Abstract. In a classical work of the 1950’s, Lee and Yang proved that for fixed nonnegative temperature, the zeros of the partition functions of a ferromagnetic Ising model always lie on the unit circle in the complex magnetic field. Zeros of the partition function in the complex temperature were then considered by Fisher, when the magnetic field is set to zero. Limiting distributions of Lee-Yang and of Fisher zeros are physically important as they control phase transitions in the model. One can also consider the zeros of the partition function simultaneously in both complex magnetic field and complex temperature. They form an algebraic curve called the Lee-Yang-Fisher (LYF) zeros. In this paper we study their limiting distribution for the Diamond Hierarchical Lattice (DHL). In this case, it can be described in terms of the dynamics of an explicit rational function $R$ in two variables (the Migdal-Kadanoff renormalization transformation). We prove that the Lee-Yang-Fisher zeros are equidistributed with respect to a dynamical $(1,1)$-current in the projective space. The free energy of the lattice gets interpreted as the pluripotential of this current. We also describe some of the properties of the Fatou and Julia sets of the renormalization transformation.

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1. Introduction

1.1. Phenomenology of Lee-Yang-Fisher zeros. We will begin with a brief survey of the Ising model of magnetic matter. More background is given in Part I of the series [BLR1]; see also [Ba, R2].

The matter in a certain scale is represented by a graph $\Gamma$. Let $V$ and $E$ stand respectively for the set of its vertices (representing atoms) and edges (representing magnetic bonds between the atoms). A magnetic state of the matter is represented by a spin configuration $\sigma \colon V \to \{\pm 1\}$ on $\Gamma$. The spin $\sigma(v)$ represents a magnetic momentum of an atom $v \in V$.

The total energy of such a configuration is given by the Hamiltonian
\begin{equation}
H(\sigma) = -J \sum_{\{v,w\} \in E} \sigma(v)\sigma(w) - h M(\sigma),
\end{equation}
where $M(\sigma) = \sum_{v \in V} \sigma(v)$.

Here, $M(\sigma)$ is the total magnetic moment of the configuration $\sigma$, $J$ is a constant describing the interaction between neighboring spins and $h$ is a constant describing an external magnetic field applied to the matter.

By the Gibbs Principle, a configuration $\sigma$ occurs with probability proportional to its Gibbs weight $W(\sigma) = \exp(-H(\sigma)/T)$, where $T \geq 0$ is the temperature. The total Gibbs weight $Z = \sum W(\sigma)$ is called the partition function. It is a Laurent polynomial in two variables $(z, t)$, where $z = e^{-h/T}$ is a "field-like" variable and $t = e^{-J/T}$ is "temperature-like".

For a fixed $t \in [0, 1]$, the complex zeros of $Z(z, t)$ in $z$ are called the Lee-Yang zeros. Their role comes from the fact that some important observable quantities can be calculated as electrostatic-like potentials of the equally charged particles located at the Lee-Yang zeros. A celebrated Lee-Yang Theorem [YL, LY] asserts that for the ferromagnetic Ising model on any graph, the zeros of the partition function lie on the unit circle $T$ in the complex plane (corresponding to purely imaginary magnetic field $h = -iT\phi$).

Magnetic matter in various scales can be modeled by a hierarchy of graphs $\Gamma_n$ of increasing size (corresponding to finer and finer scales). For suitable models, zeros of the partition functions $Z_n$ will have an asymptotic distribution $d\mu_t = \rho_t d\phi/2\pi$ on the unit circle. This distribution supports singularities of the magnetic observables (or rather, their thermodynamical limits), and hence it captures phase transitions in the model.

Instead of freezing temperature $T$, one can freeze the external field $h$, and study zeros of $Z(z, t)$ in the $t$-variable. They are called Fisher zeros as they were first studied by Fisher for the regular two-dimensional lattice, see [F, BK]. Similarly to the Lee-Yang zeros, asymptotic distribution of the Fisher zeros is supported on the singularities of the magnetic observables, and is thus related to phase transitions in the model. However, Fisher zeros do not lie on the unit circle any more. For instance, for the regular 2D lattice at $h = 0$, the asymptotic distribution lies on the union of two Fisher circles, see Figure 1.1.

More generally, one can study the distribution of zeros for $Z(z, t)$ on other complex lines in $\mathbb{C}^2$. In order to organize the limiting distributions over all such lines into a single object, we use the theory of currents; see [dR, Le].

\footnote{i.e., with $J > 0$}
\footnote{We will use either the $z$-coordinate or the angular coordinate $\phi = \arg z \in \mathbb{R}/2\pi\mathbb{Z}$ on $T$, without a comment.}
A (1, 1)-current \( \nu \) on \( \mathbb{C}^2 \) is a linear functional on the space of (1, 1)-forms that have compact support (see Appendix A.3). A basic example is the current \([X]\) of integration over an algebraic curve \( X \). A plurisubharmonic function \( G \) is called a pluripotential of \( \nu \) if \( \frac{i}{2} \partial \bar{\partial} G = \nu \), in the sense of distributions. (Informally, this means that \( \Delta(G| L) = \nu| L \) for almost any complex line \( L \), so \( G| L \) is an electrostatic potential of the charge distribution \( \nu| L \).)

For each \( n \), the zero locus of \( Z_n(z, t) \) is an algebraic curve \( S_n^c \subset \mathbb{C}^2 \), which we call the Lee-Yang-Fisher (LYF) zeros. Let \( d_n \) be the degree of \( S_n^c \). It is natural to ask whether there exists a (1, 1) current \( \mu^c \) so that

\[
\frac{1}{d_n} [S_n^c] \rightarrow \mu^c.
\]

It will describe the limiting distribution of Lee-Yang-Fisher zeros. Within almost any complex line \( L \), the limiting distribution of zeros can be obtained as the restriction \( \mu^c| L \).

In order to justify existence of \( \mu^c \), one considers the sequence of “free energies”

\[
F_n^\#(z, t) := \log |\hat{Z}_n(z, t)|,
\]

where \( \hat{Z}_n(z, t) \) is the polynomial obtained by clearing the denominators of \( Z_n \). We will say that the sequence of graphs \( \Gamma_n \) has a global thermodynamic limit if

\[
\frac{1}{d_n} F_n^\#(z, t) \rightarrow F^\#(z, t)
\]

in \( L^1_{\text{loc}}(\mathbb{C}^2) \). In Proposition 2.1 we will show that this is sufficient for the limiting current \( \mu^c \) to exist. More details, including the relation of this notion to the classical definition of thermodynamic limit, are given in §2.2.

The support of \( \mu^c \) consists of the singularities of the magnetic observables of the model, thus describing global phase transitions in \( \mathbb{C}^2 \). Connected components of \( \mathbb{C}^2 \setminus \text{supp} \mu^c \) describe the distinct (complex) phases of the system.

This discussion can also be extended to the compactification \( \mathbb{CP}^2 \) of \( \mathbb{C}^2 \) by considering lifts \( F_n^\# \) of the free energies to \( \mathbb{C}^3 \); see §2.3.

1.2. **Diamond hierarchical model.** The Ising model on hierarchical lattices was introduced by Berker and Ostlund [BO] and further studied by Bleher & Žalys [BZ1, BZ2, BZ3] and Kaufman & Griffiths [KG1].
The diamond hierarchical lattice (DHL), illustrated on Figure 1.2, is a sequence of graphs $\Gamma_n$ with two marked vertices such that $\Gamma_0$ is an interval, $\Gamma_1 = \Gamma$ is a diamond, and $\Gamma_{n+1}$ is obtained from $\Gamma_n$ by replacing each edge of $\Gamma_n$ with $\Gamma$ so that the marked vertices of $\Gamma$ match with the vertices of $\Gamma_n$. We then mark two vertices in $\Gamma_{n+1}$ so that they match with the two marked vertices of $\Gamma_n$. Part I of the series [BLR1] and the present paper are both fully devoted to study of this lattice.

**Remark 1.1.** For the DHL we will use the following slightly different definition of magnetic momentum:

$$M(\sigma) = \frac{1}{2} \sum_{(v,w) \in \mathcal{E}} (\sigma(v) + \sigma(w)).$$

For a motivation, see Appendix F of [BLR1]. Also, we will use $t := t^2 = e^{-2J/T}$ for the temperature-like variable, as it makes formulas nicer.

There is a general physical principle that the values of physical quantities depend on the scale where the measurement is taken. The corresponding quantities are called renormalized, and the (semi-group) of transformations relating them at various scales is called the renorm-group (RG). However, it is usually hard to justify rigorously the existence of the RG, let alone to find exact formulas for RG transformations. The beauty of hierarchical models is that all this can actually be accomplished.

The Migdal-Kadanoff RG equations [M1, M2, K] (see also [BO, BZ1, KG1]) for the DHL assume the form:

$$\left(z_{n+1}, t_{n+1}\right) = \left(\frac{z_n^2 + t_n^2}{z_n^2 + t_n^2}, \frac{z_n^2 + t_n^2 + 2}{z_n^2 + t_n^2} \right) := R(z_n, t_n),$$

where $z_n$ and $t_n$ are the renormalized field-like and temperature-like variables on $\Gamma_n$. The map $R$ that relates these quantities is also called the renormalization transformation.

To study the Fisher zeros, we consider the line $\mathcal{L}_{\text{inv}} = \{ z = 1 \}$ in $\mathbb{CP}^2$. This line is invariant under $R$, and $R : \mathcal{L}_{\text{inv}} \to \mathcal{L}_{\text{inv}}$ reduces to a fairly simple one-dimensional rational map

$$R : t \mapsto \left(\frac{2t}{t^2 + 1}\right)^2.$$
The Fisher zeros at level $n$ are obtained by pulling back the point $t = -1$ under $R^n$. As shown in [BL], the limiting distribution of the Fisher zeros in this case exists and it coincides with the measure of maximal entropy (see [Br, Ly]) of $R|_L$. The limiting support for this measure is the Julia set for $R|_L^{inv}$, which is shown in Figure 1.3. It was studied by [DDI, DIL, BL, Ish] and others.

![Figure 1.3](image.png)

**Figure 1.3.** On the left is the Julia set for $R|_L^{inv}$. On the right is a zoomed-in view of a boxed region around the critical point $t_c$. The invariant interval $[0, 1]$ corresponds to the states with real temperatures $T \in [0, \infty]$ and vanishing field $h = 0$.

To study the Lee-Yang zeros for the DHL, one considers the **Lee-Yang cylinder** $C := \mathbb{T} \times I$, which is invariant under $R$. The Lee-Yang zeros for $\Gamma_n$ are the real algebraic curve $S_n := S_0 \cap C$. Equation (1.3) shows that $S_n$ is the pullback of $S_0$ under the $n$-fold iterate of $R$, i.e., $S_n = (R^n)^* S_0$. Part I of the series [BLR1] uses this dynamical approach to obtain detailed information about the limiting distribution of Lee-Yang zeros for the DHL.

In this paper, we will use the Migdal-Kadanoff RG equations to study the global limiting distribution of Lee-Yang-Fisher zeros for the DHL. Their extension to $\mathbb{C}P^2$, which we will also denote by $S_c$, is a curve of degree $2 \cdot 4^n$. Our main result is:

**Theorem (Global Lee-Yang-Fisher Current).** *The currents $\frac{1}{2\pi} \omega[S_n^c]$ converge distributionally to some $(1, 1)$-current $\mu^c$ on $\mathbb{C}P^2$ whose pluripotential coincides with the (lifted) free energy $\hat{F}^\#$ of the system.*

An important subtlety arises because the degrees of $R$ do not behave properly under iteration:

$$4^n < \deg(R^n) < (\deg(R))^n = 6^n.$$  

This *algebraic instability* of $R$ has the consequence that 

$$S_n^c \neq (R^n)^* S_0^c.$$  

The issue is resolved by working with another rational mapping $R$ coming directly from the Migdal-Kadanoff RG Equations, without passing to the “physical” $(z, t)$-coordinates. This map is semi-conjugate to $R$ by a degree two rational map $\Psi : \mathbb{C}P^2 \rightarrow \mathbb{C}P^2$. Moreover, $R$ is algebraically stable, satisfying

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3For the definition, see Appendix [BLR].
deg(R^n) = (deg(R))^n = 4^n. For each n ≥ 0, we have:

\[ S_c^n = \Psi^{-1}(R^{-n}S_c^0), \]

where \( S_c^0 \) is an appropriate projective line.

Associated to any (dominant, algebraically stable) rational mapping \( f : \mathbb{CP}^k \to \mathbb{CP}^k \) is a canonically defined invariant current \( m_c \), called the Green current of \( f \). It satisfies \( f^* m_c = d \cdot m_c \), where \( d = \text{deg} f \). If \( f \) satisfies an additional (minor) technical hypothesis, then \( \text{supp} m_c \) equals the Julia set of \( f \). (In our case, \( R \) does not satisfy this additional hypothesis, but we still have that \( \text{supp} m_c \) equals the Julia set for \( R \); see Prop. 4.4).

Such invariant currents are a powerful tool of higher-dimensional holomorphic dynamics: see Bedford-Smillie [BS], Fornaess-Sibony [FS1], and others (see [Si] for an introductory survey to this subject). In Appendix B we provide some further background on complex dynamics in several variables.

We will show that the normalized sequence of currents \( \frac{1}{4^n}(R^n)^* [S_c^0] \) converges distributionally to the Green current \( m_c \) of \( R \). Pulling everything back under \( \Psi \), this will imply the Global Lee-Yang-Fisher Theorem. In this way, the classical Lee-Yang-Fisher theory gets linked to the contemporary dynamical pluripotential theory.

Asymptotic distribution for pullbacks of algebraic curves under rational maps has been a focus of intense research in multidimensional holomorphic dynamics since the early 1990’s; see [BS, FS1, RuSh, FaJ, DS2, DDG] for a sample of papers on the subject. Our result above is very close in spirit to this work. However, the above theorem does not seem to be a consequence of any available results.

