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## To cite this version:

Yan Gao, Ya Fei Ou. The dynatomic curves for unimodel polynomials are smooth and irreducible. 2012. hal-00814484

HAL Id: hal-00814484

## https://hal.science/hal-00814484

Preprint submitted on 17 Apr 2013

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# The dynatomic periodic curves for polynomial $\mathbf{z} \mapsto \mathbf{z}^{\mathbf{d}}+\mathbf{c}$ are smooth and irreducible 

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April 17, 2013


#### Abstract

We prove here the smoothness and the irreducibility of the periodic dynatomic curves $(c, z) \in \mathbb{C}^{2}$ such that $z$ is $n$-periodic for $z^{d}+c$, where $d \geq 2$.

We use the method provided by Xavier Buff and Tan Lei in $[\mathrm{BT}]$ where they prove the conclusion for $d=2$. The proof for smoothness is based on elementary calculations on the pushforwards of specific quadratic differentials, following Thurston and Epstein, while the proof for irreducibility is a simplified version of Lau-Schleicher's proof by using elementary arithmetic properties of kneading sequence instead of internal addresses.


## 1 Introduction

For $c \in \mathbb{C}$, set $f_{c}(z)=z^{d}+c$, where $d \geq 2$. For $n \geq 1$, define

$$
X_{n}:=\left\{(c, z) \in \mathbb{C}^{2} \mid f_{c}^{n}(z)=z,\left(f_{c}^{n}\right)^{\prime}(z) \neq 1 \text { and for all } 0<m<n, \quad f_{c}^{m}(z) \neq z\right\} .
$$

The objective of this note is to give an elementary proof of the following results:
Theorem 1.1. For every $n \geq 1$, the closure of $X_{n}$ in $\mathbb{C}^{2}$ is smooth.
Theorem 1.2. For every $n \geq 1$ the closure of $X_{n}$ in $\mathbb{C}^{2}$ is irreducible.

The first example is, as $d=2$

$$
X_{1}=\left\{(c, z) \in \mathbb{C}^{2} \mid z^{2}+c=z\right\} \backslash\left\{\left(\frac{1}{4}, \frac{1}{2}\right)\right\}=\left\{(c, z) \in \mathbb{C}^{2} \mid c=z-z^{2}\right\} \backslash\left\{\left(\frac{1}{4}, \frac{1}{2}\right)\right\}
$$

[^0]and
$$
\overline{X_{1}}=\left\{(c, z) \in \mathbb{C}^{2} \mid c=z-z^{2}\right\} .
$$

In the case $d=2$, Theorem 1.1 was proved by Douady-Hubbard and Buff-Tan in different methods; Theorem 1.2 was proved by Bousch, Morton, Lau-Schleicher and BuffTan with different approaches.

Our approach here to the two Theorems is a generalisation to that used by Xavier Buff and Tan Lei in [BT], where they prove the conclusion for $d=2$. To prove Theorem 1.1, we use elementary calculations on quadratic differentials and Thurston's contraction principle. To prove Theorem 1.2, we use a dynamical method by a purely arithmetic argument on kneading sequences(Lemma 3.2 below).

Section 2 proves the smoothness and Section 3 proves the irreducibility. The two sections can be read independently.

Acknowlegement. We thank Tan Lei for helpful discussions, and Yang Fei for providing some good pictures

## 2 Smoothness of the periodic curves

For $n \geq 1$, and $(c, z) \in \mathbb{C}^{2}$, we say that $z$ is periodic of period $n$ for $f_{c}: z \mapsto z^{d}+c$, where $d \geq 2$, if $f_{c}^{\circ n}(z)=z$ and for all $0<m<n, f_{c}^{m}(z) \neq z$. In this case the multiplier of $z$ for $f_{c}$ is defined to be $\left[f_{c}^{n}\right]^{\prime}(z)$.

