# A degenerate Newton's Map in two complex variables: linking with currents 

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August 16, 2006


#### Abstract

Little is known about the global structure of the basins of attraction of Newton's method in two or more complex variables. We make the first steps by focusing on the specific Newton mapping to solve for the common roots of $P(x, y)=x(1-x)$ and $Q(x, y)=y^{2}+B x y-y$.

There are invariant circles $S_{0}$ and $S_{1}$ within the lines $x=0$ and $x=1$ which are superattracting in the $x$-direction and hyperbolically repelling within the vertical line. We show that $S_{0}$ and $S_{1}$ have local super-stable manifolds, which when pulled back under iterates of $N$ form global super-stable spaces $W_{0}$ and $W_{1}$. By blowing-up the points of indeterminacy $p$ and $q$ of $N$ and all of their inverse images under $N$ we prove that $W_{0}$ and $W_{1}$ are real-analytic varieties.

We define linking between closed 1-cycles in $W_{i}(i=0,1)$ and an appropriate closed 2 current providing a homomorphism $l k: H_{1}\left(W_{i}, \mathbb{Z}\right) \rightarrow \mathbb{Q}$. If $W_{i}$ intersects the critical value locus of $N$, this homomorphism has dense image, proving that $H_{1}\left(W_{i}, \mathbb{Z}\right)$ is infinitely generated. Using the Mayer-Vietoris exact sequence and an algebraic trick, we show that the same is true for the closures of the basins of the roots $\overline{W\left(r_{i}\right)}$.


Key Words. Complex dynamics, Newton's Method, homology, linking numbers, invariant currents. 2000 Mathematics Subject Classification. 37F20, 32Q55, 32H50, 58K15.

Newton's method is one of the fundamental algorithms of mathematics, so it is evidently important to understand its dynamics, particularly the structure of the basins of attraction of the roots. Even in one dimension, the topology of these basins can be complicated and there has been a good deal of research on this subject. In higher dimensions, next to nothing is known about the topology of the basins. In this paper we make the first steps at understanding their topology in two complex variables.

We focus on a specific system: the Newton's Method used to solve for the common roots of $P(x, y)=x(1-x)$ and $Q(x, y)=y^{2}+B x y-y$. While this is one specific and relatively simple system, we believe that some of the techniques developed in this paper can be used to study more general systems.

Dynamical systems $g: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ are often classified in terms of: (1) The number of inverse images of a generic point by $g$, which is called the topological degree $d_{t}(g)$, and (2) Whether $g$ has points of indeterminacy.

Mappings $g: \mathbb{P}^{n} \rightarrow \mathbb{P}^{n}$ with $d_{t}(g)>1$, but without points of indeterminacy, have been studied by Bonifant and Dabija [10], Bonifant and Fornæss [11], Briend [14], Briend and Duval [15], Dinh and Sibony [20], Fornaess and Sibony [25, 27, 26], Hubbard and Papadopol [38], Jonnson [39], and Ueda [52].

Meanwhile, birational maps $g: \mathbb{P}^{n} \rightsquigarrow \mathbb{P}^{n}$ (rational maps with rational inverse) are examples of systems with points of indeterminacy, but with $d_{t}(g)=1$. The famous Henon mappings $H$ : $\mathbb{P}^{2} \rightsquigarrow \mathbb{P}^{2}$ fall into this class. Such systems have been studied extensively by Bedford and Smillie $[3,4,5,6,7,9,8]$, Bedford, Lyubich and Smillie [2], Devaney and Nitecki [18], Diller [19], Dinh


Figure 1: Is the loop $\gamma$ trivial in the homology of it's basin, within $\mathbb{C}^{2}$ ?
and Sibony [21], Dujardin [22], Favre and Jonsson [23], Fornaess [24], Guedj [29], and Hubbard and Oberste-Vorth [34, 35, 36].

Not nearly as much is known about mappings $g: \mathbb{P}^{n} \rightsquigarrow \mathbb{P}^{n}$ with topological degree $d_{t}(g)>1$ and with points of indeterminacy. The work of Russakovskii and Shiffman [46] considers a measure that is obtained by choosing a generic point, taking the each of its inverse images under $g^{\circ n}$ and giving them all equal weight in order to obtain a probability measure $\mu_{n}$. Under appropriate conditions on $g$ they show that the measures $\mu_{n}$ converge to a measure $\mu$ that is independent of the initial point. In [37], the authors present a proof by A. Douady that $\mu$ does not charge points in the line at infinity, a result not obtained in [46]. In a recent paper, Guedj [30] shows that if the topological degree $d_{t}(g)$ is sufficiently large, then $\mu$ does not charge the points of indeterminacy of $g$ and does not charge any pluripolar set. He then uses these facts to establish ergodic properties of $\mu$.

Many of the papers considering mappings with both $d_{t}(g)>1$ and points of indeterminacy consider ergodic properties, invariant measures, and invariant currents, focusing less on topological properties. One paper that considers some topological properties is [37], by John Hubbard and Peter Papadopol, who consider the dynamics of the Newton Map $N$ to solve for the zeros of two quadratic equations $P$ and $Q$ in two complex variables. The basins of attraction for this system show interesting topology: for example, when drawing intersections of a 1 complex-dimensional slice with the basins of attraction one often finds "bubbles" like the ones shown in Figure 1. It is natural to ask if a loop, such as the one labeled $\gamma$ in the figure, corresponds to a non-trivial loop in the homology of its basin of attraction. Clearly $[\gamma]$ is non-trivial in the basin intersected with this slice, but it is much more difficult to determine if $[\gamma]$ is non-trivial when considered within the entire 2 complex-dimensional basin, which may reconnect in unusual ways outside of this slice.

Questions about the first homology of the basins are not answered by Hubbard and Papadopol. Using general principles they show that the basin of attraction for each of the four roots is path connected and, by resolving points of indeterminacy of $N$, they show that each basin is a Stein manifold. But, instead of addressing further questions about the topology of the individual basins, they focus on creating a compactification with manageable topology on which $N$ is well-defined. Most questions about the topology and the detailed structure of these basins of attraction within their compactification of $\mathbb{C}^{2}$ remain as mysteries.

In order to develop tools to answer some detailed questions about the topology of basins of attraction for Newton maps $N: \mathbb{C}^{2} \rightarrow \mathbb{C}^{2}$ to solve two quadratic equations $P(x, y)=0$ and
$Q(x, y)=0$, we restrict our attention to the degenerate case when the four roots lie on two parallel lines. Normalizing, we study the Newton Map $N$ for $P(x, y)=x(1-x)$ and $Q(x, y)=y^{2}+B x y-y$. In this case, the first component of $N(x, y)$ depends only on $x$, while the second component depends on both $x$ and $y$. Systems of this form are commonly referred to as skew products in the literature $[1,32,39,47,50,51]$ and they are often used as "test cases" when developing new techniques. While we rely upon the fact that $N$ becomes a skew product in this degenerate case, we hope that some of the techniques developed here can eventually be adapted to non-degenerate cases.

Our approach is the following:
We compactify $\mathbb{C}^{2}$ obtaining a rational map $N: \mathbb{P} \times \mathbb{P} \rightarrow \mathbb{P} \times \mathbb{P}$ with four points of indeterminacy at $p=\left(\frac{1}{B}, 0\right), q=\left(\frac{1}{2-B}, \frac{1-B}{2-B}\right),(\infty, \infty)$, and $\left(\infty, \frac{B}{2}\right)$. There are three invariant subspaces $X_{l}:=$ $\{(x, y): \operatorname{Re}(x)<1 / 2$ and $x \neq \infty\}, X_{1 / 2}:=\{(x, y): \operatorname{Re}(x)=1 / 2$ or $x=\infty\}$ and $X_{r}:=\{(x, y) \in$ $: \operatorname{Re}(x)>1 / 2$ and $x \neq \infty\}$. The common roots of $P$ and $Q$ are $r_{1}=(0,0), r_{2}=(0,1), r_{3}=(1,0)$, and $r_{4}=(1,1-B)$ with the basins of attraction $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ in $X_{l}$ and the basins of attraction $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$ in $X_{r}$. By restricting to parameters $B \in \Omega=\{B:|1-B|>1\}$ we can assume that both $p$ and $q$ are in $X_{l}$.

We will prove that within $X_{l}$ and $X_{r}$ there are "superstable separatrices" $W_{0}$ and $W_{1}$ consisting of the points that are attracted to invariant circles within the lines $x=0$ and $x=1$ respectively. By resolving the points of indeterminacy of $N^{k}$ in $X_{l}$ we obtain a modified space $X_{l}^{\infty}$ in which all iterates of $N$ are well-defined and in which $W_{0}$ is a real-analytic variety that provides a nice boundary between $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$. Since $B \in \Omega$ there is no such problem in $X_{r}$ : all iterates of $N$ are already well defined on $X_{r}$ and $W_{1}$ is a real-analytic variety in $X_{r}$.

In this paper we will study the topology of $W_{0}$ and $W_{1}$ in detail and we will use a Mayer-Vietoris decomposition to relate their homology to the homology of the basins of attraction of the four roots: $W\left(r_{1}\right), W\left(r_{2}\right), W\left(r_{3}\right)$, and $W\left(r_{4}\right)$ and the homology of $X_{l}^{\infty}$ and $X_{r}$.

The major emphasis of this paper is to show that loops in $W_{0}$ and $W_{1}$ that are generated by intersections of $W_{0}$ or $W_{1}$ with the critical value locus $C$ are actually homologically non-trivial. The essential difficulty is to choose a notion of linking that is well defined within the space $X_{l}^{\infty}$, which is very topologically complicated as a result of the blow-ups.

We define linking between closed 1-cycles in $W_{i}(i=0,1)$ and an appropriate closed 2 current providing a homomorphism $l k: H_{1}\left(W_{i}, \mathbb{Z}\right) \rightarrow \mathbb{Q}$. If $W_{i}$ intersects the critical value locus of $N$, this homomorphism has dense image, proving that $H_{1}\left(W_{i}, \mathbb{Z}\right)$ is infinitely generated. Using the MayerVietoris exact sequence and an algebraic trick, we show that the same is true for the closures of the basins of the roots $\overline{W\left(r_{i}\right)}$.

Our work culminates to prove:
Theorem 0.1. Let $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ be the closures in $X_{l}^{\infty}$ of the basins of attraction of $r_{1}=(0,0)$ and $r_{2}=(0,1)$ under iteration of $N$ and let $\overline{W\left(r_{3}\right)}$ and $\overline{W\left(r_{4}\right)}$ be the closures in $X_{r}$ of the basins of attraction of $r_{3}=(1,0)$ and $r_{4}=(1,1-B)$.

- $H_{1}\left(\overline{W\left(r_{1}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{2}\right)}\right)$ are infinitely generated for every $B \in \Omega$.
- For $B \in \Omega$, if $W_{1}$ intersects the critical value parabola $C(x, y)=0$, then both $H_{1}\left(\overline{W\left(r_{3}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{4}\right)}\right)$ are infinitely generated, otherwise $H_{1}\left(\overline{W\left(r_{3}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{4}\right)}\right)$ are trivial.

For $B \in \Omega_{\mathrm{reg}}$, the set of parameters for which the separatrices are genuine manifolds, the basins of the four roots and their closures in $X_{l}^{\infty}$ and $X_{r}$ have the some homotopy type. Hence:
Corollary 0.2. For $B \in \Omega_{\mathrm{reg}}$, Theorem 0.1 remains true when replacing the closures of each of the basins with the basins themselves.

## 1 Basic properties of $N$

In the first part of this section we summarize the basic results from [37].
Given two vector spaces $V$ and $W$ of the same dimension and a differentiable mapping $F: V \rightarrow$ $W$, the associated Newton map $N_{F}: V \rightarrow V$ is given by the formula

$$
\begin{equation*}
N_{F}(\mathbf{x})=\mathbf{x}-[D F(\mathbf{x})]^{-1}(F(\mathbf{x})) \tag{1}
\end{equation*}
$$

If $D F\left(r_{i}\right)$ is invertible for each root $r_{i}$ of $F$, then the roots of $F$ correspond to super attracting fixed points of $N_{F}$. Conversely, every fixed point of $N_{F}$ is a root of $F$. Since each fixed point $r_{i}$ of $N_{F}$ is super-attracting, there is some neighborhood $U_{i}$ of $r_{i}$ for which each initial guess $\mathbf{x}_{\mathbf{0}} \in U_{i}$ will converge to $r_{i}$. An explicit lower bound on the size of $U_{i}$ is given by Kantorovich's Theorem [40].

Proposition 1.1. (Transformation rules) If $A: V \rightarrow V$ is affine, and invertible, and if $L: W \rightarrow$ $W$ is linear and invertible, then:

$$
\begin{equation*}
N_{L \circ F \circ A}=A^{-1} \circ N_{F} \circ A \tag{2}
\end{equation*}
$$

The proof is a careful use of the chain rule, see [37], Lemma 1.1.1.
Proposition 1.2. Newton's Method to find the intersection of two quadratics depends only on the intersection points and not on the choice of curves.

For the proof, see Corollary 1.5.2 from [37].
In this paper, we normalize so that the roots are at $\binom{0}{0},\binom{1}{0},\binom{0}{1}$, and $\binom{\alpha}{\beta}$. If we let $A=\frac{1-\alpha}{\beta}$ and $B=\frac{1-\beta}{\alpha}$, then $F\binom{x}{y}=\binom{x^{2}+A x y-x}{y^{2}+B x y-y}=\binom{P(x, y)}{Q(x, y)}$ has these roots and the corresponding Newton Map is given by:

$$
\begin{align*}
N_{F}\binom{x}{y} & =\binom{x}{y}-\left[\begin{array}{cc}
2 x+A y-1 & A x \\
B y & 2 y+B x-1
\end{array}\right]^{-1}\binom{x^{2}+A x y-x}{y^{2}+B x y-y} \\
& =\frac{1}{\Delta}\binom{x\left(B x^{2}+2 x y+A y^{2}-x-A y\right)}{y\left(B x^{2}+2 x y+A y^{2}-B x-y\right)} \tag{3}
\end{align*}
$$

where $\Delta=2 B x^{2}+4 x y+2 A y^{2}-(2+B) x-(2+A) y+1$.
Proposition 1.3. The critical value locus of $N_{F}$ is the union of the two parabolas that go through the four roots of $F$.

Proposition 1.4. The Newton Map has topological degree 4.
See [37] for a proof of Propositions 1.3 and 1.4.
It is a classical result that the dynamics of the Newton map $N(z)$ to solve for the roots of any quadratic polynomial $p(z)$ is conjugate to the map $z \mapsto z^{2}$. For the latter, the unit circle $\mathbb{S}^{1}$ forms the boundary between the basin of attraction of 0 and of $\infty$. If $\phi$ is the map conjugating $N(z)$ to $z \mapsto z^{2}$, then $\phi^{-1}\left(\mathbb{S}^{1}\right)$ is the line in $\mathbb{C}$ that is equidistant from the roots of $p$. This line forms the boundary between the basin of the two roots of $p(z)$ and the dynamics on this line (once you add a point at infinity) are conjugate to angle doubling on the unit circle.

Proposition 1.5. (Invariant lines and invariant circles) The lines joining the roots of $F$ are invariant under Newton Map $N_{F}$ and on these lines $N_{F}$ induces the dynamics of the one dimensional Newton Map to find the roots of a quadratic polynomial.

Within each line is an invariant "circle," corresponding to the points of equal distance from the two roots in that line.
(See Proposition 1.5.3 in [37])
Proof: Given any pair of roots of $F$, there is an affine mapping taking them to $\binom{0}{0}$ and $\binom{1}{0}$ and a third root to $\binom{0}{1}$ The new system is also normalized, but with the chosen pair of roots on the $x$-axis. Using Proposition 1.1, if the $x$-axis is invariant, then we will have shown that the line connecting the chosen pair of roots is also invariant. But this is easy to see because there is a factor of $y$ in the second coordinate of Equation 3 for $N_{F}$.

The dynamics on the $x$-axis correspond to taking the first coordinate of $N_{F}$ in Equation 3 with $y=0$. One finds $x \mapsto \frac{x\left(B x^{2}-x\right)}{2 B x^{2}-(2+B) x+1}=\frac{x^{2}}{2 x-1}$. This is the Newton's Method to solve $x(1-x)=0$. Using the transformation rules from Proposition 1.1 one can show that the same is true for any other invariant line.

### 1.1 The degenerate case: $A=0$

The Newton map to find the common zeros of $P(x, y)=x(1-x)$ and $Q(x, y)=y^{2}+B x y-y$ is:

$$
\begin{equation*}
N\binom{x}{y}=\frac{1}{\Delta}\binom{x\left(B x^{2}+2 x y-x\right)}{y\left(B x^{2}+2 x y-B x-y\right)}=\binom{\frac{x^{2}}{2 x-1}}{\frac{y\left(B x^{2}+2 x y-B x-y\right)}{(2 x-1)(B x+2 y-1)}} \tag{4}
\end{equation*}
$$

with

$$
\Delta=2 B x^{2}+4 x y-(2+B) x-2 y+1=(2 x-1)(B x+2 y-1)
$$

The fixed points of $N$ are the four common roots of $P$ and $Q: r_{1}=(0,0), r_{2}=(0,1), r_{3}=(1,0)$, and $r_{4}=(1,1-B)$.

The critical value locus is the union of the two parabolas going through the four roots. One of these coincides with $P(x, y)=x(1-x)=0$, while the other is the non-degenerate parabola given by

$$
\begin{equation*}
C(x, y)=y^{2}+B x y+\frac{B^{2}}{4} x^{2}-\frac{B^{2}}{4} x-y=0 \tag{5}
\end{equation*}
$$

We will often refer to the locus $C(x, y)=0$ as the critical value parabola and denote it by $C$. Figure 2 depicts the curves $P(x, y)=0$ and $Q(x, y)=0$, the critical value parabola $C$, and the four roots, $r_{1}, r_{2}, r_{3}$, and $r_{4}$.


Figure 2: The degenerate case $A=0$.
One can check directly from Equation 4 that $N$ has topological degree 4 , since every $x \neq 0,1$ has two inverse images and the second component is an equation of degree two in $y$.

There are six invariant lines and, in this degenerate case, these lines have six points of intersection in $\mathbb{C}^{2}$. Four of these intersections correspond to the roots $r_{1}, r_{2}, r_{3}$, and $r_{4}$, while the remaining two correspond to points of indeterminacy. These are denoted $p$ and $q$ and are also shown in Figure 2.

The mapping governing the $x$ coordinate is $x \mapsto \frac{x^{2}}{2 x-1}$, which is itself the one variable Newton Map corresponding to the polynomial $x(x-1$ ), with Julia set consisting of the line $\operatorname{Re}(x)=1 / 2$. This simple dynamics in $x$ is the main reason why the degenerate Newton map is much easier to understand than those considered in [37]: here all points in $\mathbb{C}^{2}$ with $\operatorname{Re}(x)<1 / 2$ are super-attracted to the line $x=0$ and all points with $\operatorname{Re}(x)>1 / 2$ are super-attracted to the line $x=1$. The vertical line at $x=m$ is mapped to the line at $x=m^{2} /(2 m-1)$ by the second coordinate of (4), which is in fact a rational map of degree 2 , except at those values of $m$ where the numerator and the denominator in the second coordinate of (4) have a common factor. This occurs exactly when $x=1 / B, x=1 /(2-B)$, and $x=1 / 2$. The first two correspond to the points of indeterminacy $p$ and $q$.

There are three major invariant sets: $X_{l}:=\{(x, y): \operatorname{Re}(x)<1 / 2$ and $x \neq \infty\}, X_{1 / 2}:=$ $\{(x, y): \operatorname{Re}(x)=1 / 2$ or $x=\infty\}$ and $X_{r}:=\{(x, y) \in: \operatorname{Re}(x)>1 / 2$ and $x \neq \infty\}$. Figure 2 shows the case when both points of indeterminacy $p$ and $q$ are in $X_{l}$. The coordinates of $p$ and $q$ are $p=\left(\frac{1}{B}, 0\right)$ and $q=\left(\frac{1}{2-B}, \frac{1-B}{2-B}\right)$. It is easy to check that $p$ and $q$ either are both in $X_{l}$, both in the separator $X_{1 / 2}$, or both in $X_{r}$. Let $\Omega=\{B \in \mathbb{C}:|1-B|>1\}$ so that if $B \in \Omega$ then both $p$ and $q$ are in $X_{l}$. Using the transformation Rules 1.1, one sees that systems with this restriction still represent every conjugacy class except for those corresponding to both $p, q \in X_{1 / 2}$.