1.3. Structure of the paper. We begin in §2 by recalling the definitions of free energy and the classical notion of thermodynamic limit for the Ising model. We then discuss the notion of global thermodynamic limit, which is sufficient in order to guarantee that some lattice have a \((1,1)\)-current \( \mu_c \) describing its limiting distribution of LYF zeros in \( \mathbb{C}^2 \). Following this, we extend the whole description from \( \mathbb{C}^2 \) to \( \mathbb{CP}^2 \) by lifting the partition function and free energy to \( \mathbb{C}^3 \). We also give an alternative interpretation of the partition function as a section of (an appropriate tensor power of) the co-tautological line bundle over \( \mathbb{CP}^2 \) that will be central to the proof of the Global LYF Current Theorem. We conclude §2 by summarizing material on the Migdal-Kadanoff RG equations.

In §3 we summarize the global features of the mappings \( R \) and \( R \) on the complex projective space \( \mathbb{CP}^2 \) that were studied in [BLR1], including their critical and indeterminacy loci, superattracting fixed points and their separatrices.

In the next section, §4 we define the Fatou and Julia sets for \( R \) and show that the Julia set coincides with the closure of preimages of the invariant complex line \( \{z = 1\} \) (corresponding to the vanishing external field). It is based on M. Green’s criteria for Kobayashi hyperbolicity of the complements of several algebraic curves in \( \mathbb{CP}^2 \) [G1, G2] that generalize the classical Montel Theorem. We then use this result to prove that points in the interior of the solid cylinder \( \mathbb{D} \times I \) are attracted to a superattracting fixed point \( \eta = (0,1) \) of \( R \).

We prove the Global LYF Current Theorem in §5 relying on estimates of how volume is transformed under a single iterate of \( R \) that are completed in §6.
We finish with several Appendices. In Appendix A we collect needed background in complex geometry (normality, Kobayashi hyperbolicity, currents and their pluri-potentials, line bundles over \(\mathbb{CP}^2\), etc). In Appendix B we provide background on the complex dynamics in higher dimensions, including the notion of algebraic stability and information on the Green current. In Appendix C we provide background on Whitney Folds, a normal form for the simplest critical points of our mapping. In Appendix D we collect several open problems.

1.4. Basic notation and terminology. \(\mathbb{C}^* = \mathbb{C} \setminus \{0\}\), \(T = \{|z| = 1\}\), \(D_r = \{|z| < r\}\), \(D = D_1\), and \(\mathbb{N} = \{0, 1, 2, \ldots\}\). Given two variables \(x\) and \(y\), \(x \asymp y\) means that \(c \leq |x/y| \leq C\) for some constants \(C > c > 0\).

Acknowledgment. We thank Jeffrey Diller, Han Peters, Robert Shrock, and Dror Varolin for interesting discussions and comments. The work of the first author is supported in part by the NSF grants DMS-0652005 and DMS-0969254. The work of the second author has been partially supported by NSF, NSERC and CRC funds. The work of the third author was partially supported by startup funds from the Department of Mathematics at IUPUI.

2. Description of the model

2.1. Free Energy. A configuration of spins \(\sigma : V \to \{\pm 1\}\) has energy given by the Hamiltonian

\[
H(\sigma) = -J \sum_{\{v, w\} \in E} \sigma(v)\sigma(w) - hM(\sigma), \quad \text{where} \quad M(\sigma) = \frac{1}{2} \sum_{(v, w) \in E} (\sigma(v) + \sigma(w)).
\]

The Ising model is called ferromagnetic if \(J > 0\), and anti-ferromagnetic otherwise. In this paper (and Part I) we will assume that the model is ferromagnetic.

The partition function (or the statistical sum) is the total Gibbs weight of all of the configurations of spins:

\[
Z_{\Gamma} = Z_{\Gamma}(z, t) = \sum_{\sigma} W(\sigma) = \sum_{\sigma} \exp(-H(\sigma)/T).
\]

It is a Laurent polynomial in \(z\) and \(t\). Notice that each of the Gibbs weights is unchanged under \((\sigma, h) \mapsto (-\sigma, -h)\). This basic symmetry of the Ising model makes \(Z\) invariant under the involution \(\iota : (z, t) \mapsto (z^{-1}, t)\). Therefore, it has the form

\[
Z_{\Gamma} = \sum_{n=0}^{d} a_n(t)(z^n + z^{-n}), \quad \text{where} \quad d = |E|.
\]

Moreover, \(a_d = t^{-d/2}\). Thus, for any given \(t \in \mathbb{C}^*\), \(Z_{\Gamma}(t, z)\) has \(2|E|\) roots \(z_i(t) \in \mathbb{C}\) called the Lee-Yang zeros.

A configuration \(\sigma\) occurs with probability \(P(\sigma) = W(\sigma)/Z\). Its entropy is defined as \(S(\sigma) = -\log P(\sigma) = \log Z + H(\sigma)/T\). The free energy of the system is defined as

\[
F_{\Gamma} = H(\sigma) - TS(\sigma) = -T \log |Z_{\Gamma}|.
\]

It is independent of the configuration \(\sigma\).

Equations (2.3) and (2.2) imply:

\[
F_{\Gamma} = -T \sum \log |z - z_i(t)| + |E| T \left(\log |z| + \frac{1}{2} \log |t|\right),
\]
where the summation is taken over the $2|\mathcal{E}|$ Lee-Yang zeros $z_i(t)$ of $Z(\cdot, t)$.

**Remark 2.1.** It will be useful to consider the following variant of the free energy:

$$F^#(z, t) := -\frac{1}{T} F(\Gamma(z, t)) + d \log |z| + \frac{1}{2} d \log |t| = \log |\tilde{Z}(z, t)|.$$  

Here, $\tilde{Z}(z, t) := z^d t^{d/2} Z(z, t)$ is obtained by clearing the denominators of $Z$. The advantage of using $F^#$, instead of $F$, is that it extends as a plurisubharmonic function on all of $\mathbb{C}^2$. We will also refer to $F^#$ as the “free energy”.

### 2.2. Thermodynamic limit

Assume that we have a “lattice” given by a “hierarchy” of graphs $\Gamma_n$ of increasing size with partition functions $Z_n$, free energies $F_n$ and magnetizations $M_n$. To pass to the thermodynamic limit we normalize these quantities per bond. One says that the hierarchy of graphs has a thermodynamic limit if

$$\lim_{n \to \infty} \frac{1}{|\mathcal{E}_n|} F_n(z, t) = F(z, t) \quad \text{for any } z \in \mathbb{R}_+, \quad t \in (0, 1).$$  

In this case, the function $F$ is called the free energy of the lattice. For many lattices (e.g. $\mathbb{Z}^d$), existence of the thermodynamic limit can be justified by van Hove’s Theorem [vH, R2]. If the classical thermodynamic limit exists, then one can justify existence of the limiting distribution of Lee-Yang zeros and relate it to the limiting free energy; see [BLR1, Prop. 2.2].

When considering Fisher zeros, it is more convenient to work with the variant $F^#$ of the free energy that is defined in Remark 2.1. In order to prove existence of a limiting distribution for the Fisher zeros, one needs to prove existence of a limiting free energy:

$$\lim_{n \to \infty} \frac{1}{2|\mathcal{E}_n|} F^#_n(1, t) = F^#(1, t) \quad \text{in } L^1_{\text{loc}}(\mathbb{C}).$$  

For the $\mathbb{Z}^2$ lattice this is achieved by the Onsager solution, which provides an explicit formula the limiting free energy; See, for example, [Ba]. Similar techniques apply to the triangular, hexagonal, and various homopolygonal lattices (see [MaSh1, MaSh2] for suitable references and an investigation of the distribution of Fisher zeros for these lattices). For various hierarchical lattices, (2.6) can be proved by dynamical means. It seems to be an open question for other lattices, including $\mathbb{Z}^d$, when $d \geq 3$.

The situation is similar for the Lee-Yang-Fisher zeros:

**Proposition 2.1.** Let $\Gamma_n$ be a hierarchy of graphs and suppose that

$$\lim_{n \to \infty} \frac{1}{2|\mathcal{E}_n|} F^#_n(z, t) = F^#(z, t) \quad \text{in } L^1_{\text{loc}}(\mathbb{C}^2).$$

Then, there is a closed positive $(1,1)$-current $\mu^c$ on $\mathbb{C}^2$ describing the limiting distribution of Lee-Yang-Fisher zeros. Its pluripotential coincides with the free energy $F^#(z, t)$.

For the DHL, we will prove existence of the limit (2.7) in the Global LYFC Current Theorem. It is an open question whether such a limit exists for other lattices, including the classical $\mathbb{Z}^d$ lattices for $d \geq 2$; see Problem D.1.

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\[\text{Note that the DHL is not in this class—instead, dynamical techniques are used to justify its classical thermodynamic limit.}\]
Proof. The locus of Lee-Yang-Fisher zeros $S_n^c$ are the zero set (counted with multiplicities) of the degree $2|E_n|$ polynomial $\tilde{Z}_n(z,t)$. The Poincaré-Lelong Formula describes its current of integration:

$$[S_n^c] = \Delta_p \log |\tilde{Z}_n(z,t)| = \Delta_p F_n^\#(z,t).$$

Here, $\Delta_p = \frac{i}{\pi} \partial \bar{\partial}$ is the pluri-Laplacian; see Appendix A.

Hypothesis (2.7) implies

$$\frac{1}{2|E_n|}[S_n^c] = \Delta_p \frac{1}{2|E_n|} F_n^\#(z,t) \to \Delta_p F_n^\#(z,t) =: \mu^c.$$

□

2.3. Global consideration of partition functions and free energy on $\mathbb{CP}^2$.

It will be convenient for us to extend the partition functions $Z_n$ and their associated free energies $F_n^\#$ from $\mathbb{C}^2$ to $\mathbb{CP}^2$. We will use the homogeneous coordinates $[Z : T : Y]$ on $\mathbb{CP}^2$, with the “physical” copy of $\mathbb{C}^2$ given by the affine coordinates $(z,t) \mapsto [z : t : 1]$.

For each $n$, we clear the denominators of $Z_n(z,t)$, obtaining a polynomial $\tilde{Z}_n(z,t)$ of degree $d_n := 2|E_n|$. It lifts to a unique homogeneous polynomial $\hat{Z}_n(Z,T,Y)$ of the same degree that satisfies $\tilde{Z}_n(z,t) = \hat{Z}_n(z,t,1)$. The associated free energy becomes a plurisubharmonic function

$$\hat{F}_n^\#(Z,T,Y) := \log |\hat{Z}_n(Z,T,Y)|$$
on $\mathbb{C}^3$. It is related by the Poincaré-Lelong Formula to the current of integration over the Lee-Yang-Fisher zeros: $\pi^*[S_n^c] := \Delta_p \hat{F}_n^\#(Z,T,Y)$.

Both of these extensions are defined on $\mathbb{C}^3$, rather than $\mathbb{CP}^2$. In the proof of the Global LYF Current Theorem, it will be useful for us to interpret the partition function as an object defined on $\mathbb{CP}^2$. Instead of being a function on $\mathbb{CP}^2$, it gets interpreted as a section $\sigma_{\hat{Z}_n}$ of an appropriate tensor power of the co-tautological line bundle; See Appendix A.4. The Lee-Yang-Fisher zeros $S_n^c$ are described as the zero locus of this section. (See Remark A.1 for an explanation of why a similar interpretation of the free energy as a section of a suitable real line bundle is not used.)

2.4. Migdal-Kadanoff renormalization for the DHL. Migdal-Kadanoff renormalization will allow us to write recursive formulæ for the partition function $Z_n(z,t)$. Restricting the spins at the marked vertices $\{a, b\}$ we obtain three conditional partition functions, $U_n$, $V_n$ and $W_n$ as follows:

$$U_n := Z_n\left(\begin{array}{c} \oplus \\
\end{array}\right), \quad V_n := Z_n\left(\begin{array}{c} \oplus \\
\end{array}\right) = Z_n\left(\begin{array}{c} \oplus \\
\end{array}\right), \quad W_n := Z_n\left(\begin{array}{c} \oplus \\
\end{array}\right).$$

The total partition function is equal to

$$Z_n = Z_n^{\Gamma_n} = U_n + 2V_n + W_n.$$

Migdal-Kadanoff RG Equations:

$$U_{n+1} = (U_n^2 + V_n^2)^2, \quad V_{n+1} = V_n^2(U_n + W_n)^2, \quad W_{n+1} = (V^2 + W_n^2)^2.$$
See Part I [BLR1] for a proof. They give a homogeneous degree 4 polynomial map
\begin{equation}
(2.8) \quad \hat{R} : (U, V, W) \mapsto ((U^2 + V^2)^2, V^2(U + W), (W^2 + V^2)^2),
\end{equation}
called the Migdal Kadanoff Renormalization, with the property that \((U_n, V_n, W_n) = \hat{R}^n(U_0, V_0, W_0)\). We will also call the induced mapping \(R : \mathbb{C}P^2 \to \mathbb{C}P^2\) the Migdal Kadanoff Renormalization.

In order to express \(U_n, V_n, \) and \(W_n\) in terms of \(z\) and \(t\), one can use
\begin{equation}
(2.9) \quad U_0 = z^{-1}t^{-1/2}, \quad V_0 = t^{1/2}, \quad W_0 = zt^{-1/2}.
\end{equation}
as initial conditions for the iteration. According to (2.9), these coordinates are related to the “physical” \((z, t)\)-coordinates as follows:
\begin{equation}
(2.10) \quad (U : V : W) = \Psi(z, t) = (z^{-1}t^{-1/2} : t^{1/2} : zt^{-1/2}).
\end{equation}
Therefore, in the \((z, t)\)-coordinates, the renormalization transformation assumes the form:
\begin{equation}
(2.11) \quad R : (z, t) \mapsto \left( \frac{z^2 + t^2}{z^{-2} + t^2}, \frac{z^2 + z^{-2} + 2}{z^2 + z^{-2} + t^2 + t^{-2}} \right).
\end{equation}
Physically, the iterates \((z_n, t_n)\) are interpreted as the renormalized field-like and temperature-like variables.

By the basic symmetry of the Ising model, the change of sign of \(h\) interchanges the conditional partition functions \(U_n\) and \(W_n\) keeping \(V_n\) and the total sum \(Z_n\) invariant. Consequently, the RG transformation \(R\) commutes with the involution \((U : V : W) \mapsto (W : V : U)\), which is also obvious from the explicit formula (2.8), while \(R\) commutes with \((z, t) \mapsto (z^{-1}, t)\).