We define

$$
X_{n}:=\left\{(c, z) \in \mathbb{C}^{2} \mid z \text { is of period } n \text { for } f_{c} \text { and of multiplier distinct from } 1\right\}
$$

The objective of here is to give an elementary proof of the following result:
Theorem 2.1. For every $n \geq 1$, the closure $\overline{X_{n}}$ of $X_{n}$ in $\mathbb{C}^{2}$ is smooth. More precisely, the boundary $\partial X_{n}$ is the finite set of $(c, z) \in \mathbb{C}^{2}$ such that $z$ is of period $m \leq n$ dividing $n$ for $f_{c}$ whose multiplier is of the form $e^{2 \pi i u / v}$ with $u, v \geq 1$ co-prime and $v=n / m$. In a neighborhood of a point $\left(c_{0}, z_{0}\right) \in \overline{X_{n}}$, the set $\overline{X_{n}}$ is locally the graph of a holomorphic $\operatorname{map}\left\{\begin{array}{ll}c \mapsto z(c) \text { with } z\left(c_{0}\right)=z_{0} & \text { if }\left(c_{0}, z_{0}\right) \in X_{n} \\ z \mapsto c(z) \text { with } c\left(z_{0}\right)=c_{0} \text { and } c^{\prime}\left(z_{0}\right)=0 & \text { if }\left(c_{0}, z_{0}\right) \in \partial X_{n}\end{array}\right.$.

The idea is to prove that some partial derivative of some defining function of $\overline{X_{n}}$ is non vanishing. Following A. Epstein, we will express this derivative as the coefficient of a quadratic differential of the form $\left(f_{c}\right)_{*} \mathcal{Q}-\mathcal{Q}$. Thurston's contraction principle gives $\left(f_{c}\right)_{*} \mathcal{Q}-\mathcal{Q} \neq 0$, therefore the non-nullness of our partial derivative.

### 2.1 Quadratic differentials and contraction principle

A meromorphic quadratic differential (or in short, a quadratic differential) $\mathcal{Q}$ on $\mathbb{C}$ takes the form $\mathcal{Q}=q \mathrm{dz}^{2}$ with $q$ a meromorphic function on $\mathbb{C}$.

We use $\mathcal{Q}(\mathbb{C})$ to denote the set of meromorphic quadratic differentials on $\mathbb{C}$ whose poles (if any) are all simple. If $\mathcal{Q} \in \mathcal{Q}(\mathbb{C})$ and $U$ is a bounded open subset of $\mathbb{C}$, the norm

$$
\|\mathcal{Q}\|_{U}:=\iint_{U}|q|
$$

is well defined and finite.
For example

$$
\left\|\frac{\mathrm{dz}^{2}}{z}\right\|_{\{|z|<R\}}=\int_{0}^{2 \pi} \int_{0}^{R} \frac{1}{r} r d r d \theta=2 \pi R
$$

For $f: \mathbb{C} \rightarrow \mathbb{C}$ a non-constant polynomial and $\mathcal{Q}=q \mathrm{dz}^{2}$ a meromorphic quadratic differential on $\mathbb{C}$, the pushforward $f_{*} \mathcal{Q}$ is defined by the quadratic differential

$$
f_{*} \mathcal{Q}:=T q \mathrm{dz}^{2} \quad \text { with } \quad T q(z):=\sum_{f(w)=z} \frac{q(w)}{f^{\prime}(w)^{2}}
$$

If $Q \in \mathcal{Q}(\mathbb{C})$, then $f_{*} \mathcal{Q} \in \mathcal{Q}(\mathbb{C})$ also.
The following lemma is a weak version of Thurston's contraction principle.
Lemma 2.2 (contraction principle). For a non-constant polynomial $f$ and a round disk $V$ of radius large enough so that $U:=f^{-1}(V)$ is relatively compact in $V$, we have

$$
\left\|f_{*} \mathcal{Q}\right\|_{V} \leq\|\mathcal{Q}\|_{U}<\|\mathcal{Q}\|_{V}, \quad \forall \mathcal{Q} \in \mathcal{Q}(\mathbb{C})
$$