Let $S_{0}$ and $S_{1}$ be the invariant circles in the fixed lines $x=0$ and $x=1$, respectively. Because the lines $x=0$ and $x=1$ are super-attracting in the $x$-direction, $S_{0}$ and $S_{1}$ are super-attracting in the $x$-direction, as well. In Section 3 we will show that these circles have local superstable manifolds $W_{0}^{l o c}$ and $W_{1}^{l o c}$. Pulling $W_{0}^{l o c}$ and $W_{1}^{l o c}$ back under the Newton map we generate superstable spaces $W_{0}$ and $W_{1}$ that form the boundary between the basin $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ and between the basin $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$, respectively. Figure 3 shows an illustration of these separatrices.


Figure 3: Superstable separatrices in the degenerate case, $A=0$.

Proposition 1.6. (Axis of symmetry) Let $\tau$ denote the vertical reflection about the line $B x+$ $2 y-1=0$, that is: $\tau(x, y)=(x, 1-B x-y)$. Then, $\tau$ is a symmetry of $N$ :

$$
\tau \circ N=N \circ \tau
$$

Furthermore, $N$ maps this axis of symmetry to the line $y=\infty$.
Proof: The map $\tau$ is affine and interchanges $r_{1}$ with $r_{2}$ and $r_{3}$ with $r_{4}$. Let $F\binom{x}{y}=\binom{P(x, y)}{Q(x, y)}$ so that $r_{1}, r_{2}, r_{3}$, and $r_{4}$ are the roots of $F$. By Proposition 1.2, the Newton map $N_{F \circ \tau}$ for finding the
roots of $F \circ \tau$ is the same as $N_{F}$, since they have the same roots. By the transformation rules of the Newton Map under affine coordinate changes, $N_{F \circ \tau}=\tau^{-1} \circ N \circ \tau$. Hence:

$$
\tau \circ N=\tau \circ N_{F \circ \tau}=\tau \circ \tau^{-1} \circ N \circ \tau=N \circ \tau
$$

The axis $B x+2 y-1=0$ is mapped to the line $y=\infty$ due to the factor $B x+2 y-1=0$ in the denominator the second component of $N$.

## 2 Computer exploration of $N$

In this section we show computer images of the basins of attraction for the four common zeros of $P$ and $Q$ for the parameters $B=0.7857+1.1161 i$, and $B=-0.7902+1.7232 i$. All of the computer images displayed in this paper were made using the wonderful program FractalAsm [44].

The separatrices $W_{0}$ and $W_{1}$ are clearly visible in these images, forming the smooth boundary between $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ and between $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$, respectively. The boundary between $W\left(r_{1}\right) \cup W\left(r_{2}\right)$ and $W\left(r_{3}\right) \cup W\left(r_{4}\right)$, when visible, corresponds to points $(x, y)$ with $\operatorname{Re}(x)=1 / 2$.
Case 1: $B=0.7857+1.1161 i$
The first kind of slice that we consider is given by the critical value parabola $C$, which is parameterized by a single complex variable, the offset from the axis of $C$. Figure 4 shows part of this slice on the left while the image on the right shows a zoomed in view corresponding to the region enclosed in the small rectangle from the image on the left. The center of the symmetry $\tau$ is the center of the image on the left of Figure 4, but is outside of the image on the right.


Figure 4: The critical value parabola $C$ for $B=0.7857+1.1161 i$. The boundary between $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ is $W_{0} \cap C$ and the boundary between the $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$ is $W_{1} \cap C$. The image on the right is a zoomed in view of the boxed region in the image on the left.

Figure 5 shows the vertical line $x=a$, where $a$ is labeled in Figure 4, as well as the vertical lines through three inverse images of $a$. We have placed the center of the symmetry $\tau$ at the center of these images so that reflection across this point perfectly interchanges the basins.

Notice how the first inverse image of $x=a$ is divided into two regions that are in $W\left(r_{1}\right)$ and two regions in $W\left(r_{2}\right)$. This is because we chose $a$ on the superstable separatrix $W_{0}$. The lines at


Figure 5: Vertical line through point $a$ from Figure 4 and three inverse images of this line. The boundary between $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ is the intersection of $W_{0}$ with these vertical lines. Notice that there are many closed loops in $W_{0}$ within these vertical lines. The center of the symmetry $\tau$ is at the center of these images.
second and third inverse images of $x=a$ are divided into three regions in $W\left(r_{1}\right)$ and in $W\left(r_{2}\right)$ and five regions in $W\left(r_{1}\right)$ and in $W\left(r_{2}\right)$, respectively.
Case 2: $B=-0.7902+1.7232 i$
Figure 6 shows the intersections of the basins of attraction for $W\left(r_{1}\right), W\left(r_{2}\right), W\left(r_{3}\right)$, and $W\left(r_{4}\right)$ with the critical value parabola $C$. Notice that there are clearly intersections of the superstable separatrix $W_{0}$ with $C$ forming the visible boundary between $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$. However, we see no boundaries between $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$, indicating that $W_{1}$ might not intersect $C$. All of the further zoom-ins that we have done show no evidence of intersections between $W_{1}$ and $C$. We cannot prove that there are values of $B$ for which $W_{1} \cap C=\emptyset$, however it seems likely, based on computer experiments.

As for the previous value of $B$, the vertical lines above points of intersection of $W_{0}$ with $C$ and the vertical lines mapped to them by $N$ contain many interestingly loops that are in $W_{0}$.

## 3 Superstable separatrices $W_{0}$ and $W_{1}$.

The invariant circle $S_{0}$ is the set of points in the line $x=0$ equidistant from $r_{1}$ and $r_{2}$ and the invariant circle $S_{1}$ is the set of points in the line $x=1$ equidistant from $r_{3}$ and $r_{4}$.

Proposition 3.1. The invariant circles $S_{0}$ and $S_{1}$ have multiplier 0 in the $x$-direction and they have multiplier 2 within the vertical line in the direction normal to the circle.

Proof: The vertical lines $x=0$ and $x=1$ are superattracting in the $x$-direction, hence the circles $S_{0}$ and $S_{1}$ are as well. Within these vertical lines, $N$ is the Newton's method for the quadratic polynomial with roots $r_{1}$ and $r_{2}$ (or $r_{3}$ and $r_{4}$ ), so the invariant circle is repelling in this line with multiplier 2 .


Figure 6: Critical value parabola $C$ for $B=-0.7902+1.7232 i$. The boundary between the $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ is $W_{0} \cap C$. We see no boundaries between $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$, indicating that $W_{1}$ might not intersect $C$.

Proposition 3.2. The invariant circles $S_{0}$ and $S_{1}$ have local superstable manifolds $W_{0}^{\text {loc }}$ and $W_{1}^{\text {loc. }}$. More specifically, there are neighborhoods $U_{0}, U_{1} \subset \mathbb{C}$ of $x=0$ and $x=1$ and subsets $W_{0}^{\text {loc }} \subset X_{l}$, $W_{1}^{l o c} \subset X_{r}$ so that:

- $N\left(W_{0}^{l o c}\right) \subset W_{0}^{l o c}$ and $N\left(W_{1}^{l o c}\right) \subset W_{1}^{l o c}$
- $W_{0}^{\text {loc }}$ is the image of some $\Phi_{0}: U_{0} \times S_{0} \rightarrow X_{l}$ which is analytic in the first coordinate and quasiconformal in the second.
- $W_{1}^{\text {loc }}$ is the image of some $\Phi_{1}: U_{1} \times S_{1} \rightarrow X_{r}$ which is analytic in the first coordinate and quasiconformal in the second.

We use a technique due to John Hubbard and Sebastien Krief which allows us to use the theory of holomorphic motions and the $\lambda$-Lemma of Mañe, Sad, and Sullivan [41], instead of the more standard graph transformation approach. A somewhat different stable manifold theorem for the invariant circles in the non-degenerate case $(A \neq 0)$ is also proved using this technique in [37]. While points in the manifolds obtained in our proof are genuinely attracted to the circles $S_{0}$ and $S_{1}$, the situation in [37] is much more complicated, with dense sets of points that are not attracted to the invariant circles.
Proof: To simplify computations we will make the change of variables $z(x)=\frac{x}{x-1}$ and $w(y)=\frac{y}{y-1}$ which conjugates the first coordinate of $N$ to $z \mapsto z^{2}$ and places the invariant circle $S_{0}$ at $\{z=$ $0,|w|=1\}$. In the new coordinates $(z, w)$, the Newton map becomes:

$$
\begin{equation*}
N\binom{z}{w}=\binom{z^{2}}{\frac{w^{2}+\left(B w-B w^{2}\right) z-w^{2} z^{2}}{1+(B-B w) z+\left(B w^{2}+B-1-2 B w\right) z^{2}}} \tag{6}
\end{equation*}
$$

and the critical value locus of $N$ in these coordinates is the image of $C$ under the change of variables, which we denote by $C^{\prime}$.

$$
\Delta_{\epsilon, \delta}=\left\{(z, w) \in X_{l}:|z|<\epsilon \text { and } 1-\delta<|w|<1+\delta\right\}
$$

so that $\Delta_{\epsilon, \delta}$ is an open neighborhood of $S_{0}$. The boundary of $\Delta_{\epsilon, \delta}$ consists of the vertical boundary $\partial^{V} \Delta_{\epsilon, \delta}=\{|z|=\epsilon\}$ and the horizontal boundary $\partial^{H} \Delta_{\epsilon, \delta}=\{|w|=1 \pm \delta\}$.

We must choose $\epsilon$ and $\delta$ so that:

1. $\Delta_{\epsilon, \delta}$ is disjoint from the critical value locus $C^{\prime}$, and
2. $N\left(\partial^{H} \Delta_{\epsilon, \delta}\right)$ is entirely outside of $\Delta_{\epsilon, \delta}$ and $N\left(\partial^{V} \Delta_{\epsilon, \delta}\right)$ is entirely inside of $|z|<\epsilon$.


Figure 7: Here $\Delta_{\epsilon, \delta}$ is shown in light grey and $N\left(\Delta_{\epsilon, \delta}\right)$ in dark grey.
The critical value locus $C^{\prime}$ intersects the vertical line $z=0$ transversely at $w=0$ and $w=\infty$, so we can choose $\epsilon$ sufficiently small so that $C^{\prime}$ intersects $\mathbb{D}_{\epsilon} \times \mathbb{P}$ outside of $\Delta_{\epsilon, \frac{1}{2}}$. Now, we reduce $\epsilon$ and $\delta$ so that the second condition holds. Because the first coordinate of $N$ is $z \mapsto z^{2}, N\left(\partial^{V} \Delta_{\epsilon, \delta}\right)$ is automatically inside of $|z|<\epsilon$. In the line $z=0, N(z, w)=w^{2}$, so by continuity we can choose $\epsilon$ and $\delta$ small enough that $N\left(\partial^{H} \Delta_{\epsilon, \delta}\right)$ is entirely outside of $\Delta_{\epsilon, \delta}$.

Let $\mathbb{D}_{\epsilon}$ be the open disc $|z|<\epsilon$ in $\mathbb{C}$ for this $\epsilon$. Conditions 1 and 2 on $\epsilon$ and $\delta$ were chosen so that the following lemma is true:

Lemma 3.3. Suppose that $D \subset \Delta_{\epsilon, \delta}$ is a complex disc that is the graph of an analytic function $\eta: \mathbb{D}_{\epsilon} \rightarrow \mathbb{P}$. Then $N^{-1}(D) \cap \Delta_{\epsilon, \delta}$ is the union of two disjoint complex discs, each given as the graph of analytic functions $\zeta_{1}, \zeta_{2}: \mathbb{D}_{\epsilon} \rightarrow \mathbb{P}$.
Proof of Lemma 3.3: The locus $N^{-1}(D) \cap \Delta_{\epsilon, \delta}$ satisfies the equation $N(z, w) \in D$, which is equivalent to $N_{2}(z, w)=\eta\left(z^{2}\right)$, because $D$ is the graph of $\eta$. Since $D \subset \Delta_{\epsilon, \delta}, D$ is disjoint from $C^{\prime}$, so $\frac{\partial}{\partial w} N_{2}(z, w)$ is non-zero in a neighborhood of $N^{-1}(D)$, and we can use the implicit function theorem to solve for $w=\zeta_{1}(z)$ and $w=\zeta_{2}(z)$. There are exactly two branches because $N_{2}(z, w)$ is degree 2 in $w$.

The graphs of $\zeta_{1}$ and $\zeta_{2}$ form the two complex discs $N^{-1}(D) \cap \Delta_{\epsilon, \delta} . \square$ Lemma 3.3.
The line $w=1$ is invariant under $N$ and attracted to the point $(0,1) \in S_{0}$. Let $D_{0}=$ $\{(z, w):|z|<\epsilon, w=1\}$. We will form $W_{0}^{l o c}$ by taking inverse images of $D_{0}$.

Since $D_{0} \subset \Delta_{\epsilon, \delta}$ satisfies the conditions of Lemma 3.3, letting $D_{1}=N^{-1}\left(D_{0}\right) \cap \Delta_{\epsilon, \delta}$ we obtain two complex discs in $\Delta_{\epsilon, \delta}$ each of which is given by the graph of an analytic function $\eta: \mathbb{D}_{\epsilon} \rightarrow \mathbb{P}$ and each of which is mapped within $D_{0}$ by $N$. These discs intersect $S_{0}$ at $w=1$ and -1 .

Because each of the discs in $D_{1}$ satisfies the hypotheses of Lemma 3.3 we can repeat this process, letting $D_{2}=N^{-1}\left(D_{1}\right) \cap \Delta_{\epsilon, \delta}$, which this lemma guarantees is the union of four disjoint discs in
$\Delta_{\epsilon, \delta}$, each of which is the graph of some analytic function $\eta: \mathbb{D}_{\epsilon} \rightarrow \mathbb{P}$. These four discs intersect $S_{0}$ at the fourth roots of 1 . Repeating this process, we obtain $D_{n}$ consisting of $2^{n}$ disjoint complex discs in $\Delta_{\epsilon, \delta}$, each given by the graph of an analytic function intersecting $S_{0}$ at the $2^{n}$-th roots of 1 .

Let $D_{\infty}=\bigcup_{n=0}^{\infty} D_{n}$, which consists of a union of disjoint complex discs through each of the dyadic points $\mathcal{D}$ on $S_{0}$. Each of these discs is the graph of an analytic function from $\mathbb{D}_{\epsilon}$ to $\mathbb{P}$, and $D_{\infty}$ is forward invariant to $S_{0}$ under $N$.

From a different perspective, $D_{\infty}$ prescribes a holomorphic motion:

$$
\phi: \mathbb{D}_{\epsilon} \times \mathcal{D} \rightarrow \mathbb{P}
$$

where $\phi(z, \theta)$ is given by $\eta(z)$ where $\eta: \mathbb{D}_{\epsilon} \rightarrow \mathbb{P}$ is the analytic function whose graph is the disc in $D_{\infty}$ containing $\theta \in \mathcal{D} \subset S_{0}$.

By the $\lambda$-lemma of Mañe-Sad-Sullivan [41], $\phi$ extends continuously to a holomorphic motion on $S_{0}$, the closure of $\mathcal{D}$.

$$
\phi: \mathbb{D}_{\epsilon} \times S_{0} \rightarrow \mathbb{P}
$$

We define $W_{0}^{\text {loc }}$ to be the image of $(z, w) \mapsto(z, \phi(z, w))$. Clearly $N\left(W_{0}^{\text {loc }}\right) \subset W_{0}^{\text {loc }}$ and every point in $W_{0}^{\text {loc }}$ is forward invariant to $S_{0}$.

The construction of of $W_{1}^{l o c}$ is nearly identical and we omit it. $\square$ Proposition 3.2.
Because the local superstable manifolds $W_{0}^{\text {loc }}$ and $W_{1}^{\text {loc }}$ are forward invariant under $N$, we can define global invariant sets $W_{0}$ and $W_{1}$ by:

$$
W_{0}=\bigcup_{n=0}^{\infty} N^{-n}\left(W_{0}^{l o c}\right), \quad W_{1}=\bigcup_{n=0}^{\infty} N^{-n}\left(W_{1}^{l o c}\right) .
$$

One might expect that $W_{0}$ and $W_{1}$ are manifolds, since the Inverse Function Theorem gives that the pull-back $N^{-k}\left(W_{i}^{l o c}\right)\left(\right.$ or $\left.N^{-k}\left(W_{i}^{l o c}\right)\right)$ by $N$ is "locally manifold" at points where $N^{-(k-1)}\left(W_{i}^{\text {loc }}\right)$ is disjoint from or transverse to the critical value locus $C$. However, we do expect that there will be some values of the parameter $B$ for which there is a tangency between $N^{-k}\left(W_{i}^{l o c}\right)$ and $C$. In fact, our computer images show that this must be the case, because we see the topology of $W_{0} \cap C$ and of $W_{1} \cap C$ change as we change $B$. For these parameter values $W_{i}$ will not be a manifolds. To make this distinction, we will call $W_{0}$ and $W_{1}$ separatrices because they separate $W\left(r_{1}\right)$ from $W\left(r_{2}\right)$ and separate $W\left(r_{3}\right)$ from $W\left(r_{4}\right)$.

The following proposition requires that all iterates of $N$ be defined for all points in a neighborhood of $W_{0}$ in $X_{l}$ and in a neighborhood of $W_{1}$ in $X_{r}$. This will require a modification $X_{l}^{\infty}$ of $X_{l}$ that is obtained by blowing-up the points of indeterminacy $p$ and $q$ and all of their inverse images under $N$. We will prove Proposition 3.4, temporarily thinking that we are working in $X_{l}$ and $X_{r}$, and then explain why it is necessary to blow-up the points of indeterminacy. The entire construction of $X_{l}^{\infty}$ is given in the following section.
Proposition 3.4. For every $B \in \Omega$, the separatrices $W_{0}$ and $W_{1}$ are real analytic subspaces of $X_{l}^{\infty}$ and $X_{r}$, each defined as the zero set of a single non-constant real-analytic equation in an neighborhood of $W_{0}$ and in a neighborhood of $W_{1}$, respectively.
The proof is similar of that of Böttcher's Theorem in one variable dynamics.
Proof: We express $N$ in the variables $z=\frac{x}{x-1}$ and $w=\frac{y}{y-1}$ so that $S_{0}$ is given by $\{z=0,|w|=1\}$. We will show that

$$
\phi(z, w)=\lim _{n \rightarrow \infty}\left(N_{2}^{n}(z, w)\right)^{1 / 2^{n}}
$$

converges to a non-constant analytic function on a neighborhood of $W_{0}$. Then, for every $(z, w) \in W_{0}$, $\left|N_{2}^{n}(z, w)\right|$ converges to 1 because $S_{0}=\{|w|=1\}$, hence $\omega(z, w):=\log |\phi(z, w)|=\log \left|\left(N_{2}^{n}(z, w)\right)^{1 / 2^{n}}\right|$ converges to 0 on $W_{0}$ and to non-zero values away from $W_{0}$.

If $\phi$ converges, then $\phi$ and $\omega$ transform nicely under the involution $\tau$. For $|z|$ small and $|w|$ close to $1, \tau$ is close to $(z, w) \mapsto(z, 1 / w)$. Using this approximation, we have $\phi(\tau(z, w))=$ $\lim _{n \rightarrow \infty}\left(N_{2}^{n}(\tau(z, w))\right)^{1 / 2^{n}}=\lim \left(\tau\left(N^{n}(z, w)\right)_{2}\right)^{1 / 2^{n}} \approx \lim \left(1 / N_{2}^{n}(z, w)\right)^{1 / 2^{n}}=1 / \phi(z, w)$. Consequently, $\omega(\tau(z, w))=-\omega(z, w)$.

We can write $\phi(z, w):=\lim _{n \rightarrow \infty}\left(N_{2}^{n}(z, w)\right)^{1 / 2^{n}}$ as a telescoping product:

$$
\begin{equation*}
\phi(z, w)=N_{2}(z, w)^{1 / 2} \cdot \frac{N_{2}^{2}(z, w)^{1 / 4}}{N_{2}(z, w)^{1 / 2}} \cdot \frac{N_{2}^{3}(z, w)^{1 / 8}}{N_{2}^{2}(z, w)^{1 / 4}} \cdots \tag{7}
\end{equation*}
$$

We now check that we can restrict the neighborhood of $W_{0}$ where $\phi$ is defined so that we can use the binomial formula

$$
\begin{equation*}
(1+u)^{\alpha}=\sum_{n=0}^{\infty} \frac{\alpha(\alpha-1) \cdots(\alpha-n+1)}{n!} u^{n}, \text { when }|u|<1 \tag{8}
\end{equation*}
$$

to define the $\frac{1}{2^{n}}$-th root in the $n$-th term of this product. We do this first for points in a neighborhood of $W_{0}$ in $\overline{W\left(r_{1}\right)}$ and a similar proof shows that the same works in a neighborhood of $W_{0}$ in $\overline{W\left(r_{2}\right)}$.