Existence of the renormalization mapping \(R\) makes discussion of the free energies \(\hat{F}_n^#\) especially clear for the DHL. Consider the linear form
\begin{equation}
Y_0(U, V, W) := U + 2V + W,
\end{equation}
which is chosen so that \(S_0^\#\) is the zero divisor of \(\sigma_{Y_0}\). For each \(n\),
\begin{equation}
(2.12) \quad \hat{Z}_n = Y_0 \circ \hat{R}^n \quad \text{and} \quad \hat{F}_n^\# = \log |\hat{Z}_n|.
\end{equation}
Here, \(\hat{Z}_n\) and \(\hat{F}_n^\#\) are expressed in terms of the initial values of \((U_0, V_0, W_0)\). To rewrite (2.12) in terms of \((z, t)\), we can pull back these expressions under \(\Psi\).

### 3. Global properties of the RG transformation in \(\mathbb{C}P^2\).

We will now summarize (typically without proofs) results from [BLR1] about the global properties of the RG mappings.

#### 3.1. Preliminaries.

The renormalization mappings \(R\) and \(R\) are semi-conjugate by the degree two rational mapping \(\Psi : \mathbb{C}P^2 \to \mathbb{C}P^2\) given by (2.10).

Both mappings have topological degree 8 (see Prop. 4.3 from Part I). However, as noted in the introduction, their algebraic degrees behave differently: \(R\) is algebraically stable, while \(R\) is not. Since \(\deg(R^n) = 4^n\), for any algebraic curve \(D\) of degree \(d\), the pullback \((R^n)^* D\) is a divisor of degree \(d \cdot 4^n\). (For background on divisors, see Appendix A.) For this reason, we will focus most of our attention on the dynamics of \(R\).
The semiconjugacy $\Psi$ sends Lee-Yang cylinder $C$ to a Mobius band $C$ that is invariant under $R$. It is obtained as the closure in $\mathbb{C}P^2$ of the topological annulus
\begin{equation}
C_0 = \{(u, w) \in \mathbb{C}^2 : w = u, \ |u| \geq 1\}.
\end{equation}

Let $T = \{(u, \bar{u}) : |u| = 1\}$ be the “top” circle of $C$, while $B$ be the slice of $C$ at infinity. In fact, $\Psi : C \to C$ is a conjugacy, except that it maps the bottom $B$ of $C$ by a 2-to-1 mapping to $B$ (see Prop. 3.1 from Part I).

3.2. Indeterminacy points for $R$. In homogeneous coordinates on $\mathbb{C}P^2$, the map $R$ has the form:
\begin{equation}
R : [U : V : W] \mapsto [(U^2 + V^2)^2 : V^2(U + W)^2 : (V^2 + W^2)^2)].
\end{equation}

One can see that $R$ has precisely two points of indeterminacy $a_+ := [i : 1 : -i]$ and $a_- := [-i : 1 : i]$. Resolving all of the indeterminacies of $R$ by blowing-up the two points $a_+$ (see Appendix A.3), one obtains a holomorphic mapping $\tilde{R} : \mathbb{C}P^2 \to \mathbb{C}P^2$.

In coordinates $\xi = u - i$ and $\chi = (w + i)/(u - i)$ near $a_+ = (i, -i)$, we obtain the following expression for the map $\tilde{R} : \mathbb{C}P^2 \to \mathbb{C}P^2$ near $L_{\text{exc}}(a_+)$:
\begin{equation}
u = \left(\frac{\xi + 2i}{1 + \chi}\right)^2, \quad w = \left(\frac{\chi^2 - 2i\chi}{1 + \chi}\right)^2.
\end{equation}

(Similar formulas hold near $a_- = (-i, i)$.) The exceptional divisor $L_{\text{exc}}(a_+)$ is mapped by $\tilde{R}$ to the conic
\begin{equation}G := \{(u - w)^2 + 8(u + w) + 16 = 0\}.
\end{equation}

3.3. Superattracting fixed points and their separatrices. We will often refer to $L_0 := \{V = 0\} \subset \mathbb{C}P^2$ as the line at infinity. It contains two symmetric superattracting fixed points, $e = (1 : 0 : 0)$ and $e' = (0 : 0 : 1)$. Let $W^e(e)$ and $W^e(e')$ stand for the attracting basins of these points. It will be useful to consider local coordinates $(\xi = W/U, \eta = V/U)$ near $e$.

The line at infinity $L_0 = \{\eta = 0\}$ is $R$-invariant, and the restriction $R|L_0$ is the power map $\xi \mapsto \xi^4$. Thus, points in the disk $\{|\xi| < 1\}$ in $L_0$ are attracted to $e$, points in the disk $\{|\xi| > 1\}$ are attracted to $e'$, and these two basins are separated by the unit circle $B$. We will also call $L_0$ the fast separatrix of $e$ and $e'$.

Let us also consider the conic
\begin{equation}L_1 = \{\xi = \eta^2\} = \{V^2 = UW\}
\end{equation}

passing through points $e$ and $e'$. It is an embedded copy of $\mathbb{C}P^1$ that is invariant under $R$, with $R|L_1(w) = w^2$, where $w = W/V = \xi/\eta$. Thus, points in the disk $\{|w| < 1\}$ in $L_1$ are attracted to $e$, points in the disk $\{|w| > 1\}$ are attracted to $e'$, and these two basins are separated by the unit circle $T$ (see §3.1 from Part I). We will call $L_0$ the slow separatrix of $e$ and $e'$.

If a point $x$ near $e$ (resp. $e'$) does not belong to the fast separatrix $L_0$, then its orbit is “pulled” towards the slow separatrix $L_1$ at rate $\rho^d$, with some $\rho < 1$, and converges to $e$ (resp. $e'$) along $L_1$ at rate $r^2$, with some $r < 1$.

The strong separatrix $L_0$ is transversally superattracting: all nearby points are pulled towards $L_0$ uniformly at rate $r^2$ (see also the proof of Lemma 5.3). It follows that these points either converge to one of the fixed points, $e$ or $e'$, or converge to the circle $B$. 
Given a neighborhood $\Omega$ of $B$, let

\[(3.5) \quad W_{s,\text{loc}}^s(B) = \{ x \in \mathbb{C}P^2 : \mathbb{R}^n x \in \Omega \ (n \in \mathbb{N}) \text{ and } \mathbb{R}^n x \to B \text{ as } n \to \infty \}\]

(where $\Omega$ is implicit in the notation, and an assertion involving $W_{s,\text{loc}}^s(B)$ means that it holds for arbitrary small suitable neighborhoods of $B$). It is shown in Part I (§9.2) that $W_{s,\text{loc}}^s(B)$ has the topology of a 3-manifold that is laminated by the union of holomorphic stable manifolds $W_{s,\text{loc}}^s(x)$ of points $x \in B$.

We conclude:

**Lemma 3.1.** $W^s(e) \cup W^s(e') \cup W_{s,\text{loc}}^s(B)$ fills in some neighborhood of $L_0$.

3.4. **Regularity of $W_{s,\text{loc}}^s(x)$**. For a diffeomorphism the existence and regularity of the local stable manifold for a hyperbolic invariant manifold $N$ has been studied extensively in [HPS]. In order to guarantee a $C^1$ local stable manifold $W_{s,\text{loc}}^s(N)$, a strong form of hyperbolicity known as *normal hyperbolicity* is assumed. Essentially, $N$ is normally hyperbolic for $f$ if the expansion of $Df$ in the unstable direction dominates the maximal expansion of $Df$ tangent to $N$ and the contraction of $Df$ in the stable direction dominates the maximal contraction of $Df$ tangent to $N$. See [HPS, Theorem 1.1]. If, furthermore, the expansion in the unstable direction dominates the $r$-th power of the maximal expansion tangent to $N$ and the contraction in the stable direction dominates the $r$-th power of the maximal contraction tangent to $N$, this guarantees that the stable manifold is of class $C^r$. The corresponding theory for endomorphisms is less developed, although note that other aspects of [HPS], related to persistence of normally hyperbolic invariant laminations, have been generalized to endomorphisms in [Be].

In our situation, $B$ is not normally hyperbolic because it lies within the invariant line $L_0$ and $R$ is holomorphic. This forces the expansion rates tangent to $B$ and transverse to $B$ (within this line) to coincide. Therefore, the following result does not seem to be part of the standard hyperbolic theory:

**Lemma 3.2.** $W_{s,\text{loc}}^s(x)$ is a $C^\infty$ manifold and the stable foliation is a $C^\infty$ foliation by complex analytic discs.

**Proof.** In Proposition 9.11 from Part I, we showed that within the cylinder $\mathcal{C}$ the stable foliation of $\mathcal{B}$ has $C^\infty$ regularity and that the stable curve of each point is real analytic. Mapping forward under $\Psi$, we obtain the same properties for the stable foliation of $B$ within $C$.

Let us work in the local coordinates $\xi = W/U$ and $\eta = V/U$. In these coordinates, $B = \{ \eta = 0, |\xi| = 1 \}$. The stable curve of some $\xi_0 \in B$ can be given by expressing $\xi$ as the graph of a holomorphic function of $\eta$:

\[(3.6) \quad \xi = h(\eta, \xi_0) = \sum_{i=0}^{\infty} a_i(\xi_0)\eta^i.\]

The right hand side is a convergent power series with coefficients depending on $\xi_0$, having a uniform radius of convergence over every $\xi_0 \in B$. The series is uniquely determined by its values on the real slice $C$, in which the leaves depend with $C^\infty$ regularity on $\xi_0$. Therefore, each of the coefficients $a_i(\xi_0)$ is $C^\infty$ in $\xi_0$. This gives that $W_{s,\text{loc}}^s(x)$ is a $C^\infty$ manifold.

\[\square\]
Figure 3.1. Critical locus for $R$ shown with the separatrix $L_0$ at infinity.

Remark 3.1. The technique from the proof of Lemma 3.2 applies to a more general situation: Suppose that $M$ is a real analytic manifold and $f : M \to M$ is a real analytic map. Let $N \subset M$ be a compact real analytic invariant submanifold for $f$, with $f|N$ expanding and with $N$ transversally attracting under $f$. Then, $N$ will have a stable foliation $W^s_{loc}(N)$ of regularity $C^r$, for some $r > 0$ (see the beginning of this subsection), with the stable manifold of each point being real-analytic. The stable manifold $W^s_{C,loc}(N)$ for the extension of $f$ to the complexification $M_C$ of $M$ will then also have $C^r$ regularity.

3.5. Critical locus. We showed in Part I [BLR1] that the critical locus of $R$ consists of 6 complex lines and one conic:

- $L_0 := \{ V = 0 \}$ = line at infinity,
- $L_1 := \{ UW = V^2 \}$ = conic $\{ uw = 1 \}$,
- $L_2 := \{ U = -W \}$ = $\{ u = -w \}$,
- $L_3^\pm := \{ U = \pm iV \}$ = $\{ u = \pm i \}$,
- $L_4^\pm := \{ W = \pm iV \}$ = $\{ w = \pm i \}$.

(Here the curves are written in the homogeneous coordinates $(U : V : W)$ and in the affine ones, $(u = U/V, w = W/V)$.) The critical locus is schematically depicted on Figure 3.1 while its image, the critical value locus, is depicted on Figure 3.2.

It will be helpful to also consider the critical locus for the lift $\tilde{R} : \mathbb{CP}^2 \to \mathbb{CP}^2$. Each of the critical curves $L_i$ lifts by proper transform (see Appendix A.3) to a critical curve $\tilde{L}_i \subset \mathbb{CP}^2$ for $\tilde{R}$. Moreover, any critical point for $\tilde{R}$ is either one of these proper transforms or lies within the exceptional divisors of $L_{exc}(a_\pm)$.

By symmetry, it is enough to consider the blow-up of $a_\pm$. We saw in Part I that there are four critical points on the exceptional divisor $L_{exc}(a_\pm)$ occurring at $\chi = -1, 1, \infty,$ and $0$, where $\chi = (w - i)/(u - 1)$. They correspond to intersections of $L_{exc}(a_\pm)$ with the collapsing line $\tilde{L}_2$, the $\tilde{L}_1$, and the critical lines $\tilde{L}_3^+$ and $\tilde{L}_4^-$, respectively.
Whitney Folds are a normal form for the simplest type of critical points of a mapping (see Appendix C). We have:

**Lemma 3.3.** All critical points of $\tilde{R}$ except the fixed points $e, e'$, the collapsing line $\tilde{L}_2$, and two points $\{\pm(i, i)\} = \tilde{L}_{2}^{\pm} \cap \tilde{L}_{4}^{\pm}$, are Whitney folds.

The only critical values obtained as images of non-Whitney folds are: $e, e'$, $b_0 = R(\tilde{L}_2) \in L_0$ and $0 := (0, 0) = R(\pm(i, i))$.

4. **Fatou and Julia sets and the measure of maximal entropy**

4.1. **Julia set.** For a rational map $R : \mathbb{CP}^n \to \mathbb{CP}^n$, the Fatou set $F_R$ is defined to be the maximal open set on which the iterates $\{R^m\}$ are well-defined and form a normal family. The complement of the Fatou set is the *Julia set* $J_R$.

If $R$ is dominant and has no collapsing varieties, Lemma A.1 gives that $R$ is locally surjective (except at indeterminate points), so that the Fatou set is forward invariant and consequently, the Julia set is backward invariant.

If $R$ has indeterminate points, then, according to this definition they are in $J_R$. In this case, $F_R$ and $J_R$ are not typically totally invariant. One can see this by blowing up an indeterminate point and observing that the image of the resulting exceptional divisor typically intersects $F_R$. Moreover, if $R$ is not algebraically stable, then any curve $C$ that is mapped by some iterate to the indeterminacy points is in $J_R$.

The Migdal-Kadinoff renormalization $R$ is not locally surjective at any $x \in L_2 \setminus L_0$. More specifically, if $U$ is a small neighborhood of $x$ that avoids $L_0$, then

$$R(U) \cap L_0 = b_0 = R(x),$$

since any point of $L_0 \setminus \{b_0\}$ only has preimages in $L_0$. However, we still have the desired invariance:

**Lemma 4.1.** The Migdal-Kadinoff renormalization $R$ has forward invariant Fatou set and, consequently, backward invariant Julia set.