Proof. The strict inequality on the right is a consequence of the fact that $U$ is relatively compact in $V$. The inequality on the left comes from

$$
\begin{aligned}
\left\|f_{*} \mathcal{Q}\right\|_{V} & =\iint_{z \in V}\left|\sum_{f(w)=z} \frac{q(w)}{f^{\prime}(w)^{2}}\right||\mathrm{d} z|^{2} \\
& \leq \iint_{z \in V} \sum_{f(w)=z}\left|\frac{q(w)}{f^{\prime}(w)^{2}}\right||\mathrm{d} z|^{2} \\
& =\iint_{w \in U}|q(w)||\mathrm{d} w|^{2}=\|\mathcal{Q}\|_{U}
\end{aligned}
$$

Corollary 2.3. If $f: \mathbb{C} \rightarrow \mathbb{C}$ is a polynomial and if $\mathcal{Q} \in \mathcal{Q}(\mathbb{C})$, then $f_{*} \mathcal{Q} \neq \mathcal{Q}$.
Remark 2.1. Thurston's contraction principal says that if $\mathcal{Q}$ is a meromorphic quadratic differential on $\mathbb{P}^{1}$ and $f: \mathbb{P}^{1} \rightarrow \mathbb{P}^{1}$ is a rational function, if one requires $f_{*} \mathcal{Q}=\mathcal{Q}$ with $\mathcal{Q} \neq 0$, then $f$ is necessarily a Lattès example.

The formulas below appeared in [L] chapter 2, we write them together as a lemma.

Lemma 2.4 (Levin). For $f=f_{c}$, we have

$$
\begin{cases}f_{*}\left(\frac{\mathrm{dz}^{2}}{z}\right)=0 & \text { if } a \neq 0  \tag{2.1}\\ f_{*}\left(\frac{\mathrm{dz}^{2}}{z-a}\right)=\frac{1}{f^{\prime}(a)}\left(\frac{\mathrm{dz}^{2}}{z-f(a)}-\frac{\mathrm{dz}^{2}}{z-c}\right) & \text { if } a \neq 0 \\ f_{*}\left(\frac{\mathrm{dz}^{2}}{(z-a)^{2}}\right)=\frac{\mathrm{dz}^{2}}{(z-f(a))^{2}}-\frac{d-1}{a f^{\prime}(a)}\left(\frac{\mathrm{dz}^{2}}{z-f(a)}-\frac{\mathrm{dz}^{2}}{z-c}\right)\end{cases}
$$

### 2.2 Proof of Theorem 2.1

Lemma 2.5 (compare with [Mil]). Given $z \in \mathbb{C}$, for $n \geq 0$ and $d \geq 2$, define $z_{n}: c \mapsto$ $f_{c}^{\circ n}(z)$ and $\delta_{n}=f_{c}^{\prime}\left(z_{n}\right)=d z_{n}^{d-1}$. Then

$$
\frac{\mathrm{dz}}{n} \mathrm{dc}=1+\delta_{n-1}+\delta_{n-1} \delta_{n-2}+\ldots+\delta_{n-1} \delta_{n-2} \cdots \delta_{1}
$$

Proof. From $z_{n}=z_{n-1}^{d}+c, d \geq 2$, we obtain

$$
\frac{\mathrm{dz}_{n}}{d c}=1+\delta_{n-1} \frac{\mathrm{dz}_{n-1}}{d c} \quad \text { with } \quad \frac{\mathrm{dz}_{0}}{d c}=0
$$

The result follows by induction.
Proof of Theorem 2.1.
Let $P_{n}(c, z):=f_{c}^{\circ n}(z)-z$ and consider the algebraic curve

$$
Y_{n}:=\left\{(c, z) \in \mathbb{C}^{2} \mid P_{n}(c, z)=0\right\}
$$

If $(c, z) \in Y_{n}$, the point $z$ is periodic for $f_{c}$ of period $m \leq n$. Then $m$ divides $n^{1}$. Therefore $Y_{n}$ is the set of $(c, z)$ such that $z$ is periodic for $f_{c}$ of period $m \leq n$ and $m$ dividing $n$.