In the coordinates $(z, w)$ the denominator of $N_{2}$ is of the form $1+r$ with $r=(B-B w) z+$ $\left(B w^{2}+B-1-2 B w\right) z^{2}$, so for $\left|(B-B w) z+\left(B w^{2}+B-1-2 B w\right) z^{2}\right|<1$ the second coordinate of $N$ can me written as:

$$
\begin{equation*}
N_{2}(z, w)=w^{2}\left(1-B z-z^{2}\right)+w z g(z, w) \tag{9}
\end{equation*}
$$

with $g(z, w)$ analytic. We can write $N_{2}$ in form (9) in a neighborhood of $z=0$ (hence a neighborhood of $\left.S_{0}\right)$ since $\left|(B-B w) z+\left(B w^{2}+B-1-2 B w\right) z^{2}\right|$ vanishes when $z=0$. From this point on, we restrict our attention to this neighborhood of $z=0$.
The general term $\frac{N_{2}^{n+1}(z, w)^{1 / 2^{n+1}}}{N_{2}^{n}(z, w)^{1 / 2^{n}}}$ in the product (7) is of the form

$$
\begin{gathered}
\left(\frac{\left(N_{2}^{n}(z, w)\right)^{2}\left(1-B N_{1}^{n}(z, w)-\left(N_{1}^{n}(z, w)\right)^{2}\right)+N_{2}^{n}(z, w) N_{1}^{n}(z, w) \cdot g\left(N_{1}^{n}(z, w), N_{2}^{n}(z, w)\right)}{\left(N_{2}^{n}(z, w)\right)^{2}}\right)^{1 / 2^{n+1}} \\
=\left(1-B N_{1}^{n}(z, w)-\left(N_{1}^{n}(z, w)\right)^{2}+\frac{N_{1}^{n}(z, w)}{N_{2}^{n}(z, w)} \cdot g\left(N_{1}^{n}(z, w), N_{2}^{n}(z, w)\right)\right)^{1 / 2^{n+1}}
\end{gathered}
$$

We need to check that we can restrict, if necessary, the neighborhood of definition for $\phi(z, w)$ so that

$$
\begin{array}{r}
\left|-B N_{1}^{n}(z, w)-\left(N_{1}^{n}(z, w)\right)^{2}+\frac{N_{1}^{n}(z, w)}{N_{2}^{n}(z, w)} \cdot g\left(N_{1}^{n}(z, w), N_{2}^{n}(z, w)\right)\right| \\
\leq\left|B N_{1}^{n}(z, w)+\left(N_{1}^{n}(z, w)\right)^{2}\right|+\left|\frac{N_{1}^{n}(z, w)}{N_{2}^{n}(z, w)} \cdot g\left(N_{1}^{n}(z, w), N_{2}^{n}(z, w)\right)\right| \leq \frac{1}{2} \tag{11}
\end{array}
$$

The first term is not a problem because $N_{1}^{n}(z, w)=z^{2 n}$ and we are restricting to $|z|$ small. Since we are in $\overline{W\left(r_{1}\right)}$ the only difficulty can can occur if $N_{2}^{n}(z, w)$ goes to 0 fast enough to make (10) large. Detailed analysis of the behavior near $r_{1}$ resolves this concern:

In [37], the authors perform blow-ups at each of the four roots, and observe that the Newton map $N$ induces rational functions of degree 2 on each of the exceptional divisors $E_{r_{1}}, E_{r_{2}}, E_{r_{3}}$, and $E_{r_{4}}$. Let's compute the rational function $s: E_{r_{1}} \rightarrow E_{r_{1}}$. In the coordinate chart $m=\frac{z}{w}$, the extension to $E_{r_{1}}$ is obtained by:

$$
s(m)=\lim _{w \rightarrow 0} \frac{m^{2} w^{2}\left(1+(B-B w) m w+\left(B w^{2}+B-1-2 B w\right) m^{2} w^{2}\right)}{w^{2}+\left(B w-B w^{2}\right) m w-w^{2} m^{2} w^{2}}=\frac{m^{2}}{1+B m}
$$

since $w=0$ on $E_{r_{1}}$.
The rational function $s(m)$ has $m=0$ as a superattracting fixed point, so there is a neighborhood of $m=0 \in E_{r_{1}}$ within $\overline{W\left(r_{1}\right)}$ so that for any point $(z, w)$ in this neighborhood, $\lim _{n \rightarrow \infty}\left|\frac{N_{1}^{n}(z, w)}{N_{2}^{n}(z, w)}\right|=$ 0 . Pulling back this neighborhood under $N$ we find a neighborhood $V \subset \overline{W\left(r_{1}\right)}$ of the line $z=0$ in which this limit holds.

Now we consider the case when $(z, w) \in \overline{W\left(r_{2}\right)}$. The concern is that $|w|$ may grow too fast for us to find a neighborhood of $S_{0}$ inequality (10) is true. Instead of analyzing the asymptotics of $g$, we can re-write $N$ in the new coordinates $(z, s)$ with $s=1 / w$ and a nearly identical construction to that of $V$ gives the appropriate neighborhood $V^{\prime}$.

Restricting the points $(z, w) \in V \cup V^{\prime}$, the $\frac{1}{2^{n+1}}$-th root in the product (7) is well-defined. We check that the product converges on the neighborhood $\Lambda$ of $S_{0}$. It is sufficient to show that the corresponding series of logarithms converges. The general term in this series is:

$$
\log \left|\left(1-B N_{1}^{n}(z, w)-\left(N_{1}^{n}(z, w)\right)^{2}+\frac{N_{1}^{n}(z, w)}{N_{2}^{n}(z, w)} \cdot g\left(N_{1}^{n}(z, w), N_{2}^{n}(z, w)\right)\right)^{1 / 2^{n+1}}\right| \leq \frac{\log 2}{2^{n+1}}
$$

using Equation 10 and the triangle inequality so that

$$
\left|1-B N_{1}^{n}(z, w)-\left(N_{1}^{n}(z, w)\right)^{2}+\frac{N_{1}^{n}(z, w)}{N_{2}^{n}(z, w)} \cdot g\left(N_{1}^{n}(z, w), N_{2}^{n}(z, w)\right)\right|<2
$$

This sequence of logarithms converges because it is dominated by a geometric series, and hence for the product (7) converges to the analytic function on $\phi(z, w)$ on $\Lambda$. This way $\omega(z, w)=\log |\phi(z, w)|$ is a real analytic function on $\Lambda$, and by the invariance properties of $\phi$ on $\omega(z, w)$ is an analytic function on a neighborhood of $W_{0}$.

The proof that $W_{1}$ is the zero locus of a non-constant analytic function is very similar.
It is important to notice that in this proof we assumed that all iterates of $N$ are defined at every point in $X_{l}$, forgetting temporarily the points of indeterminacy $p$ and $q$ (and all of their inverse images in $X_{l}$.) This is a real problem because $W_{0}$ naturally goes through all of the points of indeterminacy: Under a high enough iterate of $N$, the line $x=\frac{1}{B(2-B)}$ is mapped by a ramified covering to a vertical line arbitrarily close to the line $x=0$. Since these lines intersect $W_{0}^{\text {loc }}$ in a topological circle, the line $x=\frac{1}{B(2-B)}$ intersects $W_{0}$ in a (possibly more complicated) curve. We will see that the resolution of the indeterminacy at $p$ and $q$ replaces $p$ and $q$ with exceptional divisors $E_{p}$ and $E_{q}$ that are mapped to $x=\frac{1}{B(2-B)}$ by isomorphisms. So, to make this proof correct, we will have to blow-up at $p$ and $q$, and, in fact, at all inverse images of $p$ and $q$.

An alternative approach would be to study $W_{0}$ on $X_{l}-\cup_{n=0}^{\infty} N^{-n}(\{p, q\})$, where we have already proven it is a real-analytic variety. However, we want to consider the topology of $W_{0}$, without all of these points removed, so we prefer to do the sequence of blow-ups.

## 4 Resolution of points of indeterminacy

By restricting to parameters $B \in \Omega$, the points of indeterminacy $p$ and $q$ are in $X_{l}$ and there are no points of indeterminacy in $X_{r}$. In this section we will describe how to resolve the indeterminacy in $N$ at $p$ and $q$ and in higher iterates of $N$ at all of the inverse images of $p$ and $q$ in $X_{l}$, obtaining a new space $X_{l}^{\infty}$ on which all iterates of $N$ are defined at every point.

Writing $W_{0}$ as a real-analytic variety is not the only motivation for the construction of $X_{l}^{\infty}$. We plan to study the detailed topology of the basins of attraction and of the separatrices $W_{0}$ and $W_{1}$. It is difficult to decide what is a reasonable alternative to the statement of Theorem 0.1 without blowing up points.

### 4.1 Construction of $X_{l}^{\infty}$ and $N_{\infty}: X_{l}^{\infty} \rightarrow X_{l}^{\infty}$.

Most of the material in this section and in the following section closely follow the works of Hubbard and Papadopol [37] and Hubbard, Papadopol, and Veselov [33].

Substitution of the points $p$ and $q$ into $C(x, y)$ yields $\frac{1}{4}(B-1)$ and $\frac{B^{2}-7 B+2}{4 B-8}$, so for values of $B$ at which these expressions are non-zero, neither $p$ nor $q$ is a critical value.

Let $S \subset \Omega$ be the subset of parameter space for which no inverse image of the point of indeterminacy $p$ or of point of indeterminacy $q$ is in the critical value locus $C$. We first describe the construction of $X_{l}^{\infty}$ for parameter values $B \in S$, and then explain the necessary modifications for special circumstance when $B \notin S$.

Theorem 4.1. The set $S$ is generic in the sense of Baire's Theorem, i.e. uncountable and dense in $\Omega$.

Because of its computational nature, the proof of Theorem 4.1 is in Appendix B.
Construction of $X_{l}^{\infty}$ when $B \in S$ :
Proposition 4.2. Let $X_{l}^{0}$ be the space $X_{l}$ blown up at the points $p$ and $q$ and let $\pi_{0}: X_{l}^{0} \rightarrow X_{l}$ be the corresponding projection.

- The mapping $N$ extends analytically to a mapping $N_{0}: X_{l}^{0} \rightarrow X_{l}$.
- $N_{0}$ maps the exceptional divisors $E_{p}$ and $E_{q}$ to the line $x=\frac{1}{B(2-B)}$ by isomorphisms.

Proof: The definition of a blow-up at a point is available in Appendix A. Further details about blow-ups are available in books on Algebraic Geometry, including [28], and, in the context of complex dynamics, in the papers [33, 37].

We will work in the chart $(x, m) \mapsto\left(x, m\left(x-\frac{1}{B}\right), m\right) \in X_{l} \times \mathbb{P}^{1}$. We write $N(x, y)=\left(N_{1}(x, y), N_{2}(x, y)\right)$ so that in the coordinates $(x, m)$ we have $N_{1}(x, m)=\frac{1}{B(2-B)}$ and

$$
N_{2}(x, m)=\frac{\frac{m}{B}\left(B x^{2}+2 x m\left(x-\frac{1}{B}\right)-B x-m\left(x-\frac{1}{B}\right)\right)}{(2 x-1)\left(\frac{2 m}{B}+1\right)}
$$

When restricted to the exceptional divisor $E_{p}\left(\right.$ set $x=\frac{1}{B}$ ) the mapping becomes $m \mapsto \frac{m(1-B)}{(2-B)(2 m+B)}$. If instead we had been working in the chart $\left(y, m^{\prime}\right) \mapsto\left(m^{\prime} y+\frac{1}{B}, y, m^{\prime}\right)$, we would have obtained $m^{\prime} \mapsto \frac{(1-B)}{(2-B)\left(2+m^{\prime} B\right)}$. This is consistent with the extension in terms of $m$ since one is obtained from the other by the change of variables $m=\frac{1}{m^{\prime}}$. Both of the expressions for $N$ restricted to $E_{p}$ are linear-fractional transformations, hence $N$ maps $E_{p}$ to the line $y=\frac{1}{B(2-B)}$ by an isomorphism.

The exceptional divisor $E_{q}$ can be treated similarly.
We will denote the vertical line $x=\frac{1}{B(2-B)}$ by $V$, since we use this line so frequently. This is the vertical line that is tangent to $C$ at its "vertex".

By construction, the extension $N_{0}: X_{l}^{0} \rightarrow X_{l}$ has no points of indeterminacy. However, since we need to iterate we must consider $N_{0}$ as a rational map $N_{0}: X_{l}^{0} \rightarrow X_{l}^{0}$. Each of the inverse images of $p$ and $q$ become points of indeterminacy of $N_{0}$ because we have blown up at $p$ and $q$. Because $B \in S$, neither $p$ nor $q$ are critical values and each has four inverse images under $N_{0}$. We can blow-up at each of these eight points obtaining the space $X_{l}^{1}$ and the projection $\pi_{1}: X_{l}^{1} \rightarrow X_{l}^{0}$. One can then extend $N_{0}$ to the exceptional divisors, obtaining $N_{1}: X_{l}^{1} \rightarrow X_{l}^{0}$.

To make iterates $N^{\circ k}$ of $N$ well-defined for all $k$ we must repeat this process for the $k$-th inverse images, obtaining successive blow-ups $\pi_{k}: X_{l}^{k} \rightarrow X_{l}^{k-1}$ for every $k$. The following proposition describes the extension of $N$ to these spaces:
Proposition 4.3. Denote by $X_{l}^{k}$ the space $X_{l}^{k-1}$ blown up at each of these $2 \cdot 4^{k} k$-th inverse images of $p$ and $q$.

- The mapping $N_{k-1}$ extends analytically to a mapping $N_{k}: X_{l}^{k} \rightarrow X_{l}^{k-1}$.
- Suppose that $z$ is one of the $k$-th inverse images of $p$ or $q$ and denote the exceptional divisor over $z$ by $E_{z}$. Then, $N_{k}$ maps $E_{z}$ to $E_{N(z)}$ by isomorphism.

Proof: As in Proposition 4.2, denote the first and second components of $N$ by $N_{1}(x, y)$ and $N_{2}(x, y)$. Then, in the coordinates $(x, m) \mapsto(x, m x, m)$ in a neighborhood of $E_{z}$ the mapping is given by: $m \mapsto \frac{\partial_{x} N_{1}+\partial_{y} N_{1} m}{\partial_{x} N_{2}+\partial_{y} N_{2} m}$. By the assumption that $B \in S, D N$ is non-singular at $z$ and this gives an isomorphism from $E_{z}$ to $E_{N(z)}$.

Hence, by repeated blow-ups we obtain a sequence of spaces and projections:

$$
\begin{equation*}
X_{l} \stackrel{\pi_{0}}{\leftarrow} X_{l}^{0} \stackrel{\pi_{1}}{\leftrightarrows} X_{l}^{1} \stackrel{\pi_{2}}{\leftrightarrows} X_{l}^{2} \stackrel{\pi_{3}}{\leftrightarrows} X_{l}^{3} \stackrel{\pi_{4}}{\leftrightarrows} X_{l}^{4} \stackrel{\pi_{5}}{\leftrightarrows} X_{l}^{5} \stackrel{\pi_{6}}{\leftrightarrows} \cdots \tag{12}
\end{equation*}
$$

The extensions of the Newton map $N$ to these spaces that we calculated in Propositions 4.2 and 4.3 give another sequence of spaces and mappings:

$$
\begin{equation*}
X_{l} \stackrel{N_{0}}{\longleftarrow} X_{l}^{0} \stackrel{N_{1}}{\longleftarrow} X_{l}^{1} \stackrel{N_{2}}{\longleftarrow} X_{l}^{2} \stackrel{N_{3}}{\longleftarrow} X_{l}^{3} \stackrel{N_{4}}{\longleftarrow} X_{l}^{4} \stackrel{N_{5}}{\longleftarrow} X_{l}^{5} \stackrel{N_{6}}{\longleftarrow} \cdots \tag{13}
\end{equation*}
$$

However, we do not have a single space $X_{l}^{\infty}$, nor a single mapping $N_{\infty}$ from this space to itself. There is a standard procedure using Inverse Limits to create such a space and mapping from a sequence of spaces (12) and the sequence of mappings like (13). That is, we will let $X_{l}^{\infty}$ be the inverse limit of the blown up spaces and projections in Sequence 12 and then use the sequence of extensions of the Newton maps 13 to define a mapping $N_{\infty}: X_{l}^{\infty} \rightarrow X_{l}^{\infty}$ which naturally corresponds to an extension of $N$.

Definition 4.4. An Inverse system, denoted $\left(M_{i}, \sigma_{i}\right)$, is a family of objects $M_{i}$ in a category $C$ indexed by the natural numbers and for every i a morphism $\sigma_{i}: M_{i} \rightarrow M_{i-1}$.
The Inverse Limit of an inverse system $\left(M_{i}, \sigma_{i}\right)$, denoted by $\underset{\rightleftarrows}{\lim }\left(M_{i}, \sigma_{i}\right)$, is the object $X$ in $C$ together with morphisms $\alpha_{i}: X \rightarrow M_{i}$ satisfying $\alpha_{i-1}=\sigma_{i} \circ \overleftarrow{\alpha_{i}}$ for each $i$ that is determined uniquely by the following universal property:

For any other pair $Y, \beta_{i}: Y \rightarrow M_{i}$ such that $\beta_{i-1}=\sigma_{i} \circ \beta_{i}$, we have a unique morphism $u: Y \rightarrow X$ so that for each $i$ we have $\beta_{i}=\alpha_{i} \circ u$.

For our uses, the category will always be analytic spaces and the morphisms holomorphic maps. While not needed here, inverse systems and inverse limits can be defined more generally, for objects $M_{i}$ indexed by a filtering partially ordered set $I$. The following proposition gives a construction of $\underset{\rightleftarrows}{\lim }\left(M_{i}, \sigma_{i}\right)$ as a subset of the product space $\Pi_{i} M_{i}$.

Proposition 4.5. Given an inverse system $\left(M_{i}, \sigma_{i}\right)$ indexed by $\mathbb{N}$ (i.e. $\sigma_{i}: M_{i} \rightarrow M_{i-1}$ ), we can construct the inverse limit as follows:

$$
\underset{\rightleftarrows}{\lim }\left(M_{i}, \sigma_{i}\right)=\left\{\left(m_{0}, m_{1}, m_{2}, m_{3}, \cdots\right) \mid m_{i} \in M_{i} \text { and } \sigma_{i}\left(m_{i}\right)=m_{i-1}\right\}
$$

We define $X_{l}^{\infty}=\lim _{\rightleftarrows}\left(X_{l}^{k}, \pi_{k}\right)$. Using Proposition 4.5 we can state more concretely that

$$
X_{l}^{\infty}=\left\{\left(x_{0}, x_{1}, x_{2}, x_{3}, \cdots\right) \mid x_{i} \in X_{l}^{i} \text { and } \pi_{i}\left(x_{i}\right)=x_{i-1}\right\} .
$$

In this description of $X_{l}^{\infty}$, the mappings $N_{k}: X_{l}^{k} \rightarrow X_{l}^{k-1}$ induce a mapping $N_{\infty}: X_{l}^{\infty} \rightarrow X_{l}^{\infty}$ given by $N_{\infty}\left(\left(x_{0}, x_{1}, x_{2}, x_{3}, \cdots\right)\right)=\left(N_{1}\left(x_{1}\right), N_{2}\left(x_{2}\right), N_{3}\left(x_{3}\right), \cdots\right)$.

Construction of $X_{l}^{\infty}$ when $B \notin S$ :
For parameter values $B \notin S$, the blow-ups done at $p$ and $q$ in Proposition 4.2 are exactly the same, since $N$ extends to these blow-ups for any value of $B$. (It is worth noticing that there is actually a


Figure 8: Blowing up a point on an exceptional divisor.
critical point of $N$ on both $E_{p}$ and on $E_{q}$ since each is mapped to the line $V$, which contains a point of $C$.) Special care needs to be taken when a $k$-th inverse image of $p$ and of $q$ is a critical point of $N$. We describe the process here, although we leave some of the details for the appendix.

