*}
\sum_{y \in X_{n}: i \in y} R_{y} \leq A_{n} \forall n \in \mathbb{N} \forall i \in V \text {, } \tag{8.18}
\end{equation*}
$$

then for fixed $x \in X_{m}$ and $i \in x, A_{n}$ upper bounds summing $R_{y}$ over $y \in X_{n}$ with $i \in x \cap y \Longrightarrow x \nsim y$, so summing over the $m i \in x$ we can get (8.17) to hold.

Suppose (8.17) holds and suppose we take $c_{x}=e^{\alpha m} R_{x} \forall x \in X_{m}$ for some suitable $\alpha>0$. Then

$$
\begin{aligned}
\exp \left[-\sum_{y \nsim x} c_{y}\right] c_{x} & =\exp \left[-\sum_{n=1}^{\infty} \sum_{y \in X_{n}: y \nsim x} e^{\alpha n} R_{y}\right] e^{\alpha m} R_{x} \geq \exp \left[-\sum_{n=1}^{\infty} e^{\alpha n} A_{n} m\right] e^{\alpha m} R_{x} \\
& =\exp \left[\left(\alpha-\sum_{n=1}^{\infty} e^{\alpha n} A_{n}\right) m\right] R_{x}
\end{aligned}
$$

and this upper bounds $R_{x}$ if and only if

$$
\begin{equation*}
\alpha \geq \sum_{n=1}^{\infty} e^{\alpha n} A_{n} \tag{8.19}
\end{equation*}
$$

This proves the following proposition:

Proposition 8.6. Suppose that $X=\bigsqcup_{n=1}^{\infty} X_{n}$ and there exists $\left\{A_{n}\right\}_{n \in \mathbb{N}}$ and $\alpha>0$ such that

1. $\sum_{y \in X_{n}: y \nsim x} R_{y} \leq A_{n} m$ for all $m, n$ and all $x \in X_{m}$ (8.17);
2. $\sum_{n=1}^{\infty} e^{\alpha n} A_{n} \leq \alpha$ (8.19).

Then the Kotechy-Preiss condition (8.16) (and thus 8.5) holds with $c_{x}=e^{\alpha m} R_{x} \forall x \in X_{m}$.

Note that since $X$ is finite, only finitely many $A_{n}$ are nonzero. However for the rest of this section we will 'forget' this fact and instead consider the general case where $\left\{A_{n}\right\}_{n \in \mathbb{N}}$ is simply an infinite sequence.

In this case, for the existence of $\alpha>0$ such that $\sum_{n=1}^{\infty} e^{\alpha n} A_{n} \leq \alpha$ to hold, it is necessary that $A_{n}$ decays exponentially, i.e. there exists $c>1$ such that $\lim \sup _{n \rightarrow \infty} c^{n} A_{n}=M<\infty$. However, since this would imply that $\left(\frac{c}{2}\right)^{n} A_{n} \leq(M+\epsilon) / 2^{n}$ after finitely many terms, let $N$ be such that $\left(\frac{c}{2}\right)^{n} A_{n} \leq(M+\epsilon) / 2^{n} \forall n \geq N$ and $\sum_{n=N}^{\infty}(M+\epsilon) / 2^{n} \leq(\log c) / 2$, then it suffices to modify the first $N-1 A_{n}$ so that $\sum_{n=1}^{N-1} c^{n} A_{n} \leq(\log c) / 2$ for the desired $\alpha=\log c$ to exist. Thus any exponential decay on $A_{n}$ is sufficient for $\alpha$ to exist, up to modifying finitely many $A_{n}$. In some applications it is thus important to estimate $A_{n}$ for small $n$.