**Proof.** It suffices to show that $L_2 \subset J_R$, since $R$ is locally surjective at any other point, by Lemma A.1. By definition, $\{a_{\pm}\} \subset J_R$, so we consider $x \in L_2 \setminus \{a_{\pm}\}$. Let $U$ be any small neighborhood of $x$. Note that $R(x) = b_0$ is a fixed point of saddle-type, with one-dimensional stable and unstable manifolds. Therefore, in
order for the iterates to form a normal family on $U$, we must have $R(U) \subset W^s(b_0)$. However, this is impossible, since there are plenty of regular points for $R$ in $U$. \hfill \Box 

Lemmas \ref{lem} and \ref{lem2} give a clear picture of $J_R$ in a neighborhood of the line at infinity $L_0$.

**Proposition 4.2.** Within some neighborhood $N$ of $L_0$, we have that $J_R \cap N = W^c_{\text{loc}}(B) \cap N$. Within this neighborhood, $J_R$ is a $C^\infty$-3-dimensional manifold.

Let us consider the locus $\{h = 0\}$ of vanishing magnetic field in $\mathbb{CP}^2$ for the DHL. In the affine coordinates, it is an $R$-invariant line $L_{\text{inv}} = \{u = w\}$; in the physical coordinates, it is an $R$-invariant line $L_{\text{inv}} = \{z = 1\}$. The two maps are conjugate by the Möbius transformation $L_{\text{inv}} \to L_{\text{inv}}, u = 1/t$. Dynamics of $R : t \mapsto \left(\frac{2t}{t^2 + 1}\right)^2$ on $L_{\text{inv}}$ was studied in [B3]. In particular, it is shown in that paper that the Fatou set for $R|_{L_{\text{inv}}}$ consists entirely of the basins of attraction of the fixed points $\beta_0$ and $\beta_1$ which are superattracting within this line: see Figure [1.C.3].

**Proposition 4.3.** $J_R = \bigcup_n R^{-n}(L_{\text{inv}})$.

**Proof.** Since $R|_{L_{\text{inv}}}$ is conformally conjugate to $R|_{L_{\text{inv}}}$, every point in the Fatou set of $R|_{L_{\text{inv}}}$ is in the basin of attraction of either $b_0$ or $b_1$. Since these points are of saddle-type in $\mathbb{CP}^2$, the family of iterates $R^n$ cannot be normal at any point on $L_{\text{inv}}$. Thus $L_{\text{inv}} \subset J_R$. It follows that $\bigcup_n R^{-n}(L_{\text{inv}}) \subset J_R$ since $J_R$ is closed and backward invariant.

We will now show that $\bigcup_n R^{-n}(L_{\text{inv}})$ is dense in $J_R$. Consider a configuration of five algebraic curves

- $X_0 := \{V = 0\}$ = the separatrix $L_0$,
- $X_1 := \{U = W\}$ = the invariant line $L_{\text{inv}}$,
- $X_2 := \{U = -W\}$ = the collapsing line $L_2 \subset R^{-1}(L_{\text{inv}})$,
- $X_3 := \{U^2 + 2V^2 + W^2 = 0\}$ = a component of $R^{-1}(L_{\text{inv}})$,
- $X_4 := \{U^4 + 2U^2V^2 + 2V^4 + 2W^2V^2 + W^4 = 0\}$ = a component of $R^{-1}(X_3)$.

We will use the results of M. Green to check that the complement of these curves, $M := \mathbb{CP}^2 \setminus \bigcup_i X_i$, is a complete Kobayashi hyperbolic manifold hyperbolically embedded in $\mathbb{CP}^2$ (see Appendix [A,6]). We will first check that $M$ is Brody hyperbolic, i.e., there are no non-constant holomorphic maps $f : \mathbb{C} \to M$. To this end we will apply Green’s Theorem [A,3]. It implies that the image of $f$ must lie in a line $L$ that is tangent to the conic $X_3$ at an intersection point with $X_i$, for one of the lines $X_i$, $i = 0, 1, 2$, and contains the intersection point $X_j \cap X_l$ of the other two lines. It is a highly degenerate situation which does not occur for a generic configuration. However, this is exactly what happens in our case, as the lines $X_0, X_1, X_2$ form a self-dual triangle with respect to the conic $X_3$ (see [A,7]). However, one can check by direct calculation that the last curve, $X_4$, must intersect each of these tangent lines $L$ in at least one point away from $X_0, \ldots, X_3$. Since any holomorphic map from $\mathbb{C}$ to $L \setminus \bigcup_i X_i$ must then omit 3 points in $L$, it must be constant.

So, $M$ is Brody hyperbolic. Moreover, for each $i = 0, \ldots, 4$ the remaining curves $\bigcup_{j \neq i} X_j$ intersect $X_i$ in at least three points so that there is no non-constant holomorphic map from $\mathbb{C}$ to $X_i \setminus \bigcup_{j \neq i} X_j$. Therefore, another of Green’s results (Theorem [A,3]) applies showing that $M$ is complete hyperbolic and hyperbolically
embedded. It then follows from Proposition 4.3 that the family \( \{ R^n \} \) is normal on any open set \( N \subset \mathbb{C} \mathbb{P}^2 \) for which \( R^n : N \to M \) for all \( n \).

Given any \( \zeta \in J_R \) and any neighborhood \( N \) of \( \zeta \), we’ll show that some preimage \( R^{-n}(L_{\text{inv}}) \) intersects \( N \). Since \( \zeta \in J_R \), the family of iterates \( R^n \) are not normal on \( N \), hence \( R^n(N) \) must intersect \( \bigcup_i X_i \) for some \( n \). If the intersection is with \( X_i \) for \( i > 0 \) then \( R^{n+2}(N) \) intersects \( L_{\text{inv}} \).

So, some iterate \( R^n(N) \) must intersect \( X_0 = L_0 \). Suppose first that \( \zeta \in L_0 \). Then, by Lemma 3.1 \( \zeta \subset W^s(e) \cup W^s(e') \cup T \). Since the first two basins are contained in the Fatou set, \( \zeta \in T \), where preimages of the fixed point \( \beta_0 \in L_{\text{inv}} \) are dense.

Finally, assume \( \zeta \notin L_0 \). By shrinking \( N \) if needed, we can make it disjoint from \( L_0 \). Hence there is \( n > 0 \) such that \( R^n(N) \) intersects \( L_0 \), while \( R^{-1}(L_0) \) is disjoint from \( L_0 \). But since \( R^{-1}(L_0) = L_0 \cup L_2 \), we conclude that \( R^{-1}(N) \) must intersect \( L_2 \). But \( L_2 \) collapses under \( R \) to the fixed point \( \beta_0 \in L_{\text{inv}} \). Hence \( R^n(N) \) intersects \( L_{\text{inv}} \). \( \square \)

We will now relate \( J_R \) to the Green current \( m^c \). (See Appendix B for the definition and basic properties of \( m^c \).)

**Proposition 4.4.** \( J_R = \text{supp } m^c \).

**Proof.** The inclusion \( \text{supp } m^c \subset J_R \) follows immediately from Theorem 3.5. We will use Proposition 4.3 to show that \( J_R \subset \text{supp } m^c \). Since \( \text{supp } m^c \) is a backward invariant closed set, it is sufficient for us to show that \( L_{\text{inv}} \subset \text{supp } m^c \).

Note that \( L_{\text{inv}} = W^s(b_0) \cup W^s(b_1) \cup J_{R\mid L_{\text{inv}}} \). The basin \( W^s(L_0) \) is open and contained within the regular set for \( R \) (see Appendix B for the definition), since every point of \( W^s(L_0) \) has a neighborhood whose forward iterates remain bounded away from \( \{ a_{\pm} \} \). Therefore, \( W^s(b_0) \subset J_R \cap N \subset \text{supp } m^c \), by Theorem 3.5. Since \( \text{supp } m^c \) is closed, we also have that \( J_{R\mid L_{\text{inv}}} \subset \text{supp } m^c \).

Every point of \( W^s(b_1) \) is irregular because

\[
b_1 \in \bigcup_{n \geq 0} R^{-n} \{ a_{\pm} \}.
\]

Therefore, we cannot directly use Theorem 3.5 to conclude that \( W^s(b_1) \subset \text{supp } m^c \).

Notice that the points of \( L_2 \setminus \{ a_{\pm} \} \) are regular, since they are in \( W^s(L_0) \). Theorem 3.5 gives that \( L_2 \setminus \{ a_{\pm} \} \subset \text{supp } m^c \), since \( L_2 \subset J_R \). Because \( \text{supp } m^c \) is closed, \( L_2 \subset \text{supp } m^c \). Let \( D_2 \subset L_2 \) be a small disc centered around \( a_+ \). Preimages of \( D_2 \) under appropriate branches of \( R^n \) will give discs intersecting \( L_1 \) transversally at a sequence of points converging to \( b_1 \). By the Dynamical \( \lambda \)-Lemma (see [PM] pp. 80-84), this sequence of discs will converge to \( W^s_0(b_1) \subset L_{\text{inv}} \), where \( W^s_0(b_1) \) is the immediate basin of \( b_1 \). Since each of the discs is in \( \text{supp } m^c \), and the latter is closed, we find that \( W^s_0(b_1) \subset \text{supp } m^c \). Further preimages show that all of \( W^s(b_1) \subset \text{supp } m^c \). \( \square \)

### 4.2. Fatou Set.

Because \( J_R = \text{supp } m^c \), we immediately have:

**Corollary 4.5.** The Fatou set of \( R \) is pseudoconvex.

For the definition of pseudoconvexity, see [Kra].
Proof. It is well-known that in the complement in \( \mathbb{CP}^2 \) of the support of a closed positive \((1,1)\)-current is pseudoconvex. See [C] Theorem 6.2 or [U] Lemma 2.4. \( \square \)

Computer experiments indicate that the Fatou set of \( R \) may be precisely the union of the basins of attraction \( W^s(e) \) and \( W^s(e') \) for the two superattracting fixed points \( e \) and \( e' \). See Problem D.34.

Consider the solid cylinders
\[
SC := \left\{ [U : V : W] : \frac{V^2}{UW} \in [0,1], \frac{W}{U} < 1 \right\} \quad \text{and} \quad SC' := \left\{ [U : V : W] : \frac{V^2}{UW} \in [0,1], \frac{W}{U} > 1 \right\}.
\]

Then, we can prove the following more modest statement:

**Proposition 4.6.** For the mapping \( R \) we have \( SC \subset W^s(e) \) and \( SC' \subset W^s(e') \).

In the proof, we will need to use an important property of \( R : C \rightarrow C \) that was proved in Part I. Recall from [C,1] that \( C = \Psi(C) \) is the invariant real Möbius and that \( C_0 = C \setminus B \) is the topological annulus obtained by removing the “core curve” \( B \).

The key property is:

(P9') Every proper path \( \gamma \) in \( C_0 \) lifts under \( R \) to at least 4 proper paths in \( C_0 \).

If \( \gamma \) crosses \( G \) at a single point, then \( R^{-1}\gamma = \cup \delta_i \).

**Proof of Prop. 4.6.** It suffices to prove the proposition for \( SC \), since the statement for \( SC' \) follows from the symmetry \( \rho \).

We will decompose \( SC \) as a union of complex discs and show that each disc is in \( W^s(e) \). Let
\[
P_e := \left\{ [U : V : W] : \frac{V^2}{UW} = c \in [0,1] \right\},
\]
and
\[
P_e^* := \left\{ [U : V : W] : \frac{V^2}{UW} = c \in [0,1], \frac{W}{U} < 1 \right\},
\]
so that \( SC = \bigcup_{c \in [0,1]} P_e^* \).

The discs \( P_e^* \) and \( P_e \) are in \( W^s(e) \) because they are each within the forward invariant critical curves \( L_0 \) and \( L_1 \), respectively, on which the dynamics is given by \( (W/U) \rightarrow (W/U)^4 \) and \( (W/U) \rightarrow (W/U)^2 \), respectively.

We now show that for any \( c \in (0,1) \) we also have \( P_e^* \subset W^s(e) \). In fact \( e \in P_e^* \), so it suffices to show that \( R^n \) forms a normal family on \( P_e^* \). Consider any \( x \in P_e^* \).

If \( x = e \), then \( x \in W^s(e) \) so that \( R^n \) is normal on some neighborhood of \( x \) in \( P_e^* \).

Now consider any \( x \in P_e^* \setminus \{e\} \). There is a neighborhood of \( N \subset P_e^* \) of \( x \) with \( e \notin N \), on which we will show that \( R^n \) forms a normal family. Recall the family of curves \( X_0, \ldots, X_4 \) from the proof of Proposition 4.3, where we showed that \( \mathbb{CP}^2 \setminus \bigcup X_i \) is complete hyperbolic and hyperbolically embedded. We will show for every \( n \) that \( R^n(N) \) is in \( \mathbb{CP}^2 \setminus \bigcup X_i \), so that \( R^n \) is normal on \( N \).

Since \( P_e^* \cap X_0 = \{e\} \), and \( e \notin N \), we have that \( N \cap X_0 = \emptyset \). Therefore, by reasoning identical to that in the proof of Proposition 4.3 if \( R^n(N) \) intersects \( X_i \) for any \( i = 0, \ldots, 4 \), we must have that some iterate \( R^m(N) \) intersects \( X_1 = L_{inv} \).

We will check that forward iterates of \( R^n(P_e^*) \) are disjoint from \( L_{inv} \), which is sufficient since \( N \subset P_e^* \). The line \( L_{inv} \) intersects the invariant annulus \( C_0 \) in two
4.3. Measure of Maximal Entropy.

There is a conjecture specifying the expected ergodic properties of a dominant rational map of a projective manifold in terms of the relationship between various “dynamical degrees” of the map; see [Gu2].