The goal is to produce a space $X_{l}^{k}$ and a projection $\pi_{k}: X_{l}^{k} \rightarrow X_{l}^{k-1}$ in such a way that $N$ extends to a map (without singularities) $N_{k}: X_{l}^{k} \rightarrow X_{l}^{k-1}$. If we can create the spaces $X_{l}^{k}$ and extensions $N_{k}$ at every "level" $k$, we can use the same inverse sequence process to make $X_{l}^{\infty}$ and $N_{\infty}: X_{l}^{\infty} \rightarrow X_{l}^{\infty}$.

Suppose for the moment that $z$ is a $k$-th inverse image of $p$ and that none of the $n$-th inverse images of $p$ for $n<k$ were in the critical locus $N^{-1}(C)$. In this case, there is a single exceptional divisor in $X_{l}^{k-1}$ above $N(z)$. Because $z$ is critical, the extension of $N$ to $E_{z}$ will map all of $E_{z}$ (except for one point) to a single point in $E_{N(z)}$. However, at the slope $m_{k e r} \in E_{z}$ corresponding to the kernel of $D N$, the extension to $E_{z}$ has another point of indeterminacy! Consequently, one has to blow-up this point on $E_{z}$, obtaining a second exceptional divisor $E_{z}^{\prime}$ above $m_{k e r}$. Figure 8 shows this situation. An easy check using Taylor series for $N$ shows that $N$ extends to $E_{z}^{\prime}$ giving isomorphism from $E_{z}^{\prime}$ to $E_{N(z)}$.

These two blow-ups above $z$ are sufficient to extend $N$.
The two exceptional divisors above $z$ result in a further complication at every point $w$ that is mapped to $z$. Suppose that we have blown up at $w$. The extension of $N$ to $E_{w}$ has a point of indeterminacy at the point that is mapped to $m_{k e r} \in E_{z}$. Because of this, one has to blow-up a second time above $w$ to resolve this point of indeterminacy. In fact, at every repeated inverse image of $z$ one will have to blow-up at least twice to resolve $N$.

There are further problems if an inverse image of $z$ is again critical. At such a point, one will have to do even more blow-ups to resolve $N$. A detailed description of this process becomes rather tedious, and we will stop here.

### 4.2 The mappings from $E_{z}$ to $V$

We saw in the previous section that $N$ maps each exceptional divisor that was newly created in $X_{l}^{k}$ to one of the exceptional divisors newly created in $X_{l}^{k-1}$ by either an isomorphism, or a constant map. Since $N$ maps each $E_{p}$ and $E_{q}$ isomorphically to the line $V$, the composition $N^{\circ k+1}$ maps each of the newly created exceptional divisors $E_{z}$ in $X_{l}^{k}$ to $V$ either by an isomorphism, or a constant map. In summary:

Proposition 4.6. Let $E_{z}$ be one of the exceptional divisors newly created in $X_{l}^{k}$ and let $V$ be the line $x=1 /(B(2-B))$. Then $N^{0 k+1}$ maps $E_{z}$ to $V$ by an isomorphism, or a constant map.

### 4.3 Homology of $X_{r}$ and of $X_{l}^{\infty}$

Our goal in this paper is to relate the homology of the basins of attraction for the four roots of $F$ to the homology of the spaces $X_{r}$ and $X_{l}^{\infty}$ and to the homology of the separatrices $W_{0}$ and $W_{1}$. In this section, we will compute the homology of $X_{r}$ and $X_{l}^{\infty}$.

Given a set $T$, we will denote by $\mathbb{Z}^{(T)}$ the submodule of the product $\mathbb{Z}^{T}$ for which each element has at most finitely many non-zero components.

We often find it necessary to encode information about the generators homology within our notation. For example, the module $\mathbb{Z}^{\{[K]\}}$ means the module $\mathbb{Z}$ that is generated by the fundamental class of $[K]$.
Proposition 4.7. We have: $H_{0}\left(X_{r}\right)=\mathbb{Z}, H_{2}\left(X_{r}\right)=\mathbb{Z}^{\left\{\left[\mathbb{P}^{1}\right]\right\}}$, and $H_{i}\left(X_{r}\right)=0$, for $i \neq 0$ or 2 .
Unfortunately homology does not behave nicely under inverse limits. So, instead of directly using the fact that $X_{l}^{\infty}$ is an inverse limit to compute its homology, we will write $X_{l}^{\infty}$ is a union of open subsets $U_{0} \subset U_{1} \subset U_{2} \subset \cdots$ in such a way that $H_{2}\left(U_{i}\right)=\mathbb{Z}^{\left(L_{i} \cup\{[V]\}\right)}$, where $L_{i}$ is the set of fundamental classes of exceptional divisors contained in $U_{i}$ and $[V]$ is the fundamental class of the vertical line $V$ given by $x=\frac{1}{B(2-B)}$.

Recall that the projection $\pi: X_{l}^{\infty} \rightarrow X_{l}$ is continuous, we will create an exhaustion of $X_{l}^{\infty}$ by open sets $U_{0} \subset U_{1} \subset U_{2} \subset \cdots$ that are inverse images of open subsets forming an exhaustion of $X_{l}$. Let $V_{k}=X_{l}-\bigcup_{n=k}^{\infty}\left\{N^{-n}(p), N^{-n}(q)\right\}$. Clearly $V_{k}$ is an open subset of $X_{l}$, so we will let $U_{k}=\pi^{-1}\left(V_{k}\right)$. It is also clear that $U_{1} \subset U_{2} \subset U_{3} \subset \cdots$ and that $\bigcup_{k=1}^{\infty} U_{k}=X_{l}^{\infty}$.
Lemma 4.8. For each $k, H_{2}\left(U_{k}\right) \cong H_{2}\left(X_{l}^{k}\right)$
Proof: Notice that $U_{k}$ canonically isomorphic to $X_{l}^{k}-\bigcup_{n=k}^{\infty}\left\{N^{-n}(p), N^{-n}(q)\right\}$. Removing a discrete set of points from a 4 (real) dimensional manifold does not affect the second homology. Hence, $H_{2}\left(U_{k}\right) \cong H_{2}\left(X_{l}^{k}\right)$.

Lemma 4.9. $H_{2}\left(X_{l}^{k}\right) \cong \mathbb{Z}^{\left(L_{k} \cup\{[V]\}\right)}$, where $L_{k}$ is the set of fundamental classes of exceptional divisors in $X_{l}^{k}$.
Proposition 4.10. $H_{2}\left(X_{l}^{\infty}\right) \cong \mathbb{Z}^{(L \cup\{[V]\})}$, where $L$ is the set of fundamental classes of exceptional divisors in $X_{l}^{\infty}$ and $[V]$ is the fundamental class of the vertical line $V$.
Proof: Since $X_{l}^{\infty}=\bigcup_{k=1}^{\infty} U_{k}$ and $H_{2}\left(U_{k}\right) \cong H_{2}\left(X_{l}^{k}\right) \cong \mathbb{Z}^{(L \cup\{[V]\})}$, we have that $H_{2}\left(X_{l}^{\infty}\right) \cong$ $\xrightarrow{\lim }\left(\mathbb{Z}^{\left(L_{k} \cup\{[V]\}\right)}\right)$, which is $\mathbb{Z}^{(L \cup\{[V]\})}$.

In the generic case $B \in S$ for which none of the inverse images of $p$ or $q$ under $N$ are in the critical value parabola $C$ we can describe $H_{2}\left(X_{l}^{\infty}\right)$ somewhat more explicitly:

Proposition 4.11. For $B \in S, H_{2}\left(X_{l}^{\infty}\right)=\mathbb{Z}^{\{[V]\}} \oplus\left(\bigoplus_{N^{k}(x)=p} \mathbb{Z}^{\left\{\left[E_{x}\right]\right\}}\right) \oplus\left(\bigoplus_{N^{k}(x)=q} \mathbb{Z}^{\left\{\left[E_{x}\right]\right\}}\right)$.
Proof: This is merely a restatement of Proposition 4.10 using that when $B \in S$, only a single blow-up is necessary at each $k$-th inverse image of $p$ and of $q$ for every $k$.

We will need the following proposition about the intersection of classes in $H_{2}\left(X_{l}^{\infty}\right)$ :
Proposition 4.12. Let $[V]$ and $\left[E_{z}\right]$ be the fundamental classes of $a$ vertical line $V$ and an exceptional divisor $E_{z}$ in $H_{2}\left(X_{l}^{\infty}\right)$. Then $[V] \cdot[V]=0$ and $\left[E_{z}\right] \cdot\left[E_{z}\right] \leq-1$.

Proof: We have chosen the vertical line $V$ so that points on it are never blown up, hence within $X_{l}^{\infty}$ it has self-intersection number 0 , just as it did in $X_{l}$.

If no points on the exceptional divisor $E_{z}$ have been blown up, then it is a classical result that $\left[E_{z}\right] \cdot\left[E_{z}\right]=-1$. Otherwise, if points in $E_{z}$ have been blown up, it is a classical result that each blow-up reduces $\left[E_{z}\right] \cdot\left[E_{z}\right]$ by 1 , hence $\left[E_{z}\right] \cdot\left[E_{z}\right] \leq-1$. (See [28].)

### 4.4 Why we work in $X_{l}$

Suppose for a moment that we did this sequence of blow-ups in $X=\mathbb{P}^{1} \times \mathbb{P}^{1}$, instead of just within the invariant subspace $X_{l}$. The repeated inverse images of the points of indeterminacy must accumulate because $X$ is compact. The topology of the inverse limit $X^{\infty}$ becomes very complicated
near these points of accumulation. Hubbard and Papadopol develop elaborate techniques including Farey Blow-ups and Real-oriented Blow-ups to "tame" the topology at these points of accumulation. We avoid this problem caused by accumulation by working in the space $X_{l}$ since the inverse images of $p$ and $q$ go to the "end" of $X_{l}=\{\operatorname{Re}(x)<1 / 2\} \times \mathbb{P}^{1}$ instead of accumulating to a finite point.

## 5 Mayer-Vietoris sequences

We will study the topology of $W_{0}$ and $W_{1}$ in detail in order to prove Theorem 0.1 . The following Mayer-Vietoris calculations will allow us to relate their homology to that of the basins of attraction for the four roots $r_{1}, r_{2}, r_{3}$, and $r_{4}$.

Let $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ be the closures of $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ in $X_{l}^{\infty}$ and let $\overline{W\left(r_{3}\right)}$ and $\overline{W\left(r_{4}\right)}$ be the closures of $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$ in $X_{r}$. Since $W_{0}$ and $W_{1}$ are real-analytic varieties in $X_{l}^{\infty}$ and $X_{r}$, respectively, there are neighborhoods in $X_{l}^{\infty}$ and $X_{r}$ of $W_{0}$ and $W_{1}$ that deformation retract onto $W_{0}$ and $W_{1}$. Hence, we can consider the Mayer-Vietoris exact sequence (see [13, 31]) for the decompositions $\overline{W\left(r_{1}\right)} \cup \overline{W\left(r_{2}\right)}=X_{l}^{\infty}, \overline{W\left(r_{1}\right)} \cap \overline{W\left(r_{2}\right)}=W_{0}$ and $\overline{W\left(r_{3}\right)} \cup \overline{W\left(r_{4}\right)}=X_{r}$, $\overline{W\left(r_{3}\right)} \cap \overline{W\left(r_{4}\right)}=W_{1}$. We find that

$$
\begin{align*}
& 0 \rightarrow H_{2}\left(W_{0}\right) \xrightarrow{i_{1 *} \oplus i_{2 *}} H_{2}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{2}\left(\overline{W\left(r_{2}\right)}\right) \\
& \xrightarrow{j_{1 *}-j_{2 *}} H_{2}\left(X_{l}^{\infty}\right) \xrightarrow{\partial} H_{1}\left(W_{0}\right) \xrightarrow{i_{1 *} \oplus i_{2 *}} H_{1}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{2}\right)}\right) \rightarrow 0 \tag{14}
\end{align*}
$$

is exact, where $i_{1}$ and $i_{2}$ are the inclusions of $W_{0}$ into $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ and $j_{1}$ and $j_{2}$ are the inclusions of $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ into $X_{l}^{\infty}$. Slightly more work shows that $\partial[V]=\left[S_{0}\right]$, where $[V]$ is the fundamental class of a vertical line in $X_{l}^{\infty}$ and $\left[S_{0}\right]$ is the class of the invariant circle.

We have $H_{2}\left(X_{r}\right)=\mathbb{Z}^{\{[\mathbb{P}]\}}$ with $\partial([\mathbb{P}])=\left[S_{1}\right]$. Using that $H_{i}\left(X_{r}\right)=0$ for $i \neq 2,0$, we find that the map

$$
\begin{equation*}
H_{2}\left(W_{1}\right) \xrightarrow{i_{3 *} \oplus i_{4 *}} \quad H_{2}\left(\overline{W\left(r_{3}\right)}\right) \oplus H_{2}\left(\overline{W\left(r_{4}\right)}\right) \tag{15}
\end{equation*}
$$

is an isomorphism and the sequence

$$
\begin{equation*}
0 \rightarrow \mathbb{Z}^{\{[\mathbb{P}]\}} \xrightarrow{\partial} H_{1}\left(W_{1}\right) \xrightarrow{i_{3 *} \oplus i_{4 *}} H_{1}\left(\overline{W\left(r_{3}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{4}\right)}\right) \rightarrow 0 \tag{16}
\end{equation*}
$$

is exact, where where $i_{3}$ and $i_{4}$ are the inclusions of $W_{1}$ into $\overline{W\left(r_{3}\right)}$ and $\overline{W\left(r_{4}\right)}$.

## 6 Morse Theory for $W_{1}$ and $W_{0}$

In this section we prove that if there are parameter values $B$ for which $W_{1} \cap C=\emptyset$, then $H_{1}\left(\overline{W\left(r_{3}\right)}\right)$ and $H_{1}\left(\bar{W}\left(r_{4}\right)\right)$ are trivial. Our computer experiments suggest that such $B$ exist, but we have not proven their existence. This proves the second half of the second part of Theorem 0.1 , which we will finish proving in the following two sections.

Consider the function $h: \mathbb{C} \times \mathbb{P} \rightarrow \mathbb{R}$ given by $h\binom{x}{y}=\left|\frac{x}{x-1}\right|$ which is chosen so that

$$
\begin{equation*}
h\left(N\binom{x}{y}\right)=\left|\frac{x^{2}}{x^{2}-2 x+1}\right|=h\left(\binom{x}{y}\right)^{2} . \tag{17}
\end{equation*}
$$

We will consider the restriction of $h$ to the super-stable separatrices $W_{0}$ and $W_{1}$ and use it as a Morse function to study their topology. Because $W_{0}$ and $W_{1}$ intersect the critical value parabola $C$ in real-analytic sets, the following geometric description of the critical points of $h$ makes sense:


Figure 9: Level curves of the Morse function $h$ in part of the critical value parabola $C$. The points labeled $a$ and $b$ on $W_{0}$ are in $K$ and the point labeled $c$ on $W_{1}$ is in $K$. Repeated inverse images of these points under $N$ are critical points of $h$ on $W_{0}$ and $W_{1}$. Clearly we have only labeled a few of the points in $K$ that are visible.

Proposition 6.1. Let $K$ be the set of points in $C$ where $W_{0} \cap C$ (or $W_{1} \cap C$ ) is parallel to the level curves of $h \mid C$. Then the set of critical points of $h$ on $W_{0}$ and $W_{1}$ is $\bigcup_{k=1}^{\infty} N^{-k}(K)$.

Proof: Applying the chain rule to Equation 17 we find:

$$
\begin{equation*}
D h\left(N\binom{x}{y}\right) \cdot D N\binom{x}{y}=2 h\binom{x}{y} \cdot D h\binom{x}{y} \tag{18}
\end{equation*}
$$

Because $\left|h\binom{x}{y}\right|=0$ only when $x=0$, Equation 18 gives that if $D h\binom{x}{y}=0$ for a point $\binom{x}{y}$ on $W_{0}$ then either:

1. $\operatorname{Dh}\left(N\binom{x}{y}\right)=0$ giving that $\binom{x}{y}$ is an inverse image (possibly an $n$-th inverse image) of another critical point of $h$. Or,
2. $D N\binom{x}{y}$ is singular and $D h\left(N\binom{x}{y}\right)$ is killed by $D N\binom{x}{y}$.

The condition in the second case says that $(x, y)$ is on the critical points locus of $N$ and that the curve $W_{0} \cap C$ is tangent to the level curves of $h \mid C$ at $N(x, y)$.

Notice that if $h: W_{0} \rightarrow \mathbb{R}$ or if $h: W_{1} \rightarrow \mathbb{R}$ has no critical points aside from at $x=0$ or 1 , then the negative gradient flow $-\nabla h$ gives a deformation retraction of $W_{0}$ to $S_{0}$ or the gradient flow $\nabla h$ gives a deformation retraction of $W_{1}$ to $S_{1}$.

Proposition 6.2. If there are no points of intersection between $W_{1}$ and the critical value parabola $C$, then $W_{1}$ is homotopy equivalent to $S_{1}$.

Proof: By Proposition 6.1 if $W_{1}$ and $C$ are disjoint, there are no critical points of $h$
Corollary 6.3. If there are no points of intersection between $W_{1}$ and the critical value parabola $C$, then the basins of attraction $W\left(r_{3}\right)$ and $W\left(r_{4}\right)$ for the roots $r_{3}=(1,0)$ and $r_{4}=(1,1-B)$ have trivial first and second homology groups.

Proof: This follows for the second homology from the isomorphism between $H_{2}\left(W_{1}\right)$ and $H_{2}\left(\overline{W\left(r_{3}\right)}\right) \oplus$ $H_{2}\left(\overline{W\left(r_{4}\right)}\right)$. For the first homology, $H_{1}\left(W_{1}\right) \cong \mathbb{Z}^{\left\{\left[S_{1}\right]\right\}}=$ Image $(\partial)$ and exactness of (16) gives that
$H_{1}\left(\overline{W\left(r_{3}\right)}\right)=0=H_{1}\left(\overline{W\left(r_{4}\right)}\right)$. Because $W_{1}$ is disjoint from $C, B$ is not in the bifurcation locus, we can replace $\overline{W\left(r_{3}\right)}$ and $\overline{W\left(r_{4}\right)}$ with the basins themselves.

Proposition 6.4. There are always critical points of $h: W_{0} \rightarrow \mathbb{R}$.
Proof: Using implicit differentiation of $C(x, y)=0$, one can check that there is a unique critical point of $h \mid C$ at the intersection of $C$ with the line $B x+2 y-1=0$. Since this line is the axis of symmetry for $\tau$, it is in $W_{0}$.

Instead of studying the Morse function $h$ when $W_{0}$ or $W_{1}$ intersects $C$, in the next two sections we will use linking numbers to show that such intersections result in infinitely many loops corresponding to distinct generators of homology.

## 7 Many loops in $W_{0}$ and $W_{1}$

Denote the vertical line in $X_{l}$ at a fixed value of $x$ by $\mathbb{P}_{x}$. Such vertical lines in $X_{l}$ correspond naturally to lines in $X_{l}^{\infty}$ by means of the "proper transform" that is induced by the blow-up operation.

The Newton Map $N$ maps $\mathbb{P}_{x}$ to $\mathbb{P}_{x^{2} /(2 x-1)}$ by the rational map:

$$
R_{x}(y)=\frac{y\left(B x^{2}+2 x y-B x-y\right)}{(2 x-1)(B x+2 y-1)} .
$$

Notice that when $x=\frac{1}{B}$ and when $x=\frac{1}{2-B}$, a common term cancels from the numerator and denominator of $R_{x}$, giving $R_{x}(y)=\frac{y}{2}+\frac{1-B}{2(2-B)}$ and $R_{x}(y)=\frac{y}{2}$, respectively. The critical values of $R_{x}$ are the intersections of the critical value parabola $C$ with the line $\mathbb{P}_{x^{2} /(2 x-1)}$. There are two distinct critical values, except when $x=\frac{1}{B}$ or $\frac{1}{2-B}$.
Lemma 7.1. There are $\epsilon_{0}>0$ and $\epsilon_{1}>0$ so that if $|x-0|<\epsilon_{0}$, then $W_{0} \cap \mathbb{P}_{x}$ forms a simple closed curve and so that if $|x-1|<\epsilon_{1}$, then $W_{1} \cap \mathbb{P}_{x}$ forms a simple closed curve.

Proof: This is a direct consequence of the existence of $W_{0}^{\text {loc }}$ and $W_{1}^{\text {loc }}$ and the fact that there is no possible recurrence in the dynamics for $x$ within $X_{l}^{\infty}$ or $X_{r}$.