Let $\delta=\liminf _{n \rightarrow \infty}\left(-\log A_{n}\right) / n$ and let

$$
F(\alpha)=\alpha^{-1} \sum_{n=1}^{\infty} e^{\alpha n} A_{n}
$$

Then since $\left(e^{\alpha}\right)^{n}<A_{n}$ for large enough $n$ as long as $0<\alpha<\delta$, $e^{\alpha}$ is within the domain where $g(x)=$ $\sum_{n=1}^{\infty} x^{n} A_{n}$ is analytic, so $F$ is real-analytic in $0<\alpha<\delta$ as a composition and then product of analytic functions. Furthermore, $F(\alpha) \rightarrow \infty$ as $\alpha \rightarrow \delta$ since $e^{\left(\left(-\log A_{n}\right) / n-\epsilon\right) n} A_{n}=e^{-\epsilon n}$ can make the series grow arbitrarily large as $\epsilon \rightarrow 0$, so since the series is increasing in $\alpha$, it diverges when $\alpha \geq \delta$ so in that case $F(\alpha)=\infty$. As $\alpha \rightarrow 0$, clearly since $e^{\alpha n} A_{n} \rightarrow A_{n}$ and $\alpha^{-1} \rightarrow \infty$ we have $F(\alpha) \rightarrow \infty$, thus by IVT the infimum of $F(\alpha)$ on $0<\alpha<\delta$ is actually attained, making (8.19) equivalent to

$$
\inf _{\alpha>0} \alpha^{-1} \sum_{n=1}^{\infty} e^{\alpha n} A_{n} \leq 1
$$

### 8.4 Some combinatorial lemmas

### 8.4.1 Reduction to trees

In order to compute $w(S)$, we need to sum over all connected subgraphs $\left(S, E^{\prime}\right) \subseteq\left(S, E_{S}\right)$. However there are often 'too many' such $\left(S, E^{\prime}\right)$ to upper bound. It is thus helpful that this sum can sometimes be bounded by simply a sum over spanning trees. This is where $\left|1+v_{e}\right| \leq 1 \forall e \in E$ comes into play.

Proposition 8.7. Let $G=(V, E)$ be equipped with complex edge weights $\left\{v_{e}\right\}_{e \in E}$ satisfying $\left|1+v_{e}\right| \leq 1$ for all e. Then

$$
\begin{equation*}
\left|\sum_{\substack{E^{\prime} \subseteq E,\left(V, E^{\prime}\right) \text { connected }}} \prod_{e \in E^{\prime}} v_{e}\right| \leq \sum_{\substack{E^{\prime} \subseteq E,\left(V, E^{\prime}\right) \text { tree }}} \prod_{e \in E^{\prime}}\left|v_{e}\right| \tag{8.20}
\end{equation*}
$$

Outline of proof. WLOG $G$ is connected since otherwise both sides are 0 . It is known that there exists a map $\mathbf{R}$ from the set of spanning trees in $G$ to its set of connected subgraphs such that $T \subseteq \mathbf{R}(T)$ and that for each connected subgraph $\left(V, E^{\prime}\right) \subseteq(V, E)$, there exists a unique tree $T$ such that $E_{T} \subseteq E^{\prime} \subseteq \mathbf{E}_{R(T)}$. Thus by factoring out $\prod_{e \in E_{T}} v_{e}$ for each such $H$ on the left side of (8.20), we get