Since the Migdal-Kadanoff renormalization $R$ is an algebraically stable map of $\mathbb{CP}^2$, there are only two relevant dynamical degrees, the topological degree $\deg_{\text{top}} R$ and the algebraic degree $\deg R$, which satisfy

$$\deg_{\text{top}} R = 8 > 4 = \deg R.$$ 

This case of high topological degree was studied by Guedj [Gu1], who made use of a bound on topological entropy obtained by Dinh and Sibony [DS1]. In our situation, his results give

**Proposition 4.8.** $R$ has a unique measure $\nu$ of maximal entropy $\log 8$ with the following properties

(i) $\nu$ is mixing;

(ii) The Lyapunov exponents of $\nu$ are bounded below by $\log \sqrt{2}$;

(iii) If $\theta$ is any probability measure that does not charge the postcritical set of $R$, then $8^{-n}(R^n)^*\theta \to \nu$;

---

5It is stated more generally in [Gu2], for meromorphic maps of compact Kähler manifolds.
(iv) If $P_n$ is the set of repelling periodic points of $R$ of period $n$ that are in $\text{supp } \nu$, then $8^{-n} \sum_{a \in P_n} \delta_a \to \nu$.

The measure $\nu$ satisfies the backwards invariance $R^* \nu = 8 \nu$, hence its support is totally invariant. In our situation, $\text{supp } \nu \subseteq J_R$ because (for example) the points in $\mathcal{W}^s(B)$ are not in $\text{supp } \nu$. It can be thought of as the “little Julia set” within $J_R$ on which the “most chaotic” dynamics occurs.

**Remark 4.1.** The statement of (iv) is slightly different in [Gu1], where it is not emphasized that the repelling periodic points being considered are actually on $\text{supp } \nu$. However, it follows from the proof in [Gu1] and the fact that $\text{supp } \nu$ is totally invariant. See [DS3, Thm 1.4.13] for the analogous argument for holomorphic $f : \mathbb{CP}^2 \to \mathbb{CP}^2$.

**Remark 4.2.** We know very little about the support of $\nu$. See Problem D.3.

5. Pluri-potential interpretation

We will now prove the Global Lee-Yang-Fisher Theorem, assuming some technical volume estimates that will be saved for §6.

**Theorem 5.1.** Let us consider a linear form $Y = pU + qV + rW$ with $p \neq 0$, $r \neq 0$. Then the limit

$$G = \lim_{n \to \infty} \frac{1}{4^n} \log |Y \circ R^n|$$

exists in $L^1_{\text{loc}}(\mathbb{C}^3)$ and is equal to the Green potential of $R$ (see Appendix B.2).

The form $Y$ determines a section $\sigma_Y$ of the co-tautological line bundle on $\mathbb{CP}^2$ (see §A.4). Let $L_Y$ stand for the line $\sigma_Y = 0$ in $\mathbb{CP}^2$ and $[L_Y]$ stand for the current of integration over $Y$. Its pullbacks $(R^n)^*[L_Y]$ under the iterates of $R$ are the currents of integration over the zero divisor of $(R^n)^* \sigma_Y = \sigma_Y \circ R^n$.

Let $m^c = \pi_*(\Delta^p G)$ stand for Green $(1,1)$-current of $R$ (see §B.3). Applying the pluri-Laplacian $\Delta_p = \frac{i}{\pi} \partial \bar{\partial}$ to (5.1) we obtain:

**Corollary 5.2.** Let $Y$ be a linear form as in Theorem 5.1. Then, the normalized currents of integration $4^{-n}(R^n)^*[L_Y]$ weakly converge to the Green current $m^c$. In particular, this current gives the asymptotic distribution of $\frac{1}{4^n}(R^n)^*[S_0^c]$ in $\mathbb{CP}^2$.

**Remark 5.1.** The Global Lee-Yang-Fisher Theorem follows from Theorem 5.1 and Corollary 5.2 after pulling everything back under $\Psi$.

**Remark 5.2.** Note that the principal Lee-Yang-Fisher locus $S_0^c$ does not lie within the Julia set $J_R$. In fact, the only holomorphic curves inside the stable manifold $\mathcal{W}^s_{C,\text{loc}}(B)$ are the stable manifolds $W^s_{\text{loc}}(\zeta)$ of various points $\zeta \in B$. Indeed, the maximal complex subspace within the tangent space $E_x := T_x(\mathcal{W}^s_{C,\text{loc}}(B))$ is the complex line $E^c_\zeta := E_x \cap iE_x$, which thus must coincide with $T_xW^s_{\text{loc}}(\zeta)$. Consequently, any holomorphic curve inside $W^s_{C,\text{loc}}(B)$ must be tangent to the line field $T_xW^s_{\text{loc}}(\zeta)$.

For the same reason, none of the LY loci $S_n = (R^n)^*[S_0^c]$ lies inside $J_R$ either. However, in the limit they are distributed within the Julia set $J_R = \text{supp } m^c$. 
5.1. **Proof of Theorem 5.1.** Let us now endow the co-tautological bundle with the standard Hermitian structure $\| \cdot \|$ (see §A.4).

**Theorem 5.3.** Let $Y$ be a linear form as in Theorem 5.1. Then

$$\frac{1}{4^n} \log \| \sigma_Y \circ R^n \| \to 0 \quad \text{in} \ L^1_{loc}(\mathbb{CP}^2) \quad \text{as} \quad n \to \infty.$$  

We will first show that Theorem 5.3 implies Theorem 5.1:

**Proof.** Let $X = (U, V, W) \in \mathbb{C}^3 \setminus \{0\}$ and $x = (U : V : W) \in \mathbb{CP}^2$. By (A.5) from Appendix B we have:

$$\| \sigma_Y(x) \| = \frac{|Y(X)|}{\|X\|} \quad \text{and} \quad \| \sigma_Y(R^n x) \| = \frac{|Y(R^n X)|}{\|R^n X\|}.$$

Outside of the measure zero set $\hat{R}^{-n}\{Y = 0\}$, we can take $4^{-n} \log$ and apply Theorem 5.3 to obtain

$$\frac{1}{4^n} \log |Y(R^n X)| = \frac{1}{4^n} \log \|R^n X\| + o(1)$$

in $L^1_{loc}(\mathbb{C}^3)$. But since $R$ is algebraically stable the quantity in the right-hand side converges to the Green potential (Theorem B.3). □

**Proof of Theorem 5.3:** We will use the following general convergence criterion:

**Lemma 5.4.** Let $\phi_n$ be a sequence of $L^2$ functions on a finite measure space $(X, m)$ with bounded $L^2$-norms. If $\phi_n \to 0$ a.e. then $\phi_n \to 0$ in $L^1$.

**Proof.** Take any $\epsilon > 0$ and $\delta > 0$. By Egorov’s Lemma, there exists a set $X' \subset X$ with $m(X \setminus X') < \epsilon$ such that $\phi_n \to 0$ uniformly on $X'$. So, eventually the sup-norms of the $\phi_n$ on $X'$ are bounded by $\delta$. Hence

$$\int |\phi_n| \, dm = \int_{X'} |\phi_n| \, dm + \int_{X \setminus X'} |\phi_n| \, dm \leq \delta \cdot m(X) + B\sqrt{\epsilon},$$

where the last estimate follows from the Cauchy-Schwarz Inequality (with $B$ the $L^2$-bound on the $\phi_n$). The conclusion follows □

Next, we will estimate how badly volume can increase under iterated pullbacks by $R$. Any volume form on $\mathbb{CP}^2$ will suffice for our discussion, but it will be helpful to normalize, so that $\text{vol}\mathbb{CP}^2 = 1$. Let us fix a small forward-invariant tubular neighborhood $\Omega \subset \mathbb{CP}^2$ of the super-attracting line $L_0 = \{V = 0\}$ at infinity, and let $\Omega'$ be a neighborhood of $0 = [0 : 1 : 0]$ such that $R(\Omega') \subset \Omega$.

**Lemma 5.5.** There exists $C_1 > 0$ so that for any measurable set $X \subset \mathbb{CP}^2 \setminus (\Omega \cup \Omega')$,

$$\text{vol}(R^{-n} X) \leq C_1 (\text{vol} X)^{1/2^n}.$$  

**Proof.** Let $\tilde{R} : \mathbb{C}\mathbb{P}^2 \to \mathbb{CP}^2$ stand for the lift of $R$ to the projective space with both indeterminacy points, $a_{\pm}$, blown up. By Lemma 5.3 all critical values of $\tilde{R}$ in $\mathbb{CP}^2 \setminus (\Omega \cup \Omega')$ are generic, that is, they come from Whitney folds only. The conclusion now follows from the volume transformation rule near folds (Lemma C.3) and the fact that the projection $\pi : \mathbb{C}\mathbb{P}^2 \to \mathbb{CP}^2$ does not increase volume. □
Let now \( D \) be a tubular neighborhood of the line \( Y \) that stays away from the fixed points \( e, e' \). Let \( D_0 = D \cap \Omega, D'_0 = D \cap \Omega' \), and \( D''_0 = D \setminus \{ \Omega \cup \Omega' \} \). We inductively define
\[
D_k = R^{-1}D_{k-1} \cap \Omega, \quad D'_k = R^{-1}D_{k-1} \setminus \Omega, \quad \text{and} \quad D''_k = R^{-1}D'_{k+1}.
\]
Since \( R(\Omega) \subset \Omega \) and \( R(\Omega') \subset \Omega \), the sets \( D''_k \) and their further pullbacks lie outside \( \Omega \cup \Omega' \). For any \( n \geq 0 \),
\[
R^{-n}D = D_n \cup D'_n \cup \bigcup_{k=0}^{n-2} R^{-(n-k)}D''_k.
\]

**Lemma 5.6.** There exists \( C_2 > 0 \) so that
\[
\text{vol } D_k \leq C_2 \text{vol } D_0, \quad \text{vol } D'_k \leq C_2 (\text{vol } D_0)^{1/4}, \quad \text{and} \quad \text{vol } D''_k \leq C_2 (\text{vol } D_0)^{1/16}.
\]

**Proof.** Recall that the map \( R \) is the fourth-power map \( \xi \mapsto \xi^4 \) on line \( L_0 \), and this line is normally super-attracting. Since \( D \) stays away from the fixed points \( e, e' \) of \( R|L_0 \), the set \( D_k \) comprises \( 4^k \) tubes of full vertical size (in the tubular neighborhood \( \Omega \)) and of horizontal size of order \( 1/4^k \). The total volume of these tubes is of order \( 4^{-k} \text{vol } D_0 \) (which is much better than needed).

The last two estimates now follow from the power transformation rule for the volume (Corollary 6.6). \( \square \)

**Lemma 5.7.** There exists \( C_3 > 0 \) so that \( \text{vol } R^{-n}D \leq C (\text{vol } D)^{\gamma/2^n} \) with \( \gamma = 1/16 \).

**Proof.** By the same lemma, \( \text{vol } D'_n \leq C_2 (\text{vol } D_0)^{1/16} \). By Lemma 5.5,
\[
\text{vol}(R^{-(n-k)}D''_k) \leq C'(\text{vol } D_0)^{1/2^{n-k+4}}.
\]
Summing these up over \( k \), we obtain the desired inequality. \( \square \)

Let
\[
\phi_n = \frac{1}{4^n} \log \| \sigma_Y \circ R^n \|.
\]
Let us estimate distribution of the tales of these “random variables”.

**Lemma 5.8.** Let \( M = \sup \log \| \sigma_Y \| \). Then, there exists \( C > 0 \) so that
\[
\text{vol } \{ |\phi_n| > r \} \leq C \exp(-\gamma r 2^n) \quad \text{for any } r > M 4^{-n},
\]
where \( \gamma \) is the same as in Lemma 5.7.

**Proof.** We have:
\[
X_n(r) := \{ |\phi_n| > r \} = \{ \log \| \sigma \circ R^n \| > r 4^n \} \cup \{ \log \| \sigma \circ R^n \| < -r 4^n \} = \{ \log \| \sigma_Y \circ R^n \| < -r 4^n \} = R^{-n} \{ \| \sigma_Y \| < \exp(-r 4^n) \}
\]
(we have used that \( \log \| \sigma_Y \| < r 4^n \)). Using Lemma 5.7 together with the fact that a tubular neighborhood \( \{ \| \sigma_Y \| < k \} \) of the projective line \( L_Y \) has volume \( \leq C_4 k^2 \), we find
\[
\text{vol } X_n(r) \leq C \exp(-\gamma r 2^n + 1).\]
\( \square \)

We are ready to show that the functions \( \phi_n \) satisfy the conditions of Lemma 5.4

**Lemma 5.9.** The sequence \( \phi_n \) is \( L^2 \)-bounded.
Proof. We have:
\[ \| \phi_n \|^2 \leq \sum_{k=0}^{\infty} (k+1)^2 \text{vol}\{|\phi_n| \geq k\}. \]

By Lemma 5.8, this sum is bounded by
\[ \text{vol}(\mathbb{CP}^2) \sum_{k=0}^{M+1} (k+1)^2 + C \sum_{k=M}^{\infty} (k+1)^2 \exp(-\gamma k 2^{n+1}) \]
\[ \leq C_0 + C \sum_{k=0}^{\infty} (k+1)^2 \exp(-\gamma k) < \infty. \]

Lemma 5.10. The sequence \( \phi_n \) exponentially converges to 0 almost everywhere.

Proof. Fix any \( \lambda \in (1, 2) \). For sufficiently large \( n \), we have \( \lambda^{-n} > M 4^{-n} \), hence Lemma 5.8 gives
\[ \text{vol}\{|\phi_n| > \lambda^{-n}\} \leq C \exp(-2\gamma (2\lambda^{-1})^n). \]
Since the sum of these volumes converges, the Borel-Cantelli Lemma gives that for a.e. \( x \in \mathbb{CP}^2 \), we eventually have \( |\phi_n(x)| \leq \lambda^{-n} \).

This completes the proof of Theorem 5.3.

\[ \square \]

6. Volume Estimates

We will now prove the estimates on how volumes transform under a single iterate of \( R \) that were used in the proof from \[ \text{§5.} \]

Lemma 6.1. Let \( X \subset \mathbb{CP}^2 \) be a measurable set such that \( d := \text{dist}(X, \{e, e', b_0, 0\}) > 0 \). Then \( \text{vol}(R^{-1}X) = O(\sqrt{\text{vol}X}) \), with the constant depending on \( d \).