Most vertical lines $\mathbb{P}_{x}$ will be divided by $W_{i}(i=0$ or 1$)$ into exactly two simply connected domains. However, if $W_{i} \cap C$ is non-empty in any forward image of $\mathbb{P}_{x}$, then $W_{i}$ will divide $\mathbb{P}_{x}$ into many more simply connected domains. These are counted in the following proposition.

Proposition 7.2. Let $\mathbb{P}_{x}$ be a vertical line whose $k$-th forward image $\mathbb{P}_{\hat{x}}$ is divided by $W_{i}$ into exactly two simply connected domains. If $W_{i} \cap C \neq \emptyset$ in $\mathbb{P}_{\hat{x}}$, then $\mathbb{P}_{x}$ is divided by $W_{i}$ into between $2^{k}+2$ and $2^{k+1}$ simply connected domains.

We prove Proposition 7.2 for $W_{0}$ in $X_{l}^{\infty}$ because the proof is identical in $X_{r}$. The following lemma and corollary are direct consequences of the Riemann-Hurwitz Theorem.

Lemma 7.3. Let $R: \mathbb{P} \rightarrow \mathbb{P}$ be a ramified covering map of degree $d$ and let $U \subset \mathbb{P}$ be a simply connected open subset of $\mathbb{P}$ containing the image of at most one point of ramification of $R$. Then, $R^{-1}(U)$ consists of a finite number of disjoint simply connected domains.

The symmetry (1.6) guarantees that there is at most one of the two critical values of $R_{x}$ is in each simply connected domain, hence the inverse image of each domain is a finite number of simply connected domains. An easy check shows that if $U$ contains one of the critical values of $R_{x}$, then $R_{x}^{-1}(U)$ is a single simply connected domain, while if $U$ does not contain a critical value of $R_{x}$ it is two simply connected domains.

Corollary 7.4. Let $\mathbb{P}_{x}$ be a vertical line whose image $\mathbb{P}_{\hat{x}}$ is divided by $W_{i}$ into $m$ simply connected domains. If $W_{i} \cap C \neq \emptyset$ in $\mathbb{P}_{\hat{x}}$, then $W_{i}$ divides $\mathbb{P}_{x}$ into $2 m$ simply connected domains. Otherwise it divides $\mathbb{P}_{x}$ into $2 m-2$ simply connected domains.

The proof of Proposition 7.2 follows from this Corollary and the fact that there is a sufficiently high $k$ so that $\left|x_{k}\right|<\epsilon$ so that $W_{0}$ divides $\mathbb{P}_{x_{k}}$ into exactly two domains.


Figure 10: Left: a vertical line divided by $W_{0}$ into 10 regions. Loop $\gamma$ has size $(\gamma)=1 / 10$ and loop $\tau$ has $\operatorname{size}(\tau)=1 / 10$. Right: a different vertical line that is divided by $W_{0}$ into 10 regions. This time, the loops bounding the regions are much more ornate.

Suppose that $W_{0}$ divides $\mathbb{P}_{x}$ into $2 m$ simply connected domains $U_{1}, \cdots, U_{m}$ in $W\left(r_{1}\right)$ and $V_{1}, \cdots, V_{m}$ in $W\left(r_{2}\right)$. Let $k$ be chosen so that $W_{0}$ forms a simple closed curve in $\mathbb{P}_{\tilde{x}}$ (where $\tilde{x}$ is the $k$-th iterate of $x$ under $x \mapsto \frac{x^{2}}{2 x-1}$.) Denote by $U$ the domain in $\mathbb{P}_{\tilde{x}}$ within $W\left(r_{1}\right)$ and by $V$ the domain in $\mathbb{P}_{\tilde{x}}$ within $W\left(r_{2}\right)$.

Under the mapping $N^{k}$, each domain $U_{i}$ covers $U$ with some degree $l_{i}$ and each domain $V_{j}$ covers $V$ with degree $p_{j}$. Then: $\sum_{i=1}^{m} l_{i}=2^{k}, \quad \sum_{i=1}^{m} p_{i}=2^{k}$ because $U$ is covered by $\cup_{i=1}^{m} U_{i} \subset \mathbb{P}_{x}$ with degree $2^{k}$.

For such a $U_{i}$ we will assign $\operatorname{size}\left(U_{i}\right)=-\frac{l_{i}}{2^{k}}$ and such a $V_{i}$ we can assign $\operatorname{size}\left(U_{i}\right)=\frac{p_{i}}{2^{k}}$. This is well defined because given $k_{1}$ and $k_{2}$ as above, the $l_{i}$ corresponding to $k_{1}$ and the $l_{i}$ corresponding to $k_{2}$ will differ by $2^{k_{1}-k_{2}}$.

$$
\sum_{i=1}^{m} \operatorname{size}\left(U_{i}\right)=-1, \quad \sum_{i=1}^{m} \operatorname{size}\left(V_{i}\right)=1
$$

In the next section we will see that $\operatorname{size}\left(U_{i}\right)$ for such a region equals the linking number between $\gamma_{i}=\partial U_{i}$ and an appropriate geometric object in $X_{l}^{\infty}$.

## 8 Linking numbers

Classically one considers the linking number of two oriented loops $c$ and $d$ in $\mathbb{S}^{3}$. The linking number $l k(c, d) \in \mathbb{Z}$ is found by taking any oriented surface $\Gamma$ with oriented boundary $c$ and defining $l k(c, d)$ to be the signed intersection number of $\Gamma$ with $d$ as in Figure 8. For this and many equivalent definitions of linking number in $\mathbb{S}^{3}$ see [45, pp. 132-133], [12, pp. 229-239], and [43, Problems 13 and 14].

To see that this linking number is well-defined notice that assigning $l k(c, d)=[\Gamma] \cdot[d]$, where . indicates the intersection product on $H_{*}\left(\mathbb{S}^{3}, c\right)$, coincides with the classical definition. If $\Gamma^{\prime}$ is any


Figure 11: Here $l k(c, d)=+2$.
other 2-chain with $\partial \Gamma^{\prime}=c$ then $\partial\left(\Gamma-\Gamma^{\prime}\right)=[c]-[c]=0$ and $\left(\Gamma-\Gamma^{\prime}\right)$ represents a homology class in $H_{2}\left(\mathbb{S}^{3}\right)$. Since $H_{2}\left(\mathbb{S}^{3}\right)=0,\left[\Gamma-\Gamma^{\prime}\right]=0$ forcing $\left[\Gamma-\Gamma^{\prime}\right] \cdot[d]=0$. Therefore: $[\Gamma] \cdot[d]=\left[\Gamma^{\prime}\right] \cdot[d]$, so that $l k(c, d)$ is well defined.
Linking kernel: $\mathcal{L} Z_{p}(M)$
Suppose that $M$ is a 3 -dimensional manifold with $H_{2}(M) \neq 0$. We can define a linking number $l k(c, d)$ so long as the second argument $d$ has $[d] \cdot \sigma=0$ for every $\sigma \in H_{2}(M)$. We define $\mathcal{L} Z_{1}(M) \subset$ $Z_{1}(M)$ to be the sub-module of one chains having this property. As before, given $d \in \mathcal{L} Z_{1}(M)$ the linking number $l k(c, d)$ is well-defined for $c$ disjoint from $d$ with $[c]=0$.

Linking numbers work in a similar way if a manifold $M$ has dimension $m$ : one requires that $c$ and $d$ have dimensions summing to $m-1$ and one must restrict to $c$ disjoint from $d$ with $[c]=0$ and restrict to $d \in \mathcal{L} Z_{p}(M)$, where $p$ is the dimension of $d$.

### 8.1 Linking kernel for $X_{l}^{\infty}$

Recall from Section 4 that

$$
H_{2}\left(X_{l}^{\infty}\right)=\mathbb{Z}^{(L \cup\{[V]\})}
$$

where $L$ is the set of exceptional divisors $E_{i}$ introduced in the sequence of blow-ups and $V$ is the vertical line $x=\frac{1}{B(2-B)}$.

Recall from Proposition 4.12 that each exceptional divisor $\left[E_{i}\right]$ has $\left[E_{i}\right] \cdot\left[E_{i}\right] \leq-1$ and that $[V] \cdot[V]=0$ so that if $\omega=a_{0}[V]+a_{1}\left[E_{1}\right]+\cdots+a_{n}\left[E_{n}\right]$ satisfies $\omega \cdot \sigma=0$ for every $\sigma \in H_{2}\left(X_{l}^{\infty}\right)$ then $a_{i}=0$ for all $i \neq 0$.

In summary, $\mathcal{L} Z_{2}\left(X_{l}^{\infty}\right)=\mathbb{Z}^{\{V\}}$. The curves $\gamma_{i}$ considered in the previous section each have linking number 0 with $V$, since each $\gamma_{i}$ is entirely within some (other) vertical line. To show that any of the curves $\gamma_{i}$ are non-trivial we will need to look for a different kind of object to link with. We do this by extending the definition of linking to closed currents.

### 8.2 Generalities on currents

Just as distributions are defined as the topological dual of smooth functions with compact support, currents are the topological dual of smooth differential forms with compact support.

Let $A_{c}^{n-q}(M)$ denote the $(n-q)$-forms with compact support on a smooth manifold $M$. The linear maps $T: A_{c}^{n-q}(M) \rightarrow \mathbb{C}$ that are continuous are the currents of degree $q$ (or, as some say, the currents of dimension $n-q$ ) and are denoted by $\mathcal{D}^{q}(M)$. If $M$ has a complex structure, one defines the currents of bi-degree $(p, q)$, denoted $\mathcal{D}^{p, q}(M)$ as the topological dual of the $(n-p, n-q)$ forms with compact support $A_{c}^{n-p, n-q}(M)$. For more background on currents, consult [28, section 3.1 and 3.2] or the articles on complex dynamics [38, 49, 48].

An exterior derivative $d: \mathcal{D}^{q}(M) \rightarrow \mathcal{D}^{q+1}(M)$ is defined as the adjoint to the classical exterior derivative on smooth forms with compact support: $d T(\eta)=-T(d \eta)$.

On a complex manifold, one has two derivatives $\partial: \mathcal{D}^{p, q}(M) \rightarrow \mathcal{D}^{p+1, q}(M)$ and $\bar{\partial}: \mathcal{D}^{p, q}(M) \rightarrow$ $\mathcal{D}^{p, q+1}(M)$, defined in the analogous way. However, the real operators $d=\partial+\bar{\partial}$ and $d^{c}=\frac{i}{2 \pi}(\bar{\partial}-\partial)$
are more often used in dynamics. Currents that satisfy $d T=0$ are referred to as $d$-closed. We will denote the $d$-closed currents of degree $q$ by $Z^{q}(M)$, and when desiring to emphasize bi-degree, we will denote the $d$-closed currents of bi-degree $p, q$ by $Z^{p, q}(M)$.

Given a smooth form $\psi \in A^{q}(M)$, there is a current $T_{\psi} \in \mathcal{D}^{q}(M)$ defined by $T_{\psi}(\eta)=\int_{M} \psi \wedge \eta$ for any $\eta \in A_{c}^{n-q}(M)$. Currents of this form are referred to as smooth currents. Using Stokes Theorem, one can check that $d T_{\psi}=T_{d \psi}$ so that the inclusion $A_{c}^{n-q}(M) \rightarrow \mathcal{D}^{q}(M)$ given by $\psi \mapsto T_{\psi}$ is a cochain map.

A piecewise smooth, oriented $(n-q)$ chain $\Gamma$ in $M$ also defines a current $T_{\Gamma} \in \mathcal{D}^{q}(M)$ given by $T_{\Gamma}(\eta)=\int_{\Gamma} \eta$ for any $\eta \in A_{c}^{n-q}(M)$. We will refer to currents that can be represented this way as currents of integration.

We will often use work with closed, positive $(\mathbf{1}, \mathbf{1})$ currents. These are $(1,1)$ currents that are locally expressed as $T=d d^{c} \phi$ for a plurisubharmonic function $\phi$. (See the $d d^{c}$-Poincaré Lemma.) We denote the closed-positive $(1,1)$ currents on $M$ by $Z_{+}^{1,1}(M)$.

Typically one cannot pull back a current under a ramified mapping $F$. One very special property of positive closed currents is that they can be pulled-back, using the potential function: $F^{*}(\lambda):=$ $d d^{c}(\phi \circ F)$, where $\phi$ is a potential function for $\lambda$. When $\lambda=T_{\eta}$ is a smooth current, this pull-back coincides with the classical pull-back of smooth forms: $F^{*}\left(T_{\eta}\right)=T_{F^{*}(\eta)}$.

### 8.3 Linking with currents

The operator $d: \mathcal{D}^{q}(M) \rightarrow \mathcal{D}^{q+1}(M)$ satisfies $d \circ d=0$ and we denote the corresponding cohomology theory by $H^{*}\left(\mathcal{D}^{*}(M), d\right)$. There is a natural map from the DeRham cohomology $H_{D R}^{*}(M)$ into $H^{*}\left(\mathcal{D}^{*}(M), d\right)$ induced by the inclusion of smooth forms into the currents.
Theorem 8.1. (Approximation by smooth currents) The map $H_{D R}^{*}(M) \rightarrow H^{*}\left(\mathcal{D}^{*}(M), d\right)$ is an isomorphism. Furthermore, the cohomology class of any closed current $L$ can be represented by a closed smooth form $\eta_{L}$ with support in an arbitrarily small neighborhood of the support of $L$.

See [28, pages 382-385] for a proof.
Given $T \in Z^{2}(M)$, and a piecewise smooth 2-chain $\sigma$ having $\partial \sigma$ disjoint from the support of $T$, there is a pairing:

$$
C_{2}(M) \times Z^{2}(M) \rightarrow \mathbb{R}
$$

defined by $\langle\sigma, T\rangle=\int_{\sigma} \eta_{T}$ were $\eta_{T}$ is a smooth form within the same cohomology class as $T$ with support is bounded away from $\partial \sigma$. The existence of $\eta_{T}$ is garunteed by the approximation by smooth currents, and the pairing is well defined since the integral depends only on the cohomology class of $\eta_{T}$.

When $T$ is a current of integration integration over a piecewise smooth chain this pairing coincides with the usual intersection number of piecewise smooth chains and when $T$ is given by a smooth form, it coincides, by definition, with the standard pairing $\int_{\sigma} T$. (In fact, our pairing is a special case of the general intersection number for closed currents of complimentary degrees [28, p. 392] and [17].)
Proposition 8.2. If $\lambda$ is a positive closed current and $F$ is a ramified mapping, we have $\left\langle F_{*} \sigma, \lambda\right\rangle=$ $\left\langle\sigma, F^{*} \lambda\right\rangle$.
Proof: Let $\eta_{\lambda}$ be a smooth approximation of $\lambda$ in the same cohomology class. Then, $\left\langle F_{*} \sigma, \lambda\right\rangle=$ $\int_{F_{*} \sigma} \eta_{\lambda}=\int_{\sigma} F^{*} \eta_{\lambda}=\left\langle\sigma, F^{*} \lambda\right\rangle$, since $F^{*} \eta_{\lambda}$ is a smooth approximation of $F^{*} \lambda$.

We define the linking kernel $\mathcal{L} Z^{2}(M)$ to be the space of closed currents $T$ having $\langle\sigma, T\rangle=0$ for every $\sigma \in H_{2}(M)$. Given $T \in \mathcal{L} Z^{2}(M)$, let $B_{1}^{T}(M)$ be the 1-boundaries in $M$ that are disjoint from the support of $T$. We can define a linking number with respect to $T$

$$
l k(\cdot, T): B_{1}^{T}(M) \rightarrow \mathbb{R}
$$

by $l k(c, T)=\langle\Gamma, T\rangle$, where $\Gamma$ is any 2-chain with $\partial \Gamma=c$. Since $T \in \mathcal{L} Z^{2}(M)$, we have that $\langle\Gamma, T\rangle=\left\langle\Gamma^{\prime}, T\right\rangle$ for any other $\Gamma^{\prime}$ with $\partial \Gamma^{\prime}=c$.

### 8.4 Finding an element of $\mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$

In this subsection, we will find an element of $\mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$ by successively determining elements of $\mathcal{L} Z^{2}\left(X_{l}\right), \mathcal{L} Z^{2}\left(X_{l}^{0}\right), \mathcal{L} Z^{2}\left(X_{l}^{1}\right), \mathcal{L} Z^{2}\left(X_{l}^{2}\right), \cdots$ where $X_{k}^{j}$ is the space $X_{k}$ after having completed the blow-ups at level $j$. In the limit, we will find an element of $\mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$, which in the next subsection will be useful for linking.

Let $L_{1}$ be the invariant line that goes through $(0,0)$ and $(1,0)$, i.e. $y=0$ and $L_{2}$ be the invariant line that goes through $(0,1)$ and $(1,1-B)$, i.e. $y+B x-1=0$. (To remember the indexing, think that $L_{1}$ contains $r_{1}$ and $L_{2}$ contains $r_{2}$.) We can use the Poincaré-Lelong formula ( $[28, \mathrm{p} .388]$ or [49]) to express the fundamental classes of these lines as positive-closed currents:

$$
\left[L_{1}\right]=\frac{1}{2 \pi} d d^{c} \log |y|, \quad\left[L_{2}\right]=\frac{1}{2 \pi} d d^{c} \log |y+B x-1| .
$$

Both $L_{1}$ and $L_{2}$ intersect any given vertical line $\mathbb{P}$ with intersection number 1 . Because $[V]$ is the sole generator of $H_{2}\left(X_{l}\right)$ we have that $\left[L_{2}\right]-\left[L_{1}\right] \in \mathcal{L} Z^{2}\left(X_{l}\right)$.

Now, suppose that we want to find an element of $\mathcal{L} Z^{2}\left(X_{l}^{0}\right)$, that is, a closed 2 current that evaluates to 0 on every element of $H_{2}\left(X_{l}^{0}\right) \cong \mathbb{Z}^{\left\{[V],\left[E_{p}\right],\left[E_{q}\right]\right\}}$. In fact:

$$
\left\langle E_{p},\left[L_{1}\right]\right\rangle=1=\left\langle E_{p},\left[L_{2}\right]\right\rangle \text { and }\left\langle E_{q},\left[L_{1}\right]\right\rangle=0=\left\langle E_{q},\left[L_{2}\right]\right\rangle,
$$

using standard intersection numbers for piecewise smooth chains, so that $\left[L_{2}\right]-\left[L_{1}\right] \in \mathcal{L} Z^{2}\left(X_{l}^{0}\right)$.
This luck will not continue. Let $z$ be one of the two preimages of $p$ that is in the invariant line $L_{1}$. Since $L_{1}$ and $L_{2}$ intersect at the single point $p$, this forces that $z \notin L_{2}$. Consequently: $\left\langle E_{z},\left[L_{1}\right]\right\rangle=1 \neq 0=\left\langle E_{z},\left[L_{2}\right]\right\rangle$ so that $\left[L_{2}\right]-\left[L_{1}\right] \notin \mathcal{L} Z^{2}\left(X_{l}^{1}\right)$.

We consider the $k$-th inverse images $N^{-k}\left(L_{1}\right)$ and $N^{-k}\left(L_{2}\right)$. If we denote by $N_{1}^{k}(x, y)$ and $N_{2}^{k}(x, y)$ the first and second coordinates of $N^{k}$, then the Poincaré-Lelong formula gives

$$
\begin{aligned}
{\left[N^{-k}\left(L_{1}\right)\right] } & =\frac{1}{2 \pi} d d^{c} \log \left|N_{2}^{k}(x, y)\right|, \\
{\left[N^{-k}\left(L_{2}\right)\right] } & =\frac{1}{2 \pi} d d^{c} \log \left|N_{1}^{k}(x, y)+B \cdot N_{2}^{k}(x, y)-1\right| .
\end{aligned}
$$

Lemma 8.3. For every $k \geq 0$ we have

$$
\left\langle V,\left[N^{-k}\left(L_{1}\right)\right]\right\rangle=\left\langle V,\left[N^{-k}\left(L_{2}\right)\right]\right\rangle
$$

Proof: The $k$-th inverse images $N^{-k}\left(L_{1}\right)$ and $N^{-k}\left(L_{2}\right)$ both have degree $2^{k}$ in $y$, so they each intersect a generic vertical line transversely exactly $2^{k}$ times. These intersection numbers coincide with the pairings.