$$
\begin{aligned}
& \left|\sum_{\substack{E^{\prime} \subseteq E,\left(V, E^{\prime}\right) \text { connected }}} \prod_{e \in E^{\prime}} v_{e}\right|=\left|\sum_{\substack{E_{T} \subseteq E,\left(V, E_{T}\right) \text { tree }}} \prod_{e \in E_{T}} v_{e} \sum_{\substack{E^{\prime}: \\
E_{T} \subseteq E^{\prime} \subseteq E_{\mathbf{R}(T)}}} \prod_{e \in E^{\prime} \backslash E_{T}} v_{e}\right| \\
& =\left|\sum_{\substack{E_{T} \subseteq E,\left(V, E_{T}\right) \text { tree }}} \prod_{e \in E_{T}} v_{e} \sum_{E^{\prime}: E^{\prime} \subseteq E_{\mathbf{R}(T)} \backslash E_{T}} \prod_{e \in E^{\prime}} v_{e}\right| \\
& =\left|\sum_{\substack{E_{T} \subseteq E,\left(V, E_{T}\right) \text { tree }}} \prod_{e \in E_{T}} v_{e} \prod_{e \in E_{\mathbf{R}(T)} \backslash E_{T}}\left(1+v_{e}\right)\right| \\
& \leq \sum_{\substack{E_{T} \subseteq E,\left(V, E_{T}\right) \text { tree }}} \prod_{e \in E_{T}}\left|v_{e}\right| \prod_{e \in E_{\mathbf{R}(T)} \backslash E_{T}}\left|1+v_{e}\right| \leq \sum_{\substack{E^{\prime} \subseteq E,\left(V, E^{\prime}\right) \text { tree }}} \prod_{e \in E^{\prime}}\left|v_{e}\right| .
\end{aligned}
$$

### 8.4.2 Connected subgraphs containing a specified vertex

Definition 8.8. Let $G=(V, E)$ be finite or countably infinite equipped with edge weights $\left\{v_{e}\right\}_{e \in E}$ and let $x \in V$. Define the weighted sum (over connected subgraphs on $n$ vertices and $m$ edges)

$$
C_{n, m}\left(G,\left\{v_{e}\right\}, x\right)=\sum_{\substack{\left(V^{\prime}, E^{\prime}\right) \subseteq G \text { connected, }, e \in E^{\prime} \\ x \in V^{\prime},\left|V^{\prime}\right|=n,\left|E^{\prime}\right|=m}} \prod_{e}\left|v_{e}\right|
$$

with special cases tree sum

$$
T_{n}\left(G,\left\{v_{e}\right\}, x\right)=C_{n, n-1}\left(G,\left\{v_{e}\right\}, x\right)=\sum_{\substack{\left(V^{\prime}, E^{\prime}\right) \subseteq G \text { tree, } \\ x \in V^{\prime},\left|V^{\prime}\right|=n}} \prod_{e \in E^{\prime}}\left|v_{e}\right|
$$

and edge-counted sum

$$
C_{\bullet, m}=\sum_{n=1}^{m+1} C_{n, m}\left(G,\left\{v_{e}\right\}, x\right)
$$

When the edge weights are all 1 , we can omit them from the notation (i.e. $C_{n, m}(G, x)=C_{n, m}(G, 1, \ldots, 1, x)$ etc.). For the infinite $r$-regular tree $\mathbf{T}_{r}$ containing vertex $y$, define the constant

$$
t_{n}^{(r)}=T_{n}\left(\mathbf{T}_{r}, y\right)
$$

Note that we trivially have

$$
\begin{equation*}
C_{n, m}\left(G,\left\{v_{e}\right\}, x\right) \leq C_{n, m}(G, x)\left(\sup _{e \in E}\left|v_{e}\right|\right)^{m} \tag{8.21}
\end{equation*}
$$

We actually as of now already have all we need to state and prove the main result. However, in this section we will first state some results that will improve our main theorem. The following was proven using the universal cover graph ${ }^{13}$ for the first inequality, and generating functions and Lagrange's implicit function theorem for the closed form of $t_{n}^{(r)}$, but it is stated here without proof.