Proof. By Lemma 3.3, \( e, e', b_0, 0 \) are the only critical values of \( \tilde{R} \) that are images of critical points more complicated than Whitney folds. By Corollary 6.3, \( \text{vol}(\tilde{R}^{-1}X) = O(\sqrt{\text{vol}X}) \). Since \( R^{-1}X = \pi(\tilde{R}^{-1}X) \), where the projection \( \pi : \mathbb{CP}^2 \to \mathbb{CP}^2 \) is regular, we are done.

\[ \square \]

6.0.1. \( R \) near the collapsing line \( L_2 \).

Lemma 6.2. Near any finite point of the collapsing line \( L_2 \), except two indeterminacy points \( a_ \pm \), we have: \( \det DR \propto (u + w)^2 \).

Proof. We will use coordinate
\[ \lambda = \frac{u + w}{1 + u^2} \]
near the collapsing line \( L_2 = \{u + w = 0\} \) and outside the indeterminacy points \( a_ \pm \).

Recall that \( R(L_2) = b_0 \). In local coordinates \( (z = 1/u, \, \zeta = u/w) \) near \( b_0 \), the map \( R \) assumes form
\[ z = \lambda^2, \quad \zeta = \sigma^2, \quad \text{where} \quad \sigma = \frac{1 + u^2}{1 + w^2}. \]

So, in local coordinates \( R : (\lambda, u) \mapsto (z, \zeta) \) we have:
\[ \det DR = 4\lambda \sigma \frac{\partial \sigma}{\partial u}. \]

(6.1)
Moreover, in these coordinates, \( w = (1 + u^2)\lambda - u, \) \( \partial w/\partial u = 2u\lambda - 1, \) and an elementary calculation yields:

\[
\frac{\partial \sigma}{\partial u} = 2 - \frac{1 + u^2}{(1 + w^2)^2} (1 - uw) \lambda \sim \lambda,
\]
as long as we stay near \( L_2 \) and away from \( a_{\pm}. \) The conclusion follows. \( \square \)

A similar estimate holds in the blow-up coordinates near the indeterminacy points \( a_{\pm} \in L_2. \) Recall that \( \tilde{R} \) is given by (3.3) in the blow-up coordinates \((\xi, \chi).\) Moreover, \( \tilde{L}_2 = \{\chi = -1\}. \)

**Lemma 6.3.** Near \( \tilde{L}_2 \cap L_{\text{exc}}(a_{\pm}) \) we have:

\[
\det D\tilde{R} \asymp (1 + \chi)^2.
\]

**Proof.** It is sufficient to treat \( \tilde{R} \) near \( \tilde{L}_2 \cap L_{\text{exc}}(a_+). \) In the local coordinates

\[
\lambda = \frac{1 + \chi}{\xi + 2i}, \quad \xi = u - i
\]

near \( L_{\text{exc}}(a) \) and local coordinates \((z = 1/u, \zeta = u/w)\) near \( b_1, \) the map \( \tilde{R} \) assumes form

\[
z = \lambda^2, \quad \zeta = \sigma^2, \quad \text{where} \quad \sigma = \frac{\xi + 2i}{\chi^2 - 2i\chi}.
\]

So

\[
\det DR = 4\lambda\sigma \frac{\partial \sigma}{\partial \xi}.
\]

Moreover, \( \frac{\partial \chi}{\partial \xi} = \lambda, \) and an elementary calculation yields: \( \frac{\partial \sigma}{\partial \xi} \asymp \lambda \) near \( \tilde{L}_2 \cap L_{\text{exc}}(a_+). \) \( \square \)

We can now estimate how the volume is transformed near \( L_2: \)

**Lemma 6.4.** Any point \( x \in \tilde{L}_2 \setminus L_0 \) has a neighborhood \( D \) such that for any measurable set \( X \subset \mathbb{CP}^2, \) we have:

\[
\text{vol}(\tilde{R}^{-1}X \cap D) \leq C(\text{vol}X)^{1/3}
\]

**Proof.** By Lemmas 6.2 and 6.3 there are coordinates \((\lambda, \xi)\) near \( x \) such that \( \tilde{L}_2 = \{\lambda = 0\}, \) \( \xi(x) = 0, \) and \( \det DR \asymp \lambda^2. \) Let us take a bidisk neighborhood of \( x: \)

\[
D = \{ |\lambda| < \epsilon \} \times \{ |\xi| < \epsilon \}.
\]

For any measurable set \( Y \) that is sufficiently close to \( x \) we have:

\[
\text{vol}(\tilde{R}(Y)) \asymp \int_{Y^v} \int_{|\lambda| < \epsilon} |\lambda|^4 (\text{area} Y^h_\xi) \ d\text{area}(\lambda) \ d\text{area}(\xi),
\]

where \( Y^v \) is the projection of \( Y \) onto \( L_2 \) (in the bidisk coordinates) and \( Y^h_\xi \) are the slices of \( Y \) by the horizontal sections of \( D. \) But the inner integral above is exactly \((1/9)\text{area}(Q_3(Y^h_\xi)), \) where \( Q_3(\lambda) = \lambda^3. \) By Lemma C.2 it is bounded from below by \((1/9)(\text{area} Y^h_\xi)^3). \) By the Hölder inequality,

\[
(\text{area} Y^v)^2 \int_{Y^v} (\text{area} Y^h_\xi)^3 \ d\text{area}(\xi) \geq \left( \int_{Y^v} (\text{area} Y^h_\xi) \ d\text{area}(\xi) \right)^3 = (\text{vol}Y)^3.
\]

The conclusion follows, since \( \text{area} Y^v > (\text{area} Y^v)^2. \) \( \square \)
Let us also take care of one special point:

**Lemma 6.5.** The intersection point $L_2 \cap L_0$ has a neighborhood $D$ such that for any measurable set $X \subset \mathbb{CP}^2$, we have:

$$\text{vol}(R^{-1}X \cap D) \leq C(\text{vol } X)^{1/3}$$

**Proof.** The proof will be very similar to Lemma 6.4.

Let $\xi = W/U$ and $\eta = V/U$. Near $\{x\} = L_2 \cap L_0$, we will use the local coordinates $(\tau = \xi + 1, \eta)$ and near $b_0$ we use $(\xi, \eta)$. We have

$$\xi' = \left(\frac{\eta^2 + (\tau - 1)^2}{1 + \eta^2}\right)^2, \quad \eta' = \frac{\tau^2 \eta^2}{(1 + \eta^2)^2},$$

where $(\xi', \eta') = R(\tau, \eta)$. Sufficiently close to $x$ we have $\det DR \propto \eta \tau^2$. Let $D = \{|\tau| < \epsilon\} \times \{|\eta| < \epsilon\}$. If $Y \subset D$, we have

$$\text{vol}(R(Y)) \approx \int_{Y^{\tau}} |\eta|^2 \int_{|\tau| < \epsilon} |\tau|^4 (\text{area } Y^h_{\eta}) d \text{area}(\tau) d \text{area}(\eta).$$

As in the proof of Lemma 6.4, we must show that

$$\int_{Y^{\tau}} |\eta|^2 (\text{area } Y^h_{\eta})^3 d \text{area}(\eta) \geq C(\text{vol } Y)^3.$$

The Hölder Inequality gives that

$$ \left( \int_{Y^\eta} \frac{1}{|\eta|} d \text{area}(\eta) \right)^2 \int_{Y^\eta} |\eta|^2 (\text{area } Y^h_{\eta})^3 d \text{area}(\eta) \geq \left( \int_{Y^\eta} (\text{area } Y^h_{\eta}) d \text{area}(\eta) \right)^3 = (\text{vol } Y)^3.$$

The result follows since $1/|\eta|$ is locally integrable. \qed

**Corollary 6.6.** For any measurable set $X$ on distance at least $d$ from the fixed points $e, e'$, we have:

$$\text{vol}(R^{-1}X) \leq C(d)(\text{vol } X)^{1/4}.$$

**Proof.** The only critical value in $X$ not covered by Lemmas 6.1, 6.4 and 6.5 is $0$. But at the corresponding critical points, $\pm(i, i)$, the map $R$ is composed of a diffeomorphism and the squaring map $(u, w) \mapsto (u^2, w^2)$. It follows that at these points the volume is transformed with exponent 4. \qed

**Appendix A. Elements of complex geometry**

We are primarily interested in rational maps between complex projective spaces in two dimensions. However, greater generality will be useful in many circumstances (for example, in order to study a rational mapping near its indeterminate points). Much of the below material can be found with greater detail in [Da, De, GH, Shaf, K, La].
A.1. Projective varieties and rational maps. Let \( \pi : \mathbb{C}^{k+1} \setminus \{0\} \to \mathbb{CP}^k \) denote the canonical projection. Given \( z \in \mathbb{CP}^k \), any \( \tilde{z} \in \pi^{-1}(z) \) is called a lift of \( z \). One calls \( V \subset \mathbb{CP}^k \) a (projective) algebraic hypersurface if there is a homogeneous polynomial \( P : \mathbb{C}^{k+1} \to \mathbb{C} \) so that
\[
V = \{ z \in \mathbb{CP}^k : \ P(\tilde{z}) = 0 \}.
\]
More generally, a (projective) algebraic variety is the locus satisfying a finite number of homogeneous polynomial equations. Any algebraic variety \( V \) has the structure of a smooth manifold away from a proper subvariety \( V_{\text{sing}} \subset V \) and the dimension of \( V \setminus V_{\text{sing}} \) is called the dimension of \( V \). One calls \( V \) a projective manifold if \( V_{\text{sing}} = \emptyset \).

A rational map \( R : \mathbb{CP}^k \to \mathbb{CP}^l \) is given by a homogeneous polynomial map \( \hat{R} : \mathbb{C}^{k+1} \to \mathbb{C}^{l+1} \) for which we will assume the components have no common factors. One defines \( R(z) := \pi(\hat{R}(\tilde{z})) \) if \( \hat{R}(\tilde{z}) \neq 0 \), and otherwise we say that \( z \) is an indeterminacy point for \( R \). Since \( \hat{R} \) is homogeneous, the above notions are well-defined. Because the components of \( \hat{R} \) have no common factors, the set of indeterminate points \( I(R) \) is a projective variety of codimension greater than or equal to two.

Given two projective varieties, \( V \subset \mathbb{CP}^k \) and \( W \subset \mathbb{CP}^l \), a rational map \( R : V \to W \) is the restriction of a rational map \( \hat{R} : \mathbb{CP}^k \to \mathbb{CP}^l \) such that \( R(V \setminus I(R)) \subset W \). As above, \( I(R) \subset V \) is a projective subvariety of codimension greater than or equal to two in \( V \). If \( I(R) = \emptyset \), we say that \( R \) is a (globally) holomorphic (regular) map.

A rational mapping \( R : V \to W \) between projective manifolds is dominant if there is a point \( z \in V \setminus I(R) \) such that \( \text{rank } DR(z) = \dim W \).

We will call a subvariety \( U \subset V \) a collapsing variety for \( R \) if \( \dim(R(U)) < \dim(U) \).

Lemma A.1. Let \( R : V \to W \) be a dominant rational map between projective manifolds of the same dimension. If \( z \) is not an indeterminate point for \( R \) and not on any collapsing variety for \( R \), then \( R \) is locally surjective at \( z \).

Proof. The result will follow from the basic local properties of complex analytic subsets of \( \mathbb{C}^k \); see, for example, [De] Ch. II, §4.2 or [GH] Ch. 0.1.

Let \( z^* \in V \) satisfy the hypotheses of the Lemma. Choose local coordinates \( z = (z_1, \ldots, z_m) \) in a neighborhood of \( z^* \) and \( w = (w_1, \ldots, w_m) \) in a neighborhood \( R(z^*) \). Together, form local coordinates on \( V \times W \) in some small neighborhood \( N = N_z \times N_w \) of \( (z^*, R(z^*)) \). Within \( N \), the graph of \( R \) is a complex analytic variety \( G \) of dimension \( m \), since \( R \) is dominant. Without loss of generality, we suppose that \( z^* = R(z^*) = 0 \).

By choosing \( N \) sufficiently small, we can suppose that \( I(R) \cap N = \emptyset \) and that
\[
G \cap \{ w = 0 \} = \{(0,0)\},
\]
since \( z^* \) is not on a collapsing variety. In particular, \( G \) does not contain the \( z_m \)-axis. Let \( \pi_m : N \to N \cap \{ z_m = 0 \} \) be the projection. Using the Weierstrass Preparation Theorem and the resultant, as described on the top of p. 13 from [GH], we see that \( \pi_m(G) \) is a dimension \( m \) variety within \( N \cap \{ z_m = 0 \} \). Moreover, by (A.1), we have that
\[
\pi_m(G) \cap \{ w = 0 \} = \{(\tilde{z}, w) = (0,0)\}.
\]
(Here \( \tilde{z} = (z_1, \ldots, z_{m-1}) \).) Repeating this procedure \( m-1 \) more times, we see that \( \pi : G \cap N \to N_w \) is a branched covering. In particular, it is surjective. \( \square \)
A.2. Divisors. Divisors are a generalization of algebraic hypersurfaces that behave naturally under dominant rational maps. We will present an adaptation of material from [Da, Ch. 3] and [Shaf] suitable for our purposes.

A divisor $D$ on a projective manifold $V$ is a collection of irreducible hypersurfaces $C_1, \ldots, C_r$ with assigned integer multiplicities $k_1, \ldots, k_r$. One writes $D$ as a formal sum

$$D = k_1 C_1 + \cdots + k_r C_r.$$  

Alternatively, $D$ can be described by choosing an open cover $\{U_i\}$ of $V$ and rational functions $g_i : U_i \to \mathbb{C}$ with the compatibility property that $g_i/g_j$ is a non-vanishing holomorphic function on $U_i \cap U_j \neq \emptyset$. Taking zeros and poles of the $g_i$ counted with multiplicities, we obtain representation (A.2).

If $f : V \to W$ is a dominant holomorphic map, and $D = \{U_i, g_i\}$ is a divisor on $W$, the pullback $f^* D$ is the divisor on $V$ given by $\{f^{-1}U_i, f^*g_i\} \equiv \{f^{-1}U_i, g_i \circ f\}$. If $R : V \to W$ is a dominant rational map (with $I(R) \neq \emptyset$), we define $R^* D$ by first pull-backing $D$ under $R : V \setminus I(R) \to W$. Since $I(R)$ is a finite collection of points, the result (in terms of local defining functions) can be extended trivially to obtain a divisor $R^* D$ on all of $V$. Since the trivial extension of a divisor is unique, the result is well-defined.