Proposition 8.4. $\left[N^{-(k+1)}\left(L_{2}\right)\right]-\left[N^{-(k+1)}\left(L_{1}\right)\right] \in \mathcal{L} Z^{2}\left(X_{l}^{k}\right)$
Proof Let $E_{z}$ be any one of the exceptional divisors in $X_{l}^{k}$. Using Proposition 4.6, there is some $d$ and some $l \leq k+1$ so that $N^{l}$ maps $E_{z}$ to $V$ by a ramified cover of degree $d$ (possibly with $d=0$.) Just as in the discussion above:

$$
\begin{aligned}
& \left\langle E_{z},\left[N^{-(k+1)}\left(L_{1}\right)\right]\right\rangle=\left\langle N^{l}\left(E_{z}\right),\left[N^{-(k+1)+l} L_{1}\right]\right\rangle=d\left\langle V,\left[N^{-(k+1)+l} L_{1}\right]\right\rangle \\
& \left\langle E_{z},\left[N^{-(k+1)}\left(L_{2}\right)\right]\right\rangle=\left\langle N^{l}\left(E_{z}\right),\left[N^{-(k+1)+l} L_{2}\right]\right\rangle=d\left\langle V,\left[N^{-(k+1)+l} L_{2}\right]\right\rangle
\end{aligned}
$$

Here we are using Proposition 8.2 to obtain the first equality in each equation. (One must check that the Poincaré-Lelong equation gives that $\left[N^{-(k+1)}\left(L_{i}\right)\right]=\left(N^{l}\right)^{*}\left[N^{-(k+1)+l} L_{i}\right]$ for $l \leq k+1$.) Then, Lemma 8.3 gives that the two terms on the right hand side of each equation are equal.

Since $H_{2}\left(X_{l}^{k}\right)$ is generated by the fundamental classes of $V$ and the fundamental classes of each of the exceptional divisors $E_{z}$ we conclude that $\left[N^{-(k+1)}\left(L_{2}\right)\right]-\left[N^{-(k+1)}\left(L_{1}\right)\right] \in \mathcal{L} Z^{2}\left(X_{l}^{k}\right)$.

Since $X_{l}^{\infty}=\lim \left(X_{l}^{k}, \pi\right)$ and $\left[N^{-(k+1)}\left(L_{2}\right)\right]-\left[N^{-(k+1)}\left(L_{1}\right)\right] \in \mathcal{L} Z^{2}\left(X_{l}^{k}\right)$ we expect that a limit as $k \rightarrow \infty$ of $\left[N^{\overleftarrow{-(k+1)}}\left(L_{2}\right)\right]-\left[N^{-(k+1)}\left(L_{1}\right)\right]$ will be an element of $\mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$. For such a limit to converge we must normalize $\left[N^{-(k+1)}\left(L_{2}\right)\right]$ and $\left[N^{-(k+1)}\left(L_{1}\right)\right]$. Dividing by the degrees, we define:

$$
\begin{aligned}
\lambda_{1}^{k} & =\frac{1}{2^{k}}\left[N^{-k}\left(L_{1}\right)\right]=\frac{1}{2 \pi} d d^{c} \frac{1}{2^{k}} \log \left|N_{2}^{k}(x, y)\right| \\
\lambda_{2}^{k} & =\frac{1}{2^{k}}\left[N^{-k}\left(L_{2}\right)\right]=\frac{1}{2 \pi} d d^{c} \frac{1}{2^{k}} \log \left|N_{1}^{k}(x, y)+B \cdot N_{2}^{k}(x, y)-1\right|
\end{aligned}
$$

Let

$$
\begin{aligned}
& \lambda_{1}=\lim _{k \rightarrow \infty} \lambda_{1}^{k}=\frac{1}{2 \pi} d d^{c} \lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{2}^{k}(x, y)\right| \\
& \lambda_{2}=\lim _{k \rightarrow \infty} \lambda_{2}^{k}=\frac{1}{2 \pi} d d^{c} \lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{1}^{k}(x, y)+B \cdot N_{2}^{k}(x, y)-1\right|
\end{aligned}
$$

We will first check that these limits exist and define positive-closed 1-1 currents, and then we will show that $\lambda_{2}-\lambda_{1} \in \mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$.

Proposition 8.5. The limits

$$
\begin{aligned}
& G_{1}(x, y)=\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{2}^{k}(x, y)\right| \\
& G_{2}(x, y)=\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{1}^{k}(x, y)+B \cdot N_{2}^{k}(x, y)-1\right|
\end{aligned}
$$

converge and are plurisubharmonic functions in the basins of attraction $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$, respectively. Hence, $\lambda_{1}=\frac{1}{2 \pi} d d^{c} G_{1}(x, y)$ and $\lambda_{2}=\frac{1}{2 \pi} d d^{c} G_{2}(x, y)$ are positive closed 1-1 currents on $X_{l}^{\infty}$ : $\lambda_{1}, \lambda_{2} \in Z_{+}^{1,1}\left(X_{l}^{\infty}\right)$.

Proof: To see that $G_{1}(x, y)$ and $G_{2}(x, y)$ are well-defined and plurisubharmonic, we will show that $G_{1}(x, y)$ and $G_{2}(x, y)$ coincide with the potential functions that were described by Hubbard and Papadopol in [37, p. 21] and [38]. We will do this for $G_{1}(x, y)$, and leave necessary modifications for $G_{2}(x, y)$ to the reader.

Supposing that $(0,0)$ is a root, Hubbard and Papadopol [37] consider the limit

$$
G_{H P}(x, y)=\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left\|N^{k}(x, y)\right\|
$$

which they show converges to a plurisubharmonic function on the basin of $(0,0)$. The reader should notice that $G_{H P}$ does not depend on the choice of the norm $\|\cdot\|$ used in the definition because any two different norms on a finite dimensional vector space are equivalent by a finite multiplicative constant, which is eliminated by the multiplicative factor of $\frac{1}{2^{k}}$. Therefore, we can use the supremum norm.

We will show that $G_{1}=G_{H P}$ on $W\left(r_{1}\right)$, to see that $G_{1}$ is plurisubharmonic.
If $\left|N_{2}^{k}(x, y)\right| \geq\left|N_{1}^{k}(x, y)\right|$ for all $(x, y)$ as $k \rightarrow \infty$, then the supremum norm coincides with $\left|N_{2}^{k}(x, y)\right|$ giving $G_{1}(x, y)=G_{H P}(x, y)$. This condition is equivalent to the condition:

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|\frac{N_{2}^{k}(x, y)}{N_{1}^{k}(x, y)}\right| \geq 0 \tag{19}
\end{equation*}
$$

which will now show is a consequence of a standard result from the dynamics of one complex variable.
In [37], the authors perform blow-ups at each of the four roots, and observe that the Newton map $N$ induces rational functions of degree 2 on each of the exceptional divisors $E_{r_{1}}, E_{r_{2}}, E_{r_{3}}$, and $E_{r_{4}}$. Let's compute the rational function $s: E_{r_{1}} \rightarrow E_{r_{1}}$. In the coordinate chart $m=\frac{y}{x}$, the extension to $E_{r_{1}}$ is obtained by:

$$
s(m)=\lim _{x \rightarrow 0} \frac{m x\left(B x^{2}+2 m x^{2}-B x-m x\right)}{x^{2}(B x+2 m x-1)}=m(B+m)
$$

since $x=0$ on $E_{r_{1}}$.
Since condition (19) is a limit, it suffices to check it in an arbitrarily small neighborhood of the origin. In a small enough neighborhood, we can replace $\frac{N_{2}^{k}(x, y)}{N_{1}^{k}(x, y)}$ with $s\left(\frac{y}{x}\right)$ obtaining

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|\frac{N_{2}^{k}(x, y)}{N_{1}^{k}(x, y)}\right|=\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|s^{k}(m)\right|=G_{s}(m) \tag{20}
\end{equation*}
$$

where $G_{s}(m)$ is the standard Green's function from one variable complex dynamics associated to the polynomial $s(m)$. This last equality is actually a delicate but well-known result that was proved by Brolin [16]. A more friendly proof is available in [49, Section 9].

Having the last equality, it is a standard result, for example see Milnor [42] pages 95 and 96, that $G_{s}(m)=0$ on the filled Julia set $K(s)$ and that $G_{s}(m)>0$ outside of $K(s)$.

This justifies the replacement of the supremum norm from $G_{H P}$ by $\left|N_{2}^{k}(x, y)\right|$, and hence gives that $G_{1}(x, y)=G_{H P}(x, y)$.

Corollary 8.6. Let $s: E_{r_{1}} \rightarrow E_{r_{1}}$ be the polynomial induced by the Newton map $N$ and let $G_{s}$ : $E_{r_{1}} \rightarrow \mathbb{R}$ be its Green's function. In a sufficiently small neighborhood of $r_{1}$,

$$
G_{1}(x, y)=G_{s}\left(\frac{y}{x}\right)-\log \left|\frac{1}{x}\right|
$$

Proof: This comes directly from the algebra:

$$
\begin{aligned}
G_{1}(x, y) & =\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{2}^{k}(x, y)\right|=\lim _{k \rightarrow \infty} \frac{1}{2^{k}}\left(\log \left|\frac{N_{2}^{k}(x, y)}{N_{1}^{k}(x, y)}\right|+\log \left|N_{1}^{k}(x, y)\right|\right) \\
& =G_{s}\left(\frac{y}{x}\right)+\lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{1}^{k}(x, y)\right|=G_{s}\left(\frac{y}{x}\right)+\log |x|=G_{s}\left(\frac{y}{x}\right)-\log \left|\frac{1}{x}\right|
\end{aligned}
$$

because $\frac{N_{2}^{k}(x, y)}{N_{1}^{k}(x, y)} \approx s\left(\frac{y}{x}\right)$ and $N_{1}^{k}(x, y)=\frac{x^{2}}{2 x-1} \approx-x^{2}$ near $r_{1}$.

### 8.5 Nice properties of $\lambda_{2}$ and $\lambda_{1}$ :

In this subsection, we will prove some of the useful properties if $\lambda_{2}$ and $\lambda_{1}$. We will finish the subsection by showing that $\lambda_{2}-\lambda_{1} \in \mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$.
Lemma 8.7. (Normalization) Suppose that $\mathbb{P}_{x}$ is a vertical line that is divided into exactly two simply connected domains $U \subset W\left(r_{1}\right)$ and $V \subset W\left(r_{2}\right)$ by $W_{0}$. Then:

$$
\left\langle V, \lambda_{2}\right\rangle=1=\left\langle U, \lambda_{1}\right\rangle \text { and }\left\langle U, \lambda_{2}\right\rangle=0=\left\langle V, \lambda_{1}\right\rangle
$$

Proof: Because $N_{2}^{k}(x, y)$ and $B N_{1}^{k}(x, y)+N_{2}^{k}(x, y)-1$ are of degree $2^{k}$ in $y$, both $\lambda_{1}^{k}$ and $\lambda_{2}^{k}$ are normalized to that $\left\langle V, \lambda_{1}^{k}\right\rangle=1$ and $\left\langle V, \lambda_{2}^{k}\right\rangle=1$. Since the potentials for $\lambda_{1}^{k}$ and $\lambda_{2}^{k}$ converge to the potentials for $\lambda_{1}$ and $\lambda_{2}$, we have

$$
\left\langle U, \lambda_{1}\right\rangle=\left\langle U, \lim _{k \rightarrow \infty} \lambda_{1}^{k}\right\rangle=\lim _{k \rightarrow \infty}\left\langle U, \lambda_{1}^{k}\right\rangle=\lim _{k \rightarrow \infty} 1=1
$$

and similarly for $\lambda_{2}$. The proof that $\left\langle U, \lambda_{2}\right\rangle=0=\left\langle V, \lambda_{1}\right\rangle$ is identical.
Corollary 8.8. Suppose that $\mathbb{P}_{x}$ is vertical line, then $\left\langle\mathbb{P}_{x}, \lambda_{2}\right\rangle=1=\left\langle\mathbb{P}_{x}, \lambda_{1}\right\rangle$.

It follows directly from the definitions of $\lambda_{1}$ and $\lambda_{2}$ that $N^{*}\left(\lambda_{1}\right)=2 \cdot \lambda_{1}$ and $N^{*}\left(\lambda_{2}\right)=2 \cdot \lambda_{2}$. (For example $N^{*}\left(\lambda_{1}\right)=\frac{1}{2 \pi} d d^{c} \lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{1}^{k} \circ N\right|=\frac{1}{2 \pi} d d^{c} \lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{1}^{k+1}\right|=2 \cdot \lambda_{1}$.) In combination with Proposition 8.2, this gives:

Lemma 8.9. (Invariance) Suppose that $\Gamma$ is a piecewise smooth 2-chain, then

$$
\left\langle N(\Gamma), \lambda_{1}\right\rangle=2 \cdot\left\langle\Gamma, \lambda_{1}\right\rangle \quad\left\langle N(\Gamma), \lambda_{2}\right\rangle=2 \cdot\left\langle\Gamma, \lambda_{2}\right\rangle
$$

Proposition 8.10. (Support disjoint from $W_{0}$ ) There is a neighborhood $\Theta$ of $W_{0}$ in $X_{l}^{\infty}$ which is disjoint from the support of $\lambda_{1}$ and $\lambda_{2}$.

Proof: By construction, $\lambda_{1}$ has support in $\overline{W\left(r_{1}\right)}$ and $\lambda_{2}$ has support in $\overline{W\left(r_{2}\right)}$. We will find a neighborhood, which we also call $\Theta$, of $W_{0}$ in $\overline{W\left(r_{1}\right)}$ that is disjoint from the support of $\lambda_{1}$. Clearly similar methods will work in $\overline{W\left(r_{2}\right)}$ and the desired neighborhood is the union of the two.

Recall from Corollary 8.6 that $G_{1}(x, y)=G_{s}\left(\frac{y}{x}\right)-\log \left|\frac{1}{x}\right|$, where $G_{s}$ is the Green's function associated to the polynomial $s: E_{r_{1}} \rightarrow E_{r_{1}}$ induced by $N$ at $r_{1}$. Recall that $s(m)=m(B+m)$ in the coordinates $m=\frac{x}{y}$ on $E_{r_{1}}$, so that $m=\infty$ is a superattracting fixed point. (This is the standard situation for a quadratic polynomial.)

It is a standard result from one-variable dynamics, for example see [42] p. 96, that $G_{s}$ is harmonic outside of the Julia set $J(s)$. In particular, $G_{s}$ is harmonic in a neighborhood of $\infty$ (not including $\infty)$. A related standard result that $G_{s}$ has the singularity

$$
G(m)=\log |m|+O(1) \text { as } m \rightarrow \infty
$$

We check that this singularity exactly cancels with $-\log \left|\frac{1}{x}\right|$ coming from $G_{1}(x, y)=G_{s}\left(\frac{y}{x}\right)-\log \left|\frac{1}{x}\right|$ :

$$
\begin{aligned}
G_{1}(x, y) & =\log \left|\frac{y}{x}\right|-\log \left|\frac{1}{x}\right|+O(1) \text { as }\left|\frac{y}{x}\right| \rightarrow \infty \\
& =\log |y|+O(1) \text { as }\left|\frac{y}{x}\right| \rightarrow \infty
\end{aligned}
$$

Therefore, $G_{1}(x, m x)$ is harmonic on a neighborhood $U$ of $m=\infty$, including the point $\infty$. Choose $\theta>0$ so that if $|m|>\theta$, then $G_{1}(x, m x)$ is harmonic.

Let $\Theta_{0}=\left\{(x, y) \in \overline{W\left(r_{1}\right)}\right.$ such that $\left.\left|\frac{y}{x}\right|>\theta\right\}$. This is the open cone of points in $W\left(r_{1}\right)$ with slope to the origin greater than $\theta$. Since the invariant circle $S_{0}$ is above $m=\infty, \Theta_{0}$ is a neighborhood of $S_{0}$ (within $\overline{W\left(r_{1}\right)}$.)

By construction, $\Theta=\bigcup_{n=0}^{\infty} N^{-n}\left(\Theta_{0}\right)$ will be invariant under $N$ and open. Because $\Theta_{0}$ is disjoint from the support of $\lambda_{1}$, the invariance properties for $\lambda_{1}$ from Lemma 8.9 give that all of $\Theta$ must be disjoint from the support of $\lambda_{1}$.

Finally, since $\Theta_{0}$ contains a neighborhood of $S_{0}$, and both $W_{0}$ and $\Theta$ are invariant under $N, \Theta$ forms an open neighborhood of $W_{0}$.

In fact, using the smooth approximation theorem, one can also choose the smooth approximations of $\lambda_{1}$ and $\lambda_{2}$ to have support bounded away from $W_{0}$.

Corollary 8.11. Given any piecewise smooth chain $\sigma \in W_{0}$, we have that $\left\langle\sigma, \lambda_{1}\right\rangle=0$ and $\left\langle\sigma, \lambda_{2}\right\rangle=$ 0 .

Proposition 8.12. $\lambda_{1}-\lambda_{2} \in \mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$

Proof: This proof will be even simpler than the proof of Proposition 8.4 because we directly use the invariance of $\lambda_{1}$ and $\lambda_{2}$ shown in Lemma 8.9.

By Corollary 8.8, we have $\left\langle V, \lambda_{1}\right\rangle=\left\langle V, \lambda_{2}\right\rangle$. Any exceptional divisor $E_{z}$ was created during the blow-ups at some level $k$, and using Proposition 4.6 there is some $l$ so that $N^{\circ}(k+1)$ maps $E_{z}$ to $V=\mathbb{P}_{1 /(B(2-B))}$ by a ramified covering mapping of degree $l$, (possibly $l=0$ ). Then:

$$
\left\langle E_{z}, \lambda_{1}\right\rangle=\frac{l}{2^{k+1}}\left\langle V, \lambda_{1}\right\rangle=\frac{l}{2^{k+1}}\left\langle V, \lambda_{2}\right\rangle=\left\langle E_{z}, \lambda_{2}\right\rangle .
$$

Hence $\left\langle E_{z}, \lambda_{2}-\lambda_{1}\right\rangle=0$ for any exceptional divisor $E_{z}$.
Since an element of $H_{2}\left(X_{l}^{\infty}\right)$ is a linear combination of the fundamental class $[V]$ with a finite number of fundamental classes of exceptional divisors $E_{z}$, we have shown that $\lambda_{2}-\lambda_{1} \in \mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$.

## 8.6 $H_{1}\left(W_{0}\right)$ is infinitely generated.

From Section 7 we have infinitely many cycles $\gamma_{i}$ in $W_{0}$ of arbitrarily small "size," and we now have $\left(\lambda_{2}-\lambda_{1}\right) \in \mathcal{L} Z^{2}\left(X_{l}^{\infty}\right)$ with which we can try to link them.

Since $H_{1}\left(X_{l}^{\infty}\right)=0$, every 1-cycle in $X_{l}^{\infty}$ is a 1-boundary in $X_{l}^{\infty}$. In particular, $Z_{1}\left(W_{0}\right) \subset$ $B_{1}\left(X_{l}^{\infty}\right)$. By Lemma 8.10, the support of $\lambda_{2}-\lambda_{1}$ is disjoint from $W_{0}$, giving that $Z_{1}\left(W_{0}\right) \subset$ $B_{1}^{\lambda_{2}-\lambda_{1}}\left(X_{l}^{\infty}\right)$. Hence, we can restrict $l k\left(\cdot, \lambda_{2}-\lambda_{1}\right)$ to 1 -cycles in $W_{0}$ :

$$
l k\left(\cdot, \lambda_{2}-\lambda_{1}\right): Z_{1}\left(W_{0}\right) \rightarrow \mathbb{R}
$$

Proposition 8.13. For every $\gamma \in Z_{1}\left(W_{0}\right)$, lk $\left(\gamma, \lambda_{2}-\lambda_{1}\right)$ depends only on $[\gamma] \in H_{1}\left(W_{0}\right)$.
Proof: Suppose that $\gamma_{1}-\gamma_{2}=\partial \sigma$, with $\sigma \in C_{2}\left(W_{0}\right)$. Then, Corollary 8.11 gives that $\left\langle\sigma, \lambda_{2}-\lambda_{1}\right\rangle=$ 0 , hence $l k\left(\gamma_{1}, \lambda_{2}-\lambda_{1}\right)=l k\left(\gamma_{1}, \lambda_{2}-\lambda_{1}\right)$.

Proposition 8.14. The image of $l k\left(\cdot, \lambda_{2}-\lambda_{1}\right): H_{1}\left(W_{0}\right) \rightarrow \mathbb{R}$ is contained in the rationals $\mathbb{Q}$.
Proof: Recall from Section 3 that there is an $\epsilon>0$ for which $W_{0}$ restricted to $|x|<\epsilon$ is homeomorphic to the product $S_{0} \times \mathbb{D}_{\epsilon}$. Because any $\gamma \in Z_{1}\left(W_{0}\right)$ is compact, there exists a sufficiently high iterate $N^{k}$ so that $N^{k}(\gamma)$ lies within $|x|<\epsilon$. Then $\left[N^{k}(\gamma)\right]=n \cdot\left[S_{0}\right]$ for some appropriate $n$. Using the invariance property, this gives $l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)=\frac{n}{2^{k}} \in \mathbb{Q}$.