Proposition 8.9. Let $G=(V, E)$ be finite or countably infinite of maximum degree $r$ and equipped with edge weights $\left\{v_{e}\right\}_{e \in E}$. Let $x \in V$ and $y$ be a vertex in $\mathbf{T}_{r}$. Then

$$
C_{\bullet, m}(G, x) \leq C_{\bullet, m}\left(\mathbf{T}_{r}, y\right)=T_{m+1}\left(\mathbf{T}_{r}, y\right)=t_{m+1}^{(r)}=r \frac{((r-1)(m+1))!}{m!((r-2) m+3)!}
$$

In consequence (8.21)

$$
C_{\bullet, m}\left(G,\left\{v_{e}\right\}, x\right) \leq t_{m+1}^{(r)}\left(\sup _{e \in E}\left|v_{e}\right|\right)^{m}
$$

a particular case of which is

$$
\begin{equation*}
T_{n}\left(G,\left\{v_{e}\right\}, x\right) \leq t_{n}^{(r)}\left(\sup _{e \in E}\left|v_{e}\right|\right)^{n-1} \tag{8.22}
\end{equation*}
$$

Another result, whose proof will be given, helps bound the expression 8.9 gives for $t_{n}^{(r)}$.

## Proposition 8.10.

$$
t_{n}^{(r)}=r \frac{((r-1)(n))!}{(n-1)!((r-2) n+2)!} \leq \frac{(r n)^{n-1}}{n!}
$$

Proof.

$$
\frac{((r-1)(n))!}{((r-2) n+2)!}=\frac{(r n-n)!}{(r n-n-(n-2))!} \leq(r n-n)^{n-2} \leq(r n)^{n-2}
$$

### 8.5 The main results

This is the main theorem of the paper:

Theorem 8.11. Let $G=(V, E)$ be loopless, finite and equipped with complex edge weights satisfying $\left\{v_{e}\right\}_{e \in E}$ satisfying $\left|1+v_{e}\right| \leq 1 \forall e \in E$. Let $Q=Q\left(G,\left\{v_{e}\right\}\right)>0$ be the smallest number satisfying

$$
\begin{equation*}
\inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n} Q^{-(n-1)} \max _{x \in V} T_{n}\left(G,\left\{v_{e}\right\}, x\right) \leq 1 \tag{8.23}
\end{equation*}
$$

[^8]Then all zeros (in s) of $Z_{G}\left(s,\left\{v_{e}\right\}\right)$ lie in $|s|<Q$.

Note that $Q<\infty$ since $T_{n} \equiv 0$ for $n>|V|$, so the series is finite.
Proof. First by 8.3 we view $Z_{G}$ as a grand-canonical gas with $w_{S}=w(S)$ and $W\left(S_{i}, S_{j}\right)$ being the exclusion function $1\left\{S_{i} \cap S_{j}=\varnothing\right\}$. Let $R_{S}=w(S)$ and let $A_{n}=\max _{x \in V} \sum_{S: x \in S,|S|=n}|w(S)|$. Note $A_{1}=0$ which is why the series in (8.23) starts at 2. Then using the partition of $\mathscr{P}_{*}(V)$ into $X_{n}=\left\{S \in \mathscr{P}_{*}(V):|S|=n\right\}$, we have

$$
\sum_{S \in X_{m}: x \in S} R_{S}=\sum_{S:|S|=n, x \in S} R_{S} \leq A_{n} \forall x \in V
$$

i.e. (8.18) and in consequence the first condition of 8.6 hold.