A.3. Blow-ups. Given a pointed projective surface $(V, p)$, the blow-up of $V$ at $p$ is another projective surface $\tilde{V}$ with a holomorphic projection $\pi : \tilde{V} \to V$ such that

- $L_{\text{exc}}(p) := \pi^{-1}(p)$ is a complex line $\mathbb{CP}^1$ called the exceptional divisor;
- $\pi : \tilde{V} \setminus L_{\text{exc}}(p) \to V \setminus \{p\}$ is a biholomorphic map.

See [GHI, Shaf].

The construction has a local nature near $p$, so it is sufficient to provide it for $(\mathbb{C}^2, 0)$. The space of lines $l \subset \mathbb{C}^2$ passing through the origin is $\mathbb{CP}^1$, by definition. Then $\mathbb{C}^2$ is realized as the surface $X$ in $\mathbb{C}^2 \times \mathbb{CP}^1$ given by equation $\{(u, v) \in l\}$ with the natural projection $(u, v, l) \mapsto (u, v)$. In this model, points of the exceptional divisor $L_{\text{exc}} = \{(0, 0, l) : l \in \mathbb{CP}^1\}$ get interpreted as the directions $l$ at which the origin is approached.

Any line $l \subset \mathbb{C}^2$ naturally lifts to the “line” $\tilde{l} = \{(u, v, l) : (u, v) \in l\}$ in $\tilde{\mathbb{C}}^2$ crossing the exceptional divisor at $(0, 0, l)$. Moreover, $\tilde{\mathbb{C}}^2 \setminus \tilde{l}$ is isomorphic to $\mathbb{C}^2$. Indeed, let $\phi(u, v) = au + bv$ a linear functional that determines $l$. It is linearly independent from one of the coordinate functionals, say with $v$ (so $a \neq 0$). Then

$$(u, v, l) \mapsto (\phi, \kappa := v/\phi)$$

is a local chart that provides a desired isomorphism. In particular, two charts corresponding to the coordinate axes in $\mathbb{C}^2$ provide us with local coordinates $(u, \kappa = v/u)$ and $(v, \kappa = u/v)$ which are usually used in calculations.

The value of this construction lies in the fact that it can be used to resolve the indeterminacy of a rational map; see [Shaf] Ch. IV, §3.3. Moreover, any analytic curve $C$ on $V$ lifts to an analytic curve $\tilde{C} := \pi^{-1}(C \setminus \{p\})$ on $\tilde{V}$, known as the proper transform of $C$, which tends to have milder singularities than $C$; see [Shaf] Ch. IV, §4.1. Taking multiplicities into consideration, the proper transform of a divisor $D$ is defined similarly.

---

6This turns $\tilde{\mathbb{C}}^2$ into a line bundle over $\mathbb{CP}^1$ known as the tautological line bundle.
A.4. Divisors on $\mathbb{CP}^m$, Degree, and the Co-tautological line bundle. Associated to any homogeneous polynomial $P : \mathbb{C}^{m+1} \to \mathbb{C}$ is a divisor $D_P$ given by $\{ U_i, P \circ \sigma_i \}$, where the $\{ U_i \}$ form an open covering of $\mathbb{CP}^m$ that admits local sections $\sigma_i : U_i \to \mathbb{C}^{m+1} \setminus \{ 0 \}$ of the canonical projection $\pi$. Furthermore, every divisor can be described as a difference $D = D_P - D_Q$ for appropriate $P$ and $Q$. The following simple formula describes the pull-back:

\begin{equation}
R^*D_P = D_{R^*P} = D_{P \circ R^*},
\end{equation}

where $R : \mathbb{C}^{m+1} \to \mathbb{C}^{m+1}$ is the lift of $R$.

The degree of a divisor $D = D_P - D_Q$ is $\deg D = \deg P - \deg Q$. On $\mathbb{CP}^2$, Bezout’s Theorem asserts that two divisors $D_1$ and $D_2$ intersect $\deg D_1 \cdot \deg D_2$ times in $\mathbb{CP}^2$, counted with appropriate intersection multiplicities. Suppose that $D_1$ and $D_2$ are irreducible algebraic curves assigned multiplicity one. Then, an intersection point $z$ is assigned multiplicity one if and only if both curves are non-singular at $z$, meeting transversally there. See [Shaf, Ch. IV].

The algebraic degree of a rational map $R : \mathbb{CP}^m \to \mathbb{CP}^m$ (denoted $\deg R$) is the common degree of the homogeneous equations in $R$. Equation (A.3) implies

**Lemma A.2.** Given a dominant rational map $R : \mathbb{CP}^m \to \mathbb{CP}^m$ and a divisor $D$ in $\mathbb{CP}^m$, we have:

$$\deg(R^*D) = \deg R \cdot \deg D.$$
More generally, any divisor $D = D_P - D_Q$ defines a section $\sigma_D$ of the $\deg(D)$-th tensor power of the co-tautological line bundle, defined by $\sigma_D(z) = (\hat{z} : P(\hat{z})/Q(\hat{z}))$. One can recover $D$ from $\sigma_D$ by taking its zero divisor.

A.5. Currents. We will now give a brief background on currents; for more details see [8] and the appendix from [5]. Currents are naturally defined on general complex (or even smooth) manifolds, however to continue our discussion of rational maps, divisors, etc, we restrict our attention to projective manifolds.

A $(1,1)$-current $\nu$ on $V$ is a continuous linear functional on $(m-1, m-1)$-forms with compact support. It can be also defined as a generalized differential $(1,1)$-form $\sum \mu_{ij} dz_i d\bar{z}_j$ with distributional coefficients.

A basic example is the current $[C]$ of integration over the regular points $C_{\text{reg}}$ of an algebraic hypersurface $C$:

$$\omega \mapsto \int_{C_{\text{reg}}} \omega,$$

where $\omega$ is a test $(m-1, m-1)$-form. The current of integration over a divisor $D$ is defined by extending linearly.

The space of currents is given the distributional topology: $\nu_n \to \nu$ if $\nu_n(\omega) \to \nu(\omega)$ for every test form $\omega$.

A differential $(m-1, m-1)$-form $\omega$ is called positive if its integral over any complex subvariety is non-negative. A $(1,1)$-current $\mu$ is called positive if $\mu(\omega) \geq 0$ for any positive $(1,1)$-form. A current $\mu$ is called closed if $d\mu = 0$, where the differential $d$ is understood in the distributional sense.

In this paper, we focus on closed, positive $(1,1)$ currents. They have a simple description in terms of local potentials, rather analogous to the definition of divisors.

Recall that $\partial$ and $\bar{\partial}$ stand for the holomorphic and anti-holomorphic parts of the external differential $d = \partial + \bar{\partial}$. Their composition $\Delta_p = \frac{1}{2} \partial \bar{\partial}$ is the \textit{pluri-Laplacian}. \footnote{Many authors introduce real operators $d = \partial + \bar{\partial}$ and $d^c = \frac{1}{2\pi}(\partial - \bar{\partial})$ and denote the pluri-Laplacian by $dd^c$. We use $\Delta_p$ to avoid confusion between the operator $d$ and the algebraic degree of a map.}

Given a $C^2$-function $h$, the restriction of $\Delta_p h$ to any non-singular complex curve $X$ is equal to the form $\Delta(h|X)dz \wedge d\bar{z}$, where $z$ is a local coordinate on $X$ and $\Delta$ is the usual Laplacian in this coordinate.

If $U$ is an open subset of $\mathbb{C}^n$ and $h : U \to [-\infty, \infty)$ is a plurisubharmonic (PSH) function, then its pluri-Laplacian $\Delta_p h$ is a closed $(1,1)$-current on $U$. Conversely, the $\partial \bar{\partial}$-Poincare Lemma asserts that every closed, positive $(1,1)$-current on $U$ is obtained this way.

Therefore, any closed positive $(1,1)$ current $\nu$ on a manifold $V$ can be described using an open cover $\{U_i\}$ of $V$ together with PSH functions $v_i : U_i \to [-\infty, \infty)$ that are chosen so that $\nu = \Delta_p v_i$ in each $U_i$. The functions $v_i$ are called \textit{local potentials} for $\nu$ and they are required to satisfy the compatibility condition that $v_i - v_j$ is pluriharmonic (PH) on any non-empty intersection $U_i \cap U_j \neq \emptyset$. The support of $\nu$ is defined by:

$$\text{supp } \nu := \{ z \in V : \text{if } z \in U_j \text{ then } v_j \text{ is not pluriharmonic at } z \}.$$
have \([D] = \Delta_p \log |g_i|\). The result is closed \((1, 1)\) current, which is positive iff \(D\) is effective (i.e. the multiplicities \(k_1, \ldots, k_r\) are non-negative).

Suppose \(R : V \to W\) is a dominant rational map and \(\nu\) is a closed-positive \((1, 1)\) current on \(W\). The pullback \(R^*\nu\) is closed positive \((1, 1)\) current on \(V\), defined as follows. First, one obtains a closed positive \((1, 1)\) current \(R^*\nu\) defined on \(V \setminus I(R)\) by pulling-back the system of local potentials defining \(\nu\) under \(R : V \setminus I(R) \to W\). One then extends \(R^*\nu\) trivially through \(I(R)\), to obtain a closed, positive \((1, 1)\) current defined on all of \(V\). (By a result of Harvey and Polking [HaPo], this extension is closed.) See [Si, Appendix A.7] for further details. Pullback is continuous with respect to the distributional topology.

Similarly to divisors, there is a particularly convenient description of closed, positive \((1, 1)\) currents on \(\mathbb{C}P^m\). Associated to any PSH function \(H : \mathbb{C}^{m+1} \to [-\infty, \infty)\), having the homogeneity (for some \(c > 0\)) that

\[
H(\lambda z) = c \log |\lambda| + H(z),
\]

is a closed, positive \((1, 1)\) current, denoted by \(\pi_* (\Delta_p H)\), given by the system of local potentials \(\{U_i, H \circ \sigma_i\}\), where the \(\{U_i\}\) form an open covering of \(\mathbb{C}P^m\) that admits local sections \(\sigma_i : U_i \to \mathbb{C}^{m+1} \setminus \{0\}\) of the canonical projection \(\pi\). (In each \(U_i\), it is defined by \(\pi_* (\Delta_p H) = \Delta_p (H \circ \sigma_i)\) Moreover, every closed positive \((1, 1)\) current on \(\mathbb{C}P^m\) is described in this way; See [Si Thm A.5.1]. The function \(H\) is called the pluripotential of \(\pi_* (\Delta_p H)\).

If \(R : \mathbb{C}P^m \to \mathbb{C}P^m\) is a rational map, the action of pull-back is described by

\[
R^* \pi_* (\Delta_p H) = \pi_* (\Delta_p (H \circ R)).
\]

**Remark A.1.** One can imagine adapting the description of divisors on \(\mathbb{C}P^m\) in terms of sections of (tensor powers) of the co-tautological bundle to develop a description of closed positive \((1, 1)\)-currents on \(\mathbb{C}P^m\) by using PSH sections of appropriate real line bundles. However, since \(\mathbb{C}P^m\) is simply-connected, every real line bundle is trivial. Hence, this approach does not really provide any further insights.

**A.6. Kobayashi hyperbolicity and normal families.** In [4] we use the Kobayashi metric in order to prove that the iterates \(R^n\) form a normal family on certain subspaces of \(\mathbb{C}P^2\). Here we recall the relevant definitions and some important results that we use. The reader can consult the books [K, La] and the original papers by M. Green [G1, G2] for more details. For more dynamical applications, see e.g. [Si].

The Kobayashi pseudometric is a natural generalization of the Poincaré metric on Riemann surfaces. Let \(\| \cdot \|\) stand for the Poincaré metric on the unit disk \(\mathbb{D}\). Let \(M\) be a complex manifold. Pick a a tangent vector \(\xi \in TM\), and let \(\mathcal{H}(\xi)\) be a family of holomorphic curves \(\gamma : \mathbb{D} \to M\) tangent to the line \(\mathbb{C} \cdot \xi\) at \(\gamma(0)\). Then \(Df(v) = \xi\) for some \(v \equiv v_\gamma \in T_0\mathbb{D}\), and the Kobayashi pseudometric is defined to be:

\[
ds_M(\xi) = \inf_{\gamma \in \mathcal{H}(\xi)} \|v_\gamma\|.
\]

The Kobayashi pseudometric is designed so that holomorphic maps are distance decreasing: if \(f : U \to M\) is holomorphic then \(ds_M(Df(\xi)) \leq ds_U(\xi)\).

The reason for “pseudo-” is that for certain complex manifolds \(M\), \(ds(\xi)\) can vanish for some non-vanishing tangent vectors \(\xi \neq 0\). For example, \(ds\) identically vanishes on \(\mathbb{C}^n\) or \(\mathbb{C}P^n\). A complex manifold \(M\) is called **Kobayashi hyperbolic** if \(ds\) is non-degenerate: \(ds(\xi) > 0\) for any non-vanishing \(\xi \in TM\). Then it induces a (Finsler) metric on \(M\).
Let $N$ be a compact complex manifold. Endow it with some Hermitian metric $|\cdot|_N$. A complex submanifold $M \subset N$ is called hyperbolically embedded in $N$ if the Kobayashi pseudometric on $M$ dominates the Hermitian metric on $N$, i.e., there exists $c > 0$ such that $ds_M(\xi) \geq c|\xi|_N$ for all $\xi \in TM$. Obviously, $M$ is Kobayashi hyperbolic in this case.

A complex manifold $M$ is called Brody hyperbolic if there are no non-constant holomorphic mappings $f: \mathbb{C} \to M$. If $M$ is Kobayashi hyperbolic, it is also Brody hyperbolic, but the converse is generally not true unless $M$ is compact.

An open subset of $\mathbb{CP}^2$ that is Kobayashi hyperbolic, but not hyperbolically embedded in $\mathbb{CP}^2$ is described in [K, Example 3.3.11] and an open subset of $\mathbb{C}^2$ that is Brody hyperbolic but not Kobayashi hyperbolic is described in [K, Example 3.6.6].