From here on we will write $l k\left(\cdot, \lambda_{2}-\lambda_{1}\right): H_{1}\left(W_{0}\right) \rightarrow \mathbb{Q}$.
Proposition 8.15. Suppose that $\gamma_{i}$ is a curve in a vertical line bounded by a simply connected domain $U_{i}$. Then: $\operatorname{lk}\left(\gamma_{i}, \lambda_{2}-\lambda_{1}\right)=\operatorname{size}\left(U_{i}\right)$, where $\operatorname{size}\left(U_{i}\right)$ was defined in Section 7.

Proof of Proposition 8.15: Recall that size $\left(U_{i}\right)$ is defined as $\pm \frac{l_{i}}{2^{k}}$ where $k$ is such that $N^{k}$ maps to a vertical line $\mathbb{P}_{x}$ that is divided by $W_{0}$ into only two domains $U \subset W\left(r_{1}\right)$ and $V \subset W\left(r_{2}\right)$ and where $l_{i}$ is the degree of this mapping to $U$ or $V$. The sign is - if $U_{i}$ is mapped to $U$ and + if $U_{i}$ is mapped to $V$. Without loss in generality, suppose that $U_{i}$ is mapped to $U$, and hence $\operatorname{size}\left(U_{i}\right)<0$. Using Lemma 8.9 we have that:

$$
\left\langle U_{i}, \lambda_{2}-\lambda_{1}\right\rangle=\frac{1}{2^{k}}\left\langle N^{k}\left(U_{i}\right), \lambda_{2}-\lambda_{1}\right\rangle=\frac{1}{2^{k}}\left\langle l_{i} U,-\lambda_{1}\right\rangle=-\frac{l_{i}}{2^{k}}\left\langle U_{i}, \lambda_{1}\right\rangle=-\frac{l_{i}}{2^{k}}=\operatorname{size}\left(U_{i}\right)
$$

where we are using that $\left\langle U, \lambda_{2}\right\rangle=0$ and $\left\langle U, \lambda_{1}\right\rangle=1$.
Corollary 8.16. The image of the homomorphism $l k\left(\cdot, \lambda_{2}-\lambda_{1}\right): H_{1}\left(W_{0}\right) \rightarrow \mathbb{Q}$ contains elements of arbitrarily small, but non-zero, absolute value.

This gives us our desired result:

Corollary 8.17. The homology group $H_{1}\left(W_{0}\right)$ is infinitely generated.
Notice that an additive subgroup of $\mathbb{Q}$ that is dense must be infinitely generated, but a dense additive subgroup of $\mathbb{R}$ typically is not infinitely generated because the generators can be incommensurable.

Recall the Mayer-Vietoris exact sequence (14):

$$
H_{2}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{2}\left(\overline{W\left(r_{2}\right)}\right) \rightarrow H_{2}\left(X_{l}^{\infty}\right) \xrightarrow{\partial} H_{1}\left(W_{0}\right) \rightarrow H_{1}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{2}\right)}\right) \rightarrow 0
$$

If Image $(\partial)=0$, or even if we knew that $|\operatorname{size}(\partial(\sigma))|$ were bounded away from 0 for every $\sigma \in$ $H_{2}\left(X_{l}^{\infty}\right)$, we would be able to conclude that $H_{1}\left(\overline{W\left(r_{1}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{2}\right)}\right)$ are infinitely generated. However, this is not the case.

Proposition 8.18. There are $\sigma \in H_{2}\left(X_{l}^{\infty}\right)$ with $\left|l k\left(\partial(\sigma), \lambda_{2}-\lambda_{1}\right)\right|>0$ arbitrarily small.
Proof: For every $k$, there exists some exceptional divisor $E$ having $N^{k}: E \rightarrow V$ an isomorphism. For generic parameter values $B \in S$, any exceptional divisor at a $(k-1)$-st inverse image of $p$ will have this property, since, for generic $B$ there is a single exceptional divisor above each point that we have blown up, and $N: E_{z} \rightarrow E_{N(z)}$ is always an isomorphism.

For the values of $B \notin S$, there may be many blow-ups done at each ( $k-1$ )-st inverse image of $p$. We take a detailed look at the sequence of blow-ups from Section 4.1 that was used to create $X_{l}^{k-1}$ from $X_{l}^{k-2}$. One must check that for each exceptional divisor $E_{N(z)}^{i}$ that occurs in the sequence of blow-ups at $N(z)$, there is an exceptional divisor in the sequence of blow-ups at $z$ that maps isomorphically to $E_{N(z)}^{i}$. Therefore, for any $k$, one can find an exceptional divisor $E$ so that $N^{k-1}: E \rightarrow E_{p}$ is an isomorphism. Since $N: E_{p} \rightarrow V$ is always an isomorphism, $E$ is the desired exceptional divisor.

Because $N^{k}$ maps $E$ isomorphically to $V$, it maps $\partial([E])$ to $\partial([V])$. The invariance property from Lemma 8.9 gives that

$$
l k\left(\partial([E]), \lambda_{2}-\lambda_{1}\right)=\frac{1}{2^{k}} l k\left(\partial([V]), \lambda_{2}-\lambda_{1}\right)=\frac{1}{2^{k}}
$$

Proposition 8.18.

## 8.7 $H_{1}\left(\overline{W\left(r_{1}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{2}\right)}\right)$ are infinitely generated.

The following idea will allow us to show that $H_{1}\left(\overline{W\left(r_{1}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{2}\right)}\right)$ are infinitely generated, despite the fact that $\left|l k\left(\partial(\sigma), \lambda_{2}-\lambda_{1}\right)\right|$ can be arbitrarily small, but non-zero, for $\sigma \in H_{2}\left(X_{l}^{\infty}\right)$.

## Even and odd parts of Homology:

Recall from Proposition 1.6 that $N$ has a symmetry of reflection $\tau$ about the line $B x+2 y-1=0$ which exchanges the basins of attraction. This $\tau$ induces an involution $\tau_{*}$ on $H_{*}\left(X_{l}^{\infty}\right), H_{*}\left(W_{0}\right)$, and $H_{*}\left(W\left(r_{1}\right)\right) \oplus H_{*}\left(W\left(r_{2}\right)\right)$. Every homology class $\sigma$ will have $\tau_{*}^{2}(\sigma)=\sigma$ and consequently the eigenvalues of $\tau$ are $\pm 1$.

We say that a homology class $\sigma$ is even if it is in the eigenspace of $\tau_{*}$ corresponding to eigenvalue +1 , and we say that $\sigma$ is odd if it is in the eigenspace of $\tau_{*}$ corresponding to eigenvalue -1 .

Because the Mayer-Vietoris exact sequence commutes naturally with induced maps, we have a decomposition of the sequence (14) into even and odd parts. We will only need the odd part:

$$
\left(H_{2}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{2}\left(\overline{W\left(r_{2}\right)}\right)\right)^{\text {odd }} \rightarrow H_{2}^{\text {odd }}\left(X_{l}^{\infty}\right) \xrightarrow{\partial} H_{1}^{\text {odd }}\left(W_{0}\right) \rightarrow\left(H_{1}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{2}\right)}\right)\right)^{\text {odd }} \rightarrow 0
$$

Lemma 8.19. If $\sigma$ is some piecewise smooth chain, then: $\left\langle\sigma, \lambda_{2}\right\rangle=\left\langle\tau(\sigma), \lambda_{1}\right\rangle$ and $\left\langle\sigma, \lambda_{1}\right\rangle=$ $\left\langle\tau(\sigma), \lambda_{2}\right\rangle$.

## Proof:

Recall the definition of $\lambda_{2}$ and $\lambda_{1}$ :

$$
\begin{aligned}
& \lambda_{1}=\frac{1}{2 \pi} d d^{c} \lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{2}^{k}(x, y)\right| \\
& \lambda_{2}=\frac{1}{2 \pi} d d^{c} \lim _{k \rightarrow \infty} \frac{1}{2^{k}} \log \left|N_{1}^{k}(x, y)+B \cdot N_{2}^{k}(x, y)-1\right|
\end{aligned}
$$

Since precomposition with $\tau$ exchanges the line $B x+y-1=0$ with the line $y=0$, Equation 8.19 holds.

Corollary 8.20. For every $[\gamma] \in H_{1}\left(W_{0}\right)$ we have: $l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)=-l k\left(\tau(\gamma), \lambda_{2}-\lambda_{1}\right)$.

## Proof:

Suppose that $\sigma$ is a piecewise smooth 2-chain with $\partial \sigma=\gamma$. Then we certainly have $\partial(\tau(\sigma))=\tau(\gamma)$.
Lemma 8.19 gives:

$$
l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)=\left\langle\sigma, \lambda_{2}-\lambda_{1}\right\rangle=\left\langle\tau(\sigma), \lambda_{1}-\lambda_{2}\right\rangle=-\left\langle\tau(\sigma), \lambda_{2}-\lambda_{1}\right\rangle=-l k\left(\tau(\gamma), \lambda_{2}-\lambda_{1}\right)
$$

Proposition 8.21. If $\gamma \in H_{1}^{\text {odd }}\left(W_{0}\right)$ is in the image of the boundary map $\partial: H_{2}^{\text {odd }}\left(X_{l}^{\infty}\right) \rightarrow$ $H_{1}^{\text {odd }}\left(W_{0}\right)$, then $l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)=0$.

We first need the following lemma:
Lemma 8.22. For any exceptional divisor $E_{z}$ we have

$$
\begin{equation*}
\partial\left(\tau_{*}\left(\left[E_{z}\right]\right)\right)=-\tau_{*}\left(\partial\left(\left[E_{z}\right]\right)\right) \tag{21}
\end{equation*}
$$

Proof: This proof will depend essentially on the explicit interpretation of the boundary map $\partial$ from the Mayer-Vietoris sequence. In the following paragraph we closely paraphrase Hatcher [31], p. 150:
The boundary map $\partial: H_{n}(X) \rightarrow H_{n-1}(A \cap B)$ can be made explicit. A class $\alpha \in H_{n}(X)$ is represented by a cycle $z$. By appropriate subdivision, we can write $z$ as a sum $x+y$ of chains in $A$ and $B$, respectively. While it need not be true that $x$ and $y$ are cycles individually, we do have $\partial x=-\partial y$ since $z=x+y$ is a cycle. The element $\partial \alpha$ is represented by the cycle $\partial x=-\partial y$.


Figure 12: Showing that $\partial\left(\tau_{*}\left(\left[E_{z}\right]\right)\right)=-\tau_{*}\left(\partial\left(\left[E_{z}\right]\right)\right)$.

We use this explicit interpretation of $\partial$ to check Equation 21. Notice that $\tau_{*}\left(\left[E_{z}\right]\right)=\left[E_{\tau(z)}\right]$ consistent with the orientation that $E_{z}$ and $E_{\tau(z)}$ have as Riemann surfaces. Therefore we have that $\partial\left(\tau_{*}\left(\left[E_{z}\right]\right)\right)=\partial\left(\left[E_{\tau(z)}\right]\right)=\left[\partial U_{1}\right]=-\left[\partial U_{2}\right]$, where $U_{1}$ is the oriented region of $E_{\tau(z)}$ that is in $\overline{W\left(r_{1}\right)}$ and $U_{2}$ is the oriented region of $E_{\tau(z)}$ that is in $\overline{W\left(r_{2}\right)}$.

Similarly $\partial\left(\left[E_{z}\right]\right)=\left[\partial V_{1}\right]=-\left[\partial V_{2}\right]$, where $V_{1}$ and $V_{2}$ are $E_{z} \cap \overline{W\left(r_{1}\right)}$ and $E_{z} \cap \overline{W\left(r_{2}\right)}$. Because $\tau$ maps $E_{z}$ to $E_{\tau(z)}$ swapping $\overline{W\left(r_{1}\right)}$ with $\overline{W\left(r_{2}\right)}$ we have:

$$
\tau_{*}\left(\partial\left(\left[E_{z}\right]\right)\right)=\left[\partial U_{2}\right]=-\partial\left(\tau_{*}\left(\left[E_{z}\right]\right)\right)
$$

## Proof of Proposition 8.21:

Since elements of the form $\left[E_{z}\right]-\left[\tau\left(E_{z}\right)\right]$ span $H_{2}^{\text {odd }}\left(X_{l}^{\infty}\right)$, we need only check that the images of differences like this under $\partial$ have 0 linking number:

$$
\begin{aligned}
l k\left(\partial\left(\left[E_{z}\right]-\left[\tau\left(E_{z}\right)\right]\right), \lambda_{2}-\lambda_{1}\right) & =\operatorname{lk}\left(\partial\left(\left[E_{z}\right]\right)-\partial\left(\tau_{*}\left(\left[E_{z}\right]\right)\right), \lambda_{2}-\lambda_{1}\right) \\
& =\operatorname{lk}\left(\partial\left(\left[E_{z}\right]\right)+\tau_{*}\left(\partial\left(\left[E_{z}\right]\right)\right), \lambda_{2}-\lambda_{1}\right)=0
\end{aligned}
$$

The last term is 0 by Lemma 8.22.
Proposition 8.23. The image of $l k\left(\cdot, \lambda_{2}-\lambda_{1}\right): H_{1}^{\text {odd }}\left(W_{0}\right) \rightarrow \mathbb{Q}$ contains elements of arbitrarily small, but non-zero absolute value.

## Proof of Proposition 8.23:

Recall from Proposition 8.16 that we can find 1-cycles $\gamma$ that have $l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)$ arbitrarily small, but non-zero. Notice that $[\gamma-\tau(\gamma)]$ is obviously odd, and using Lemma 8.22:

$$
\begin{aligned}
l k\left(\gamma-\tau(\gamma), \lambda_{2}-\lambda_{1}\right) & =l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)-l k\left(\tau(\gamma), \lambda_{2}-\lambda_{1}\right) \\
& =l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)+l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)=2 l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)
\end{aligned}
$$

Hence, by choosing $\gamma$ so that $l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)$ is arbitrarily small, but non-zero, we can make $l k(\gamma-$ $\left.\tau(\gamma), \lambda_{2}-\lambda_{1}\right)$ arbitrarily small, but non-zero with $[\gamma-\tau(\gamma)] \in H_{1}^{\text {odd }}\left(W_{0}\right)$.

Recall the last part of the exact sequence on the odd parts of homology:

$$
\rightarrow H_{2}^{\text {odd }}\left(X_{l}^{\infty}\right) \xrightarrow{\partial} H_{1}^{\text {odd }}\left(W_{0}\right) \xrightarrow{i_{1 *} \oplus i_{2 *}}\left(H_{1}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{2}\right)}\right)\right)^{\text {odd }} \rightarrow 0
$$

where $i_{1}$ and $i_{2}$ are the inclusions $W_{0} \hookrightarrow \overline{W\left(r_{1}\right)}$ and $W_{0} \hookrightarrow \overline{W\left(r_{2}\right)}$ respectively.
As a consequence of Proposition 8.21, given any $\eta \in\left(H_{1}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{2}\right)}\right)\right)^{\text {odd }}$ we can define $l k\left(\eta, \lambda_{2}-\lambda_{1}\right)=l k\left(\gamma, \lambda_{2}-\lambda_{1}\right)$ for any $\gamma \in H_{1}^{\text {odd }}\left(W_{0}\right)$ whose image under $i_{1 *} \oplus i_{2 *}$ is $\eta$. As a consequence of Proposition 8.23 we know that there are $\eta \in\left(H_{1}\left(\overline{W\left(r_{1}\right)}\right) \oplus H_{1}\left(\overline{W\left(r_{2}\right)}\right)\right)^{\text {odd }}$ with arbitrarily small $\left|l k\left(\eta, \lambda_{2}-\lambda_{1}\right)\right|$. This proves the desired result:

Theorem 8.24. Let $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ be the closures in $X_{l}^{\infty}$ of the basins of attraction of the roots $r_{1}=(0,0)$ and $r_{2}=(0,1)$ under the Newton Map $N$. Then $H_{1}\left(\overline{W\left(r_{1}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{2}\right)}\right)$ are infinitely generated.

Recall also:
Corollary 8.25. For parameter values $B \in \Omega_{r}$, we can replace $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ with $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ finding that $H_{1}\left(W\left(r_{1}\right)\right)$ and $H_{1}\left(W\left(r_{2}\right)\right)$ are also infinitely generated.

### 8.8 Linking with currents in $X_{r}$

Much of the work in the previous few subsections was to make linking numbers well defined in $X_{l}^{\infty}$, overcoming the indeterminacy from the fact that $H_{2}\left(X_{l}^{\infty}\right)$ is infinitely generated. Because $H_{2}\left(X_{r}\right) \cong \mathbb{Z}^{\{[\mathbb{P}]\}}$ it is relatively easy to find elements in $\mathcal{L} Z_{2}\left(X_{r}\right)$. However, we can just mimic the work from the previous sub-sections in an appropriate way.

The major difference is that in $X_{l}^{\infty}$ there is always an intersection of $W_{0}$ with $C$ resulting in loops in $W_{0}$ of arbitrarily small size. In $X_{r}$, we must stipulate that an intersection of $W_{1}$ with $C$ exists before proving that the homology is infinitely generated, because there appear to be parameter values for which there is no intersection.

If we define $\lambda_{3}$ and $\lambda_{4}$ in a similar way as we defined $\lambda_{1}$ and $\lambda_{2}$, then the following are proven in an easy way:

Proposition 8.26. If $W_{1}$ intersects the critical value parabola $C$, then $H_{1}\left(W_{1}\right)$ is infinitely generated.

Since there is only one generator of $H_{2}\left(X_{r}\right)$ this directly gives:
Theorem 8.27. If $W_{1}$ intersects the critical value parabola $C$, then $H_{1}\left(\overline{W\left(r_{3}\right)}\right)$ and $H_{1}\left(\overline{W\left(r_{4}\right)}\right)$ are infinitely generated.
where $\overline{W\left(r_{3}\right)}$ and $\overline{W\left(r_{4}\right)}$ are the closures in $X_{r}$ of the basins of attraction of roots $r_{3}=(1,0)$ and $r_{4}=(1,1-B)$ under $N$.
Corollary 8.28. For parameter values $B \in \Omega_{r e g}$, we can replace $\overline{W\left(r_{1}\right)}$ and $\overline{W\left(r_{2}\right)}$ with $W\left(r_{1}\right)$ and $W\left(r_{2}\right)$ finding that $H_{1}\left(W\left(r_{1}\right)\right)$ and $H_{1}\left(W\left(r_{2}\right)\right)$ are also infinitely generated.

This is the last part of the proof of Theorem 0.1.

## A Blow-ups of complex surfaces at a point.

Further material is available in [28, pp. 182-189 and 473-478] and the introduction of [33].
Suppose that $R: \mathbb{C}^{2} \rightarrow \mathbb{C}^{2}$ has a point of indeterminacy at $(0,0)$. Blowing-up at $(0,0)$ produces a new space:

$$
\begin{equation*}
\widetilde{\mathbb{C}}_{(0,0)}^{2}=\left\{(z, l) \in \mathbb{C}^{2} \times \mathbb{P}^{1}: z \in l\right\} \tag{22}
\end{equation*}
$$

to which we can often find an extension $R: \widetilde{\mathbb{C}}_{(0,0)}^{2} \rightarrow \mathbb{C}^{2}$ with no indeterminacy. Here $\mathbb{P}^{1}$ is identified with the space of directions through $(0,0)$ in $\mathbb{C}^{2}$. There is a natural projection $\rho: \widetilde{\mathbb{C}}_{(0,0)}^{2} \rightarrow \mathbb{C}^{2}$ given by $\rho(z, l)=z$ and the set $E_{(0,0)}=\rho^{-1}((0,0))$ is referred to as the exceptional divisor. A standard check shows that the blow-up is independent of the choice of coordinates hence well defined on a complex surface $M$ at a point $z$.

A rational map $R: \mathbb{C}^{2} \rightarrow \mathbb{C}^{2}$ can be lifted to a new rational mapping $\widetilde{R}: \widetilde{\mathbb{C}}_{(0,0)}^{2}-E_{(0,0)} \rightarrow \mathbb{C}^{2}$ be defining $R(z, l)=R(z)$ for $z \neq 0$. If the indeterminacy in $R$ at $(0,0)$ was reasonably tame, $\widetilde{R}$ extends by continuity to all of $E_{(0,0)}$. Otherwise, there will be points of indeterminacy of $\widetilde{R}$ on $E_{(0,0)}$ to which $\widetilde{R}$ cannot be extended. One can try further blow-ups at these points to resolve these new points of indeterminacy. The extension of $\widetilde{R}$ to $E_{(0,0)}$ is analytic except at any new points of indeterminacy because $E_{(0,0)}$ is a space of complex co-dimension 1 .