As $Y_{n}$ is a closed subset of $\mathbb{C}^{2}$, we have $\overline{X_{n}} \subset Y_{n}$.
We decompose $Y_{n}$ into

$$
\begin{aligned}
Y_{n}= & X_{n} \\
& \sqcup\left\{(c, z) \mid z \text { is of period } n \text { for } f_{c} \text { with multiplier } 1\right\} \\
& \sqcup\left\{(c, z) \mid z \text { is of period } m \text { for } f_{c} \text { with } m<n \text { and } m \text { dividing } n\right\}
\end{aligned}
$$

We will examine case by case points in $Y_{n}$, determine points in $\overline{X_{n}}$ and establish the smoothness of $\overline{X_{n}}$ at each of these points.
Case 1. Consider a point $\left(c_{0}, z_{0}\right) \in X_{n} \subset Y_{n}$.

[^1]If $(c, z) \in Y_{n}$ is close to $\left(c_{0}, z_{0}\right) \in X_{n}$, the points of the orbit of $z$ are close to points of the orbit of $z_{0}$ and there are therefore at least $n$ distinct points in the orbit of $z$. It follows that the period of $z$ is equal to $n$. This shows that in a neighborhood of $\left(c_{0}, z_{0}\right)$, the curves $X_{n}$ and $Y_{n}$ coincide. It suffices to show that $Y_{n}$ is smooth in a neighborhood of $\left(c_{0}, z_{0}\right)$. As $\left[f_{c_{0}}^{\circ n}\right]^{\prime}\left(z_{0}\right) \neq 1$, we have

$$
\frac{\partial P_{n}}{\partial z}\left(c_{0}, z_{0}\right) \neq 0
$$

The implicit function theorem implies that $Y_{n}$, therefore $X_{n}$, is smooth in a neighborhood of $\left(c_{0}, z_{0}\right)$.

Case 2. Now consider a point $\left(c_{0}, z_{0}\right) \in Y_{n}$ such that $z_{0}$ is of period equal to $n$ for $f_{c_{0}}$ with multiplier 1.

Fix any $\ell \geq n$ that is a multiple of $n$. And consider $P_{\ell}$ and $Y_{\ell}$. We know that

$$
\begin{equation*}
\left(c_{0}, z_{0}\right) \in Y_{\ell} \quad \text { and } \quad\left[f_{c_{0}}^{\ell}\right]^{\prime}\left(z_{0}\right)=1 \tag{2.2}
\end{equation*}
$$

Claim. For any triple $\left(c_{0}, z_{0}, \ell\right)$ satisfying (2.2), we have $\frac{\partial P_{\ell}}{\partial c}\left(c_{0}, z_{0}\right) \neq 0$.
Proof. For $k \geq 0$, define inductively $z_{k+1}=f_{c_{0}}\left(z_{k}\right)$ and define $\delta_{k}:=f_{c_{0}}^{\prime}\left(z_{k}\right)$. We have, by Lemma 2.5

$$
\frac{\partial P_{\ell}}{\partial c}\left(c_{0}, z_{0}\right)=\left.\frac{d}{d c}\left(f_{c}^{\circ \ell}\left(z_{0}\right)-z_{0}\right)\right|_{c_{0}}=1+\delta_{\ell-1}+\delta_{\ell-1} \delta_{\ell-2}+\ldots+\delta_{\ell-1} \delta_{\ell-2} \cdots \delta_{1} .
$$

Now consider the quadratic differential $\mathcal{Q} \in \mathcal{Q}(\mathbb{C})$ defined by

$$
\mathcal{Q}(z):=\sum_{k=0}^{\ell-1} \frac{\rho_{k}}{z-z_{k}} \mathrm{dz}^{2}, \quad \text { with } \rho_{k}=\delta_{\ell-1} \delta_{\ell-2} \cdots \delta_{k} .
$$