A family $\mathcal{F}$ of holomorphic mappings from a complex manifold $U$ to a complex manifold $M$ is called normal if every sequence in $\mathcal{F}$ either has a subsequence converging locally uniformly or a subsequence that diverges locally uniformly to infinity in $M$. In the case that $M$ is embedded into some compact manifold $Z$, a stronger condition is that $\mathcal{F}$ is precompact in $\text{Hol}(U, Z)$ (where $\text{Hol}(U, Z)$ is the space of holomorphic mappings $U \to Z$ endowed with topology of uniform convergence on compact subsets of $U$).

Proposition A.3. Let $M$ be a hyperbolically embedded complex submanifold of a compact complex manifold $N$. Then for any complex manifold $U$, the family $\text{Hol}(U, M)$ is precompact in $\text{Hol}(U, N)$.

See Theorem 5.1.11 from [K].

The classical Montel’s Theorem asserts that the family of holomorphic maps $\mathbb{D} \to \mathbb{C}\setminus\{0, 1\}$ is normal (as $\mathbb{C}\setminus\{0, 1\}$ is a hyperbolic Riemann surface). It is a foundation for the whole Fatou-Julia iteration theory. Several higher dimensional versions of Montel’s Theorem, due to M. Green [G1, G2], are now available. Though their role in dynamics is not yet so prominent, they have found a number of interesting applications. Below we will formulate two particular results used in this paper (see [G1]). The following is Theorem 2 from [G1].

Theorem A.4. Let $X$ be a union of (possibly singular) hypersurfaces $X_1, \ldots, X_m$ in a compact complex manifold $N$. Assume $N \setminus X$ is Brody hyperbolic and

$$X_{i_1} \cap \cdots \cap X_{i_k} \setminus (X_{j_1} \cup \cdots X_{j_l})$$

is Brody hyperbolic for any choice of distinct multi-indices $\{i_1, \ldots, i_k, j_1, \ldots, j_l\} = \{1, \ldots, m\}$. Then $N \setminus X$ is a complete hyperbolically embedded submanifold of $N$.

In the last section of [G1], the following result is proved:

Theorem A.5. Let $M = \mathbb{CP}^2 \setminus (Q \cup X_1 \cup X_2 \cup X_3)$, where $Q$ is a non-singular conic and $X_1, X_2, X_3$ are lines. Then any non-constant holomorphic curve $f: \mathbb{C} \to M$ must lie in a line $L$ that is tangent to $Q$ at an intersection point with one of the lines, $X_i$, and that contains the intersection point $X_j \cap X_l$ of the other two lines.

The configurations that appear in this theorem are related to amusing projective triangles:
A.7. Self-dual triangles. Let \( Q(z) = \sum q_{ij} z_i z_j \) be a non-degenerate quadratic form in \( E \cong \mathbb{C}^3 \), and \( X = \{ Q = 0 \} \) be the corresponding conic in \( \mathbb{CP}^2 \). The form \( Q \) makes the space \( E \) Euclidean, inducing duality between points and lines in \( \mathbb{CP}^2 \). Namely, to a point \( z = (z_0 : z_1 : z_2) \) corresponds the line \( L_z = \{ \zeta : Q(z, \zeta) = 0 \} \) called the polar of \( z \) with respect to \( X \) (here we use the same notation for the quadratic form and the corresponding inner product). Geometrically, this duality looks as follows. Given a point \( z \in \mathbb{CP}^2 \), there are two tangent lines from \( z \) to \( X \). Then \( L_z \) is the line passing through the corresponding tangency points. (In case \( z \in X \), the polar is tangent to \( X \) at \( z \)).

Three points \( z_i \) in \( \mathbb{CP}^2 \) in general position are called a “triangle” \( \Delta \) with vertices \( z_i \). Equivalently, a triangle can be given by three lines \( L_i \) in general position, its “sides”. Let us say that \( \Delta \) is self-dual (with respect to the conic \( X \)) if its vertices are dual to the opposite sides.

**Lemma A.6.** A triangle \( \Delta \) with vertices \( z_i \) is self-dual if and only if the corresponding vectors \( \hat{z}_i \in E \) form an orthogonal basis with respect to the inner product \( Q \).

All three sides of a self-dual triangle satisfy the condition of **Theorem A.5** so they can give us exceptional holomorphic curves \( \mathbb{C} \to \mathbb{CP}^2 \smallsetminus (Q \cup X_1 \cup X_2 \cup X_3) \).

**Remark A.2.** In a similar way one can define self-dual tetrahedra in higher dimensions.

**Appendix B. Elements of complex dynamics in several variables**

We now provide a brief background on the dynamics of rational maps in several variables. We refer the reader to the survey [Si] for more details.

**B.1. Algebraic Stability.** The following statement appears in [Si] Prop. 1.4.3:

**Lemma B.1.** Let \( R \) and \( S \) be two rational maps \( \mathbb{CP}^m \to \mathbb{CP}^m \). Then, \( \deg(S \circ R) = \deg(S) \cdot \deg(R) \) if and only if there is no algebraic hypersurface \( V \subset \mathbb{CP}^m \) that is collapsed by \( R \) to an indeterminate point of \( S \).

A rational mapping \( R : \mathbb{CP}^m \to \mathbb{CP}^m \) is called algebraically stable if there is no integer \( n \) and no collapsing hypersurface \( V \subset \mathbb{CP}^m \) so that \( R^n(V) \) is contained within the indeterminacy set of \( R \), [Si] p. 109]. A consequence of **Lemma B.1** is that \( R \) is algebraically stable if and only if \( \deg R^n = (\deg R)^n \).

A direct consequence of **Lemma A.2** is:

**Lemma B.2.** If \( R \) is an algebraically stable map, then for any divisor \( D \) we have:

\[
\text{(B.1) } \deg((R^n)^* D) = (\deg R)^n \cdot \deg D
\]

**B.2. Green potential.**

**Theorem B.3** (see [Si] Thm 1.6.1). Let \( R : \mathbb{CP}^m \to \mathbb{CP}^m \) be an algebraically stable rational map of degree \( d \). Then the limit

\[
G = \lim_{n \to \infty} \frac{1}{d^n} \log \| \hat{R}^n \|
\]

exists in \( L^1_{\text{loc}}(\mathbb{C}^3) \) and determines a plurisubharmonic function. This function satisfies the following equivariance properties:

\[
\text{(B.2) } G(\lambda z) = G(z) + \log |\lambda|, \quad \lambda \in \mathbb{C}^*,
\]
\[ G \circ \hat{R} = dG. \]

It is called the Green potential of \( R \).

B.3. **Green current.** Applying the pluri-Laplacian \( \Delta_p \) to the Green potential, we obtain:

**Theorem B.4** (see [Si], Thm 1.6.1). Let \( R : \mathbb{CP}^m \to \mathbb{CP}^m \) be an algebraically stable rational map of degree \( d \). Then \( m^c = \pi_*(\Delta_p G) \) is a closed positive \((1,1)\)-current on \( \mathbb{CP}^2 \) satisfying the equivariance relation: \( R^*m^c = d \cdot m^c \).

The current \( m^c \) is called the Green current of \( R \).

The set of regular points for an algebraically stable rational map \( R : \mathbb{CP}^m \to \mathbb{CP}^m \) is:

\[ N := \{ x \in \mathbb{CP}^m : \text{there exits neighborhoods } U \text{ of } x \text{ and } V \text{ of } I(R) \text{ so that } f^n(U) \cap V = \emptyset \text{ for every } n \in \mathbb{N} \}. \]

The regular points form an open subset of \( \mathbb{CP}^m \).

One primary interest in the Green current \( m^c \) is the following connection between its support \( \text{supp } m^c \) and the Julia set \( J_R \). (See §4 for the definitions of the Fatou and Julia sets.) Note that \( \text{supp } m^c \) is closed and backwards invariant, \( R^{-1} \text{supp } m^c \subset \text{supp } m^c \), since \( R^*m^c = d \cdot m^c \).

**Theorem B.5** (see [Si], Thm 1.6.5). Let \( f : \mathbb{CP}^m \to \mathbb{CP}^m \) be an algebraically stable rational map. Then:

\[ J_R \cap N \subset \text{supp } m^c \subset J_R. \]

An algebraically stable rational map \( R : \mathbb{CP}^2 \to \mathbb{CP}^2 \) for which \( \text{supp } m^c \subset J_R \) is given in [FS1] Example 2.1.

**Appendix C. Complex Whitney folds**

To simplify calculations near the critical points, it is convenient to bring \( R \) to a normal form. A complex Whitney fold is generic and the simplest one (see [AGV]). Let \( R : (\mathbb{C}^2,0) \to (\mathbb{C}^2,0) \) be a germ of holomorphic map with a critical point at 0.

The map \( R \) (and the corresponding critical set) is called a complex Whitney fold if:

(W1) The critical set \( L \) is a non-singular curve near 0;
(W2) \( DR(0) \) has rank 1 and \( \text{Ker } DR(0) \) is transverse to \( L \);
(W3) The second differential \( D^2R(0) \) is not vanishing in the direction of \( \text{Ker } DR(0) \).

The following is a standard result from singularity theory:

**Lemma C.1.** A Whitney fold can be locally brought to a normal form \((u, w) \mapsto (u, w^2)\) in holomorphic coordinates.

We now consider how volume is transformed under a mapping near a Whitney fold. Let us begin with the 1D power map:

**Lemma C.2.** Let \( Q : w \mapsto w^d \). For any measurable set \( Y \subset \mathbb{C} \),

\[ \text{area}(Q^{-1}X) \leq \left( \text{area } X \right)^{1/d}. \]

---

8They are called normal points in [Si], but we prefer regular to avoid confusion with the notion of normal families.
Proof. Of course, we can assume that area $X > 0$. Let us take the radius $r > 0$ such that $\pi r^2 = \text{area} X$. Let $X_+ = X \cap \mathbb{D}_r$, $X_+ = X \setminus X_-$, $X_+ = \mathbb{D}_r \setminus X_-$. Then

$$(\text{area} X)^{1/d} = (\text{area} \mathbb{D}_r)^{1/d} = \text{area}(Q^{-1} \mathbb{D}_r) = \text{area}(Q^{-1} X_-) + \text{area}(Q^{-1} X_+) \geq \text{area}(Q^{-1} X_-) + d \text{Jac} Q^{-1}(r) \text{area} X \geq \text{area}(Q^{-1} X_-) + \text{area}(Q^{-1} X_+).$$

□

Corollary C.3. Let $L$ be a complex Whitney fold for a map $R$, and let $L' \subset L$. Then for any measurable set $Y$ sufficiently close to $R(L')$, we have:

$$\text{vol}(R^{-1} X) \leq C \sqrt{\text{vol} X},$$

where $R^{-1}$ is the one-to-two branch of the inverse map associated to the fold, and $C = C(R, L, L')$.

Proof. As the assertion is local, we can use the Whitney normal coordinates $(u, w)$ near $L'$ (Lemma C.1). Let $X^h$ be the projection of $X$ (in this coordinates) onto the $u$-axis $L$, and let $X^v(u)$ be the slice of $X$ be the vertical complex line through $(u, 0) \subset L$. Then

$$\text{vol}(R^{-1} X) \leq \int_{X^h} \text{area}(Q^{-1}(X^v(u))) d \text{area}(u)$$

$$\leq \int_{X^h} \sqrt{\text{area} X^v(u)} d \text{area}(u) \leq \sqrt{\text{area} X^h \cdot \text{vol} X},$$

where the estimates follow from Fubini, Lemma C.2 and Cauchy-Schwarz respectively. □

Appendix D. Open Problems

Problem D.1 (Existence of Fisher and Lee-Yang-Fisher distributions). Consider a sequence of graphs $\Gamma_n$ for which van Hove’s Theorem [vH, R2] justifies existence of the limiting free energy (2.5). Under what circumstances does the limit (2.7) exist? As explained in § 2.2, this would justify consideration of the limiting distributions of Lee-Yang-Fisher zeros for more general lattices.

For $\mathbb{Z}^d$, existence of the limit (2.7) seems to be open if $d \geq 2$.

Problem D.2 (Geometric properties of the Lee-Yang-Fisher current). The theory of geometric currents has become increasingly useful in complex dynamics, see [RS, BLS, Di, dT, DDG] as a sample.

The Green current $m^c$ is strongly laminar in a neighborhood of $B$. The structure is given by the stable lamination of $B$ (see § 3.3) together with transverse measure obtained under holomomy from the Lebesgue measure on $B$. However, $m^c$ is not strongly laminar in a neighborhood of the topless Lee-Yang “cylinder” $C_1$.

One can see this also follows: a disc within the invariant line $L_{\text{inv}}$ centered at $L_{\text{inv}} \cap B$ is within the stable lamination of $B$. Therefore, an open neighborhood within $L_{\text{inv}}$ of $L_{\text{inv}} \cap C_1$ would have to be a leaf of the laminations. However, $m^c$ restricts to $L_{\text{inv}}$ in a highly non-trivial way, coinciding with the measure of maximal entropy for $R|L_{\text{inv}}$. (It is supported on the Julia set shown in Figure 1.3).

Does $m^c$ have a weaker geometric structure? For example, is it non-uniformly laminar [BLS, Di] or woven [Di, dT, DDG]?
Problem D.3 (Support for the measure of maximal entropy). What can be said about the support of the measure of maximal entropy $\nu$ that was discussed in §4.3? Is the critical fixed point $b_c \in L_{\text{inv}}$ within $\text{supp}\, \nu$? A positive answer to this question is actually equivalent to $C \cap \text{supp}\, \nu \neq \emptyset$ and also to $\text{supp}\, \nu \cap L_{\text{inv}} = J_R|L_{\text{inv}}$.

Problem D.4 (Fatou Set). In Theorem 4.7 we showed that certain “solid cylinders” are in $W^s(e)$ and $W^s(e')$. Computer experiments suggest a much stronger result:

Conjecture. $W^s(e) \cup W^s(e')$ is the entire Fatou set for $R$.

Problem D.5 (Julia Set). Proposition 4.2 gives that in a neighborhood of $B$, $J_R$ is a $C^\infty$ 3-manifold. What can be said about the global topology of $J_R$?

Remark D.1. Note that each of the above problems D.2 – D.5 has a natural counterpart for $\mathcal{R}$.

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