Note that from 8.7 we have, for all $S$ where $x \in S$ and $|S|=n$,

$$
\begin{align*}
& |w(S)| \leq|s|^{-(|S|-1)} \sum_{\substack{E^{\prime} \subseteq E_{S},\left(S, E^{\prime}\right) \text { tree }}} \prod_{e \in E^{\prime}}\left|v_{e}\right| \\
\Longrightarrow & \sum_{S: x \in S,|S|=n}|w(S)| \leq|s|^{-(n-1)} \sum_{S: x \in S,|S|=n} \sum_{\substack{E^{\prime} \subseteq E_{S},\left(S, E^{\prime}\right) \text { tree }}} \prod_{e \in E^{\prime}}\left|v_{e}\right|=|s|^{-(n-1)} T_{n}\left(G,\left\{v_{e}\right\}, x\right) \\
\Longrightarrow & A_{n}=\max _{x \in V} \sum_{S: x \in S,|S|=n}|w(S)| \leq|s|^{-(n-1)} \max _{x \in V} T_{n}\left(G,\left\{v_{e}\right\}, x\right) \\
\Longrightarrow & \alpha^{-1} \sum_{n=1}^{\infty} e^{\alpha n} A_{n}=\alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n} A_{n} \leq \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n}|s|^{-(n-1)} \max _{x \in V} T_{n}\left(G,\left\{v_{e}\right\}, x\right) \\
\Longrightarrow & \inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n} A_{n} \leq \inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n}|s|^{-(n-1)} \max _{x \in V} T_{n}\left(G,\left\{v_{e}\right\}, x\right) \tag{8.24}
\end{align*}
$$

Since the right hand side of (8.24) is decreasing in $|s|$ and the second condition of 8.6 would be satisfied with $|s|=Q$, it would also be satisfied with $|s|>Q$ so by $8.5 Z_{G}\left(s,\left\{v_{e}\right\}\right)$ is nonvanishing in $w_{S} \leq R_{S}=w(S)$, thus $Z_{G}\left(s,\left\{v_{e}\right\}\right) \neq 0 \forall|s|>Q$.

If, in addition, the degree of vertices in $G$ is at most $r$, then let $C(r)$ be the smallest number such that

$$
\inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n}(C(r))^{-(n-1)} t_{n}^{(r)} \leq 1
$$

Then as a corollary to 8.11 , directly from 8.9 , we have

Corollary 8.12. Let $G=(V, E)$ be loopless, finite, of maximum degree $r$ and equipped with complex edge weights satisfying $\left\{v_{e}\right\}_{e \in E}$ satisfying $\left|1+v_{e}\right| \leq 1 \forall e \in E$. Let $v_{\max }=\max _{e \in E}\left|v_{e}\right|$. Then all zeros of $Z_{G}\left(s,\left\{v_{e}\right\}\right)$ lie in $|s|<C(r) v_{\max }$.

Recall that the chromatic polynomial $C_{G}(s)=Z_{G}(s,-1, \ldots,-1)$. Thus directly from 8.12 we have

Corollary 8.13. Let $G=(V, E)$ be loopless, finite, of maximum degree $r$. Then all zeros of $C_{G}(s)$ lie in $|s|<C(r)$.

The following is a corollary of the bound on $t_{n}^{(r)}$ from 8.10.

## Corollary 8.14.

$$
C(r)=O(r)
$$

Proof. Note $\log \frac{n^{n-1}}{n!}=(n-1) \log n-\sum_{i=1}^{n} \log n \leq(n-1) \log n-n \log n+n=n-\log n$ so $\frac{n^{n-1}}{n!} \leq \frac{e^{n}}{n} \leq e^{n}$. Thus if $C(r) \geq M r$,

$$
\begin{aligned}
\inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n}(C(r))^{-(n-1)} t_{n}^{(r)} & \leq \inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n}(C(r))^{-(n-1)} \frac{(r n)^{n-1}}{n!} \\
& =\inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{\alpha n}\left(\frac{C(r)}{r}\right)^{-(n-1)} \frac{n^{n-1}}{n!} \\
& \leq \inf _{\alpha>0} \alpha^{-1} \sum_{n=2}^{\infty} e^{(\alpha+1) n}\left(\frac{C(r)}{r}\right)^{-(n-1)} \\
& \leq \inf _{\alpha>0} \alpha^{-1} M \sum_{n=2}^{\infty}\left(\frac{e^{\alpha+1}}{M}\right)^{n}=\inf _{\alpha>0}(M \alpha)^{-1} e^{2(\alpha+1)} \frac{1}{1-e^{\alpha+1} / M} .
\end{aligned}
$$

The above can be verified with $M=45$ and $\alpha=0.5$.