Proposition A.1. If $M$ is a complex surface and $z$ is any point in $M$, then the blow-up $\widetilde{M}_{z}$ has $H_{2}\left(\widetilde{M}_{z}\right) \cong H_{2}(M) \oplus \mathbb{Z}^{\left(\left\{\left[E_{z}\right]\right\}\right)}$ and $H_{i}\left(\widetilde{M}_{z}\right) \cong H_{i}(M)$ for $i \neq 2$.

The proof is an application of the Mayer-Vietoris sequence on homology and the fact that $\widetilde{\mathbb{C}}_{(0,0)}^{2}$ has the homotopy type of $\mathbb{P}$.

Further analysis shows that the fundamental class $\left[E_{z}\right]$ of an exceptional divisor has self-intersection number -1 . Meanwhile, blowing up a smooth point on any complex curve $C$ decreases the selfintersection of its homology class [C] by one. (See [28], for proof.)

## B Proof of Theorem 4.1

Let $S \subset \Omega$ be the set of parameter values $B$ for which no inverse image of the point of indeterminacy $p$ or the point of indeterminacy $q$ is in the critical value locus $C$. We are especially interested in $B \in S$ because the sequence of blow-ups from Section 4.1 is especially easy to describe for these $B$.

Theorem 4.1 states that $S$ is generic in the sense of Baire's Theorem, i.e. uncountable and dense in $\Omega$. The proof will follow as a corollary to:

Theorem. (Baire) Let $X$ be either a complete metric space, or a locally compact Hausdorf space. Then, the intersection of any countable family of dense open sets in $X$ is dense.

See Bredon [13] for a proof of Baire's Theorem.
Proof of Theorem 4.1: Let $S_{n} \subset \mathbb{C}$ be the subset of parameter values $B$ for which none of the $n$-th inverse images of $p$ or $q$ under $N$ are in the critical value locus $C$.

Lemma B.1. $S_{n}$ is a dense open set in $\mathbb{C}$
Proof: Let $R_{n}$ be the set of $B$ for which an $n$-th inverse image of $p$ is in $C$ and let $T_{n}$ be the set of $B$ for which an $n$-th inverse image of $q$ in $C$. We will show that $R_{n}$ and $T_{n}$ are finite, so that $S_{n}=\Omega-\left(R_{n} \cup T_{n}\right)$ is a dense open set.

Lemma B.2. $T_{n}$ is a finite set.
Proof: We have $B \in T_{n}$ if:

$$
\begin{equation*}
y^{2}+B x y+\frac{B^{2}}{4} x^{2}-\frac{B^{2}}{4} x-y=0, \quad N_{1}^{n}(x, y)=\frac{1}{2-B}, \quad N_{2}^{n}(x, y)=\frac{1-B}{2-B} \tag{23}
\end{equation*}
$$

has a solution. As always, $N_{1}^{n}$ and $N_{2}^{n}$ denote the first and second coordinates of $N^{n}$. By clearing the denominators in the second and third equations, condition 23 can be expressed as the common zeros of 3 polynomials $P_{1}(x, y, B), P_{2}(x, y, B)$, and $P_{3}(x, y, B)$ in the three variables $x$, $y$, and $B$. We will check there is no common divisor of $P_{1}(x, y, B), P_{2}(x, y, B)$, and $P_{3}(x, y, B)$ so that the solutions to 23 form a finite set.

First, notice that $P_{1}(x, y, B)=y^{2}+B x y+\frac{B^{2}}{4} x^{2}-\frac{B^{2}}{4} x-y$ is irreducible. It is sufficient to write an explicit biholomorphic map from $\mathbb{C}^{2}$ to $\left\{P_{1}=0\right\} \subset \mathbb{C}^{3}$. At a given $B$, the line $B x+2 y=t$ intersects $\left\{P_{1}=0\right\}$ at a single point which we denote by $f_{B}(t)$. It is easy to check that $(t, B) \mapsto\left(f_{B}(t), B\right)$ provides the desired isomorphism.

Hence $P_{1}$ has a factor in common with $P_{2}$ or $P_{3}$ if and only if $P_{1}$ divides $P_{2}$ or $P_{3}$. We will show that this is impossible by examining the lowest degree terms of $P_{2}$ and $P_{3}$. If $P_{1}$ divides $P_{2}$ or $P_{3}$, then the lowest degree term, $-y$, of $P_{1}$ must divide the lowest degree term of $P_{2}$ or the lowest degree term of $P_{3}$.

We check by induction that the lowest degree term of $P_{2}$ is $\pm 1$ for every $n$. To simplify notation, let $a_{k}(x, y, B)$ be the polynomial obtained by clearing the denominators from $N_{1}^{k}(x, y)=\frac{1}{2-B}$.

By clearing denominators of $N_{1}(x, y)=\frac{1}{2-B}$, we find $a_{1}(x, y, B)=x^{2}(2-B)-1(2 x-1)=2 x^{2}-$ $B x^{2}-2 x+1$, so $a_{1}(x, y, B)$ has constant term $\pm 1$. Now suppose that $a_{n}(x, y, B)$ has constant term $\pm 1$. By definition, $a_{n+1}(x, y, B)$ is obtained by clearing the denominators of $a_{n}\left(N_{1}(x, y), N_{2}(x, y), B\right)=$
0. Because the denominators of both $N_{1}(x, y)$ and $N_{2}(x, y)$ have constant term $\pm 1$ and because $a_{n}(x, y, B)$ has constant term 1 we find that $a_{n+1}(x, y, B)$ has constant term $\pm 1$.

Because $P_{2}$ has constant term $\pm 1$ for every $n P_{1}$ cannot divide $P_{2}$, and we conclude that there are no common factors between $P_{1}$ and $P_{2}$.

A nearly identical proof by induction shows that lowest degree term of $P_{3}$ is also $\pm 1$ for each $n$. Hence $P_{1}$ does not divide $P_{3}$, and we conclude that $P_{1}$ and $P_{3}$ have no common divisors.

To see that $P_{2}$ and $P_{3}$ have no common divisors, notice that $P_{2}(x, y, B)=0$ is an equation for many disjoint vertical lines, while $P_{3}(x, y, B)=0$ stipulates that the $n$-th image of this locus has constant $y=0$. Since vertical lines are mapped to vertical lines by $N, P_{2}$ and $P_{3}$ can have no common factors.

Hence, $P_{1}, P_{2}$, and $P_{3}$ are algebraically independent, so they have a finite number of common zeros, giving that $T_{n}$ is a finite set. $\square$ Lemma B.2.

Lemma B.3. $R_{n}$ is a finite set.
Proof: Now we show that $R_{n}$, the set of $B$ so that an $n$-th inverse image of $p$ under $N$ is in $C$, is finite. In terms of equations, $R_{n}$ is the set of $B$ so that:

$$
\begin{equation*}
y^{2}+B x y+\frac{B^{2}}{4} x^{2}-\frac{B^{2}}{4} x-y=0, \quad N_{1}^{n}(x, y)=\frac{1}{B}, \quad N_{2}^{n}(x, y)=0 \tag{24}
\end{equation*}
$$

has a solution. Let $Q_{1}, Q_{2}$, and $Q_{3}$ be the polynomials equations resulting from clearing the denominators in Equation 24.

The proof is the same as for $T_{n}$ except that a different proof is needed to see that $Q_{1}$ does not divide $Q_{3}$. An adaptation of the proof that $P_{1}$ does not divide $P_{3}$ fails because the lowest degree term of $Q_{3}$ has positive degree in $y$. We will check that $Q_{1}$ does not divide $Q_{3}$ and leave the remainder of the proof to the reader.

The $x$-axis, $y=0$, is one of the invariant lines of $N$ and it intersects the basins $W\left(r_{1}\right), W\left(r_{3}\right)$ and the separator $\operatorname{Re}(x)=1 / 2$. Therefore it is disjoint from the two basins $W\left(r_{2}\right)$ and $W\left(r_{4}\right)$. By definition, $Q_{3}(x, y, B)$ is the equation for the $n$-the inverse image of the $x$-axis. So, for a given $B$, the locus $Q_{3}(x, y, B)=0$ is also disjoint from the two basins $W\left(r_{2}\right)$ and $W\left(r_{4}\right)$.

For every $B$, the critical value parabola $C$ goes through the four roots $r_{1}, r_{2}, r_{3}$, and $r_{4}$, so it intersects all four basins of attraction. By definition, $C$ is the zero locus $Q_{1}(x, y, B)=0$. Therefore, if $Q_{1}$ divides $Q_{3}$, there is a component of the zero locus $Q_{3}(x, y, B)=0$ intersecting all four basins $W\left(r_{1}\right), W\left(r_{2}\right), W\left(r_{3}\right)$ and $W\left(r_{4}\right)$ for every $B$. This is impossible, so $Q_{1}$ cannot divide $Q_{3}$. $\square$ Lemma B. 3 and $\square$ Lemma B.1.

Since $S_{n}$ is a dense open set in $\Omega$ for each $n$ and $S=\cap_{n=0}^{\infty} S_{n}$, so it follows from Baire's Theorem that $S$ is uncountable and dense in the parameter space $\Omega$. $\square$ Theorem 4.1.

## Aknowledgements

This paper is a condensed version of my Ph.D. thesis for Cornell University. I thank the Department of Defense for generous financial support by means of a National Defense Sciences and Engineering Graduate (NDSEG) fellowship and I thank the National Science Foundation for further financial support by an Interdisciplinary Graduate Education and Research Traineeship (IGERT) fellowship.

My advisor John H. Hubbard suggested that I study the topology of the basins of attraction for Newton's Method. He provided mathematical guidance and enthusiasm about this work and introduced to me the use of blow-ups to resolve points of indeterminacy. John Smillie and Eric Bedford provided a number of suggestions including encouraging me to learn about currents, which became a key technique in this paper. Allen Hatcher provided helpful comments about the topology. Alexey Glutsyuk provided very helpful discussions and suggestions during the final stages of this project.

The computer program FractalAsm [44] written by Karl Papadantonakis was used to generate all of the images of the basins of attraction and was invaluable for gaining an intuition about the topology of these basins.

The referee caught an error in my pairing between closed currents and piecewise smooth chains (allowing for a correction) and also provided many valuable suggestions for streamlining and clarifying the final manuscript. I thank the referee greatly for his or her detailed consideration.

## References

[1] L. M. Abramov and V. A. Rohlin. Entropy of a skew product of mappings with invariant measure. Vestnik Leningrad. Univ., 17(7):5-13, 1962.
[2] Eric Bedford, Mikhail Lyubich, and John Smillie. Polynomial diffeomorphisms of $\mathbb{C}^{2}$. IV. The measure of maximal entropy and laminar currents. Invent. Math., 112(1):77-125, 1993.
[3] Eric Bedford and John Smillie. Polynomial diffeomorphisms of $\mathbf{C}^{2}$ : currents, equilibrium measure and hyperbolicity. Invent. Math., 103(1):69-99, 1991.
[4] Eric Bedford and John Smillie. Polynomial diffeomorphisms of C ${ }^{2}$. II. Stable manifolds and recurrence. J. Amer. Math. Soc., 4(4):657-679, 1991.
[5] Eric Bedford and John Smillie. Polynomial diffeomorphisms of $\mathbb{C}^{2}$. III. Ergodicity, exponents and entropy of the equilibrium measure. Math. Ann., 294(3):395-420, 1992.
[6] Eric Bedford and John Smillie. Polynomial diffeomorphisms of $\mathbb{C}^{2}$. V. Critical points and Lyapunov exponents. J. Geom. Anal., 8(3):349-383, 1998.
[7] Eric Bedford and John Smillie. Polynomial diffeomorphisms of $\mathbb{C}^{2}$. VI. Connectivity of J. Ann. of Math. (2), 148(2):695-735, 1998.
[8] Eric Bedford and John Smillie. External rays in the dynamics of polynomial automorphisms of $\mathbb{C}^{2}$. In Complex geometric analysis in Pohang (1997), volume 222 of Contemp. Math., pages 41-79. Amer. Math. Soc., Providence, RI, 1999.
[9] Eric Bedford and John Smillie. Polynomial diffeomorphisms of $\mathbb{C}^{2}$. VII. Hyperbolicity and external rays. Ann. Sci. École Norm. Sup. (4), 32(4):455-497, 1999.
[10] Araceli M. Bonifant and Marius Dabija. Self-maps of $\mathbb{P}^{2}$ with invariant elliptic curves. In Complex manifolds and hyperbolic geometry (Guanajuato, 2001), volume 311 of Contemp. Math., pages 1-25. Amer. Math. Soc., Providence, RI, 2002.
[11] Araceli M. Bonifant and John Erik Fornæss. Growth of degree for iterates of rational maps in several variables. Indiana Univ. Math. J., 49(2):751-778, 2000.
[12] Raoul Bott and Loring W. Tu. Differential forms in algebraic topology, volume 82 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1982.
[13] Glen E. Bredon. Topology and geometry, volume 139 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1993.
[14] Jean-Yves Briend. Exposants de Liapounof et points périodiques d'endomorphismes holomorphes de $\mathbb{C P}^{k}$. PhD thesis, Toulouse, 1997.
[15] Jean-Yves Briend and Julien Duval. Deux caractérisations de la mesure d'équilibre d'un endomorphisme de $\mathrm{P}^{k}(\mathbf{C})$. Publ. Math. Inst. Hautes Études Sci., (93):145-159, 2001.
[16] Hans Brolin. Invariant sets under iteration of rational functions. Ark. Mat., 6:103-144 (1965), 1965.
[17] Jean-Pierre Demailly. Courants positifs et théorie de l'intersection. Gaz. Math., (53):131-159, 1992.
[18] R. Devaney and Z. Nitecki. Shift automorphisms in the Hénon mapping. Comm. Math. Phys., 67(2):137-146, 1979.
[19] Jeffrey Diller. Dynamics of birational maps of $\mathbf{P}^{2}$. Indiana Univ. Math. J., 45(3):721-772, 1996.
[20] Tien-Cuong Dinh and Nessim Sibony. Dynamique des applications d'allure polynomiale. J. Math. Pures Appl. (9), 82(4):367-423, 2003.
[21] Tien-Cuong Dinh and Nessim Sibony. Dynamics of regular birational maps in $\mathbb{P}^{k}$. J. Funct. Anal., 222(1):202-216, 2005.
[22] Romain Dujardin. Hénon-like mappings in $\mathbb{C}^{2}$. Amer. J. Math., 126(2):439-472, 2004.
[23] C. Favre and M. Jonsson. Brolin's theorem for curves in two complex dimensions. Ann. Inst. Fourier (Grenoble), 53(5):1461-1501, 2003.
[24] John Erik Fornæss. Dynamics in several complex variables, volume 87 of CBMS Regional Conference Series in Mathematics. Published for the Conference Board of the Mathematical Sciences, Washington, DC, 1996.
[25] John Erik Fornæss and Nessim Sibony. Complex dynamics in higher dimension. I. Astérisque, (222):5, 201-231, 1994. Complex analytic methods in dynamical systems (Rio de Janeiro, 1992).
[26] John Erik Fornæss and Nessim Sibony. Complex dynamics in higher dimensions. In Complex potential theory (Montreal, PQ, 1993), volume 439 of NATO Adv. Sci. Inst. Ser. C Math. Phys. Sci., pages 131-186. Kluwer Acad. Publ., Dordrecht, 1994. Notes partially written by Estela A. Gavosto.
[27] John Erik Fornaess and Nessim Sibony. Complex dynamics in higher dimension. II. In Modern methods in complex analysis (Princeton, NJ, 1992), volume 137 of Ann. of Math. Stud., pages 135-182. Princeton Univ. Press, Princeton, NJ, 1995.
[28] Phillip Griffiths and Joseph Harris. Principles of algebraic geometry. Wiley-Interscience [John Wiley \& Sons], New York, 1978. Pure and Applied Mathematics.
[29] Vincent Guedj. Dynamics of quadratic polynomial mappings of $\mathbb{C}^{2}$. Michigan Math. J., 52(3):627-648, 2004.
[30] Vincent Guedj. Ergodic properties of rational mappings with large topological degree. Annals of Mathematics, 161:1589-1607, 2005.
[31] Allen Hatcher. Algebraic topology. Cambridge University Press, Cambridge, 2002.
[32] Stefan-M. Heinemann. Julia sets of skew products in C ${ }^{2}$. Kyushu J. Math., 52(2):299-329, 1998.
[33] John Hubbard, Peter Papadopol, and Vladimir Veselov. A compactification of Hénon mappings in $\mathbf{C}^{2}$ as dynamical systems. Acta Math., 184(2):203-270, 2000.
[34] John H. Hubbard and Ralph W. Oberste-Vorth. Hénon mappings in the complex domain. I. The global topology of dynamical space. Inst. Hautes Études Sci. Publ. Math., (79):5-46, 1994.
[35] John H. Hubbard and Ralph W. Oberste-Vorth. Hénon mappings in the complex domain. II. Projective and inductive limits of polynomials. In Real and complex dynamical systems (Hillerød, 1993), volume 464 of NATO Adv. Sci. Inst. Ser. C Math. Phys. Sci., pages 89-132. Kluwer Acad. Publ., Dordrecht, 1995.
[36] John H. Hubbard and Ralph W. Oberste-Vorth. Linked solenoid mappings and the nontransversality locus invariant. Indiana Univ. Math. J., 50(1):553-566, 2001.
[37] John H. Hubbard and Peter Papadopol. Newton's method applied to two quadratic equations in $\mathbb{C}^{2}$ viewed as a global dynamical system. To appear in the Memoires of the AMS.
[38] John H. Hubbard and Peter Papadopol. Superattractive fixed points in $\mathbf{C}^{n}$. Indiana Univ. Math. J., 43(1):321-365, 1994.
[39] Mattias Jonsson. Dynamics of polynomial skew products on C ${ }^{2}$. Math. Ann., 314(3):403-447, 1999.
[40] L. V. Kantorovič. On Newton's method. Trudy Mat. Inst. Steklov., 28:104-144, 1949.
[41] R. Mañé, P. Sad, and D. Sullivan. On the dynamics of rational maps. Ann. Sci. École Norm. Sup. (4), 16(2):193-217, 1983.
[42] John Milnor. Dynamics in one complex variable. Friedr. Vieweg \& Sohn, Braunschweig, 1999. Introductory lectures.
[43] John W. Milnor. Topology from the differentiable viewpoint. Based on notes by David W. Weaver. The University Press of Virginia, Charlottesville, Va., 1965.
[44] Karl Papadantonakis. Fractalasm. A fast, generalized fractal exploration program. http://www.math.cornell.edu/~dynamics/FA/index.html.
[45] Dale Rolfsen. Knots and links, volume 7 of Mathematics Lecture Series. Publish or Perish Inc., Houston, TX, 1990. Corrected reprint of the 1976 original.
[46] Alexander Russakovskii and Bernard Shiffman. Value distribution for sequences of rational mappings and complex dynamics. Indiana Univ. Math. J., 46(3):897-932, 1997.
[47] Olivier Sester. Hyperbolicité des polynômes fibrés. Bull. Soc. Math. France, 127(3):393-428, 1999.
[48] Nessim Sibony. Dynamique des applications rationnelles de $\mathbf{P}^{k}$. In Dynamique et géométrie complexes (Lyon, 1997), volume 8 of Panor. Synthèses, pages ix-x, xi-xii, 97-185. Soc. Math. France, Paris, 1999.
[49] John Smillie. Complex dynamics in several variables. In Flavors of geometry, volume 31 of Math. Sci. Res. Inst. Publ., pages 117-150. Cambridge Univ. Press, Cambridge, 1997. With notes by Gregery T. Buzzard.
[50] Hiroki Sumi. A correction to the proof of a lemma in: "Dynamics of sub-hyperbolic and semi-hyperbolic rational semigroups and skew products". Ergodic Theory Dynam. Systems, 21(4):1275-1276, 2001.
[51] Hiroki Sumi. Dynamics of sub-hyperbolic and semi-hyperbolic rational semigroups and skew products. Ergodic Theory Dynam. Systems, 21(2):563-603, 2001.
[52] Tetsuo Ueda. Fatou sets in complex dynamics on projective spaces. J. Math. Soc. Japan, 46(3):545-555, 1994.