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[^1]:    ${ }^{2} \mathrm{We}$ always assume the vertices to be labeled, i.e. distinct, so $G_{1}$ and $G_{2}$ could be isomorphic but not equal.

[^2]:    ${ }^{3}$ By at most $\left(k(k-1)^{r / 2-1}+1\right)^{2}$ since $u \in V\left(H_{r}(v)\right) \Longleftrightarrow v \in V\left(H_{r}(u)\right)$ and each $H_{r}(v)$ has at most $k(k-1)^{r / 2-1}+1$ nodes, so $\left|v: V\left(H_{r}(v)\right) \cap V\left(H_{r}\left(v_{0}\right)\right) \neq \varnothing\right|=\left|\left\{v: v \in V\left(H_{r}(u)\right): u \in V\left(H_{r}\left(v_{0}\right)\right)\right\}\right| \leq\left(k(k-1)^{r / 2-1}+1\right)^{2}$.

[^3]:    ${ }^{4} n_{1}=k+1$ and $n_{i+1}=n_{i}+2$ if $k$ is odd; otherwise $n_{i+1}=n_{i}+1$
    ${ }^{5}$ Note $(a+\epsilon)^{k}$ (where $\epsilon \rightarrow 0$ and $a \geq 0$ constant) is upper bounded by $a^{k}+k^{2} a^{k-1} \epsilon$, since the $i$-th term after $k a^{k-1} \epsilon$ in the binomial expansion is exactly $\frac{\epsilon^{i}\binom{k}{i+1}}{k a^{i}} \xrightarrow{\epsilon \rightarrow 0} 0$ multiplied onto $k a^{k-1} \epsilon . a=1-\frac{t}{n_{i}-k+1} \leq 1$ and $\epsilon=\frac{t}{n_{i}-k+1}$.

[^4]:    ${ }^{6}$ The $s=2$ case is known as the Ising model.
    ${ }^{7} \beta=1 /\left(k_{B} T\right)$ where $k_{B}$ is the Boltzmann constant and $T$ is the temperature.

[^5]:    ${ }^{8 \prime}$ nondegenerate' excludes cases $v=-1$ (using (8.1)), which causes some $\sigma$ to have weight 0 , and $s=0$ (using (8.2)); 'physical' means that $v$ is such that there are no negative terms in the outmost summation, allowing a probability distribution to exist
    ${ }^{9}$ In particular, theorems that guarantee a certain complex region is free of zeros are known as Lee-Yang theorems.
    ${ }^{10} \mathrm{~A}$ similar argument can be made for parallel edges $e_{1}, \ldots, e_{n}$, which can be replaced by a single edge $e$ with weight $\prod_{i=1}^{n}\left(1+v_{e_{i}}\right)-1$ and get the same $Z_{G}$. However this does not simplify much so we will not be assuming there be no parallel edges.

[^6]:    ${ }^{11}$ Note that this is equivalent to setting $w_{x}=0$ for $x \notin \Lambda$.

[^7]:    ${ }^{12} \log z$ denotes the principal value of $z \neq 0$, i.e. the $\operatorname{logarithm}$ where $\operatorname{Im} \log z \in[-\pi, \pi]$, i.e. $\log (x+y i)=\log \sqrt{x^{2}+y^{2}}+$ $i \operatorname{atan} 2(x, y) \forall x, y \in \mathbb{R}, x y \neq 0$.

[^8]:    ${ }^{13}$ The universal cover graph $U$ is a tree constructed from the set of all walks starting at $x$ and that do not using any edge consecutively. $U$ is a subgraph of $\mathbf{T}_{r}$ so $C_{\bullet}, m(U, x) \leq C_{\bullet}, m\left(\mathbf{T}_{r}, y\right)$ is clear. However proving $C_{\bullet}, m(G, x) \leq C_{\bullet}, m(U, x)$ requires more work and is omitted.