Applying Lemma 2.4, and writing $f$ for $f_{c_{0}}$, we obtain

$$
f_{*} \mathcal{Q}(z)=\sum_{k=0}^{\ell-1} \frac{\rho_{k}}{\delta_{k}}\left(\frac{\mathrm{dz}^{2}}{z-z_{k+1}}-\frac{\mathrm{dz}^{2}}{z-c_{0}}\right)=\mathcal{Q}(z)-\frac{\partial P_{\ell}}{\partial c}\left(c_{0}, z_{0}\right) \cdot \frac{\mathrm{dz}^{2}}{z-c_{0}}
$$

By Corollary 2.3, we can not have $f_{*} \mathcal{Q}=\mathcal{Q}$. It follows that

$$
\frac{\partial P_{\ell}}{\partial c}\left(c_{0}, z_{0}\right) \neq 0
$$

This ends the proof of the claim.
Now let $\ell=n$, by implicit function theorem, there exists unique locally holomorphic function $c(z)$ with $f_{c(z)}^{n}(z)=z, c\left(z_{0}\right)=z_{0}$ and $c^{\prime}\left(z_{0}\right)=0\left(\right.$ for $\left.\frac{\partial P_{\ell}}{\partial z}\left(c_{0}, z_{0}\right)=0\right)$. Then there is neighborhood $U$ of $\left(c_{0}, z_{0}\right)$ in $\mathbb{C}^{2}$ such that

$$
Y_{n} \cap U=\left\{(c(z), z)| | z-z_{0} \mid<\varepsilon\right\}
$$

As $z_{0}$ is a $n$ periodic point of $f_{c_{0}}$ and the map $z \mapsto\left[f_{c(z)}^{\circ n}\right]^{\prime}(z)$ is holomorphic and can not $\mathbf{b e}^{2}$ constantly 1 , we can choose $\varepsilon$ small enough such that $z$ is $n$ periodic point of $f_{c(z)}$ with multiplier $\neq 1$ for $\left|z-z_{0}\right|<\varepsilon$. Then

$$
U \cap Y_{n} \backslash\left\{\left(c_{0}, z_{0}\right)\right\} \subset U \cap X_{n} \subset U \cap Y_{n}
$$

It follows $\left(c_{0}, z_{0}\right) \in \partial X_{n}$ and $U \cap Y_{n}$ is a neighborhood of $\left(c_{0}, z_{0}\right)$ on $\overline{X_{n}}$. Then $\overline{X_{n}}$ is smooth at $\left(c_{0}, z_{0}\right)$ and parametered locally by $z$.

Case 3. Finally consider $\left(c_{0}, z_{0}\right) \in Y_{n}$ so that $z_{0}$ is of period $m<n$ for $f_{c_{0}}$ with $m$ dividing $n$.

Note that $Y_{m} \subset Y_{n}$.
If $\left[f_{c_{0}}^{n}\right]^{\prime}\left(z_{0}\right) \neq 1$ then $\left[f_{c_{0}}^{m}\right]^{\prime}\left(z_{0}\right) \neq 1$. By the existence and the unicity of the implicit function theorem the local solutions of $f_{c}^{n}(z)-z=0$ and $f_{c}^{m}(z)=z$ coincide, that is, $Y_{m}$ and $Y_{n}$ coincide locally. So at point $\left(c_{0}, z_{0}\right), Y_{n}$ is locally the graph of a holomorphic function $z(c)$ with $z\left(c_{0}\right)=z_{0}$ and $z(c)$ is $m$ periodic point of $f_{c}$. It follows that $\left(c_{0}, z_{0}\right) \notin$ $\overline{X_{n}}$.

If $\left[f_{c_{0}}^{n}\right]^{\prime}\left(z_{0}\right)=1$ and $\left[f_{c_{0}}^{m}\right]^{\prime}\left(z_{0}\right)=1$, then both triples $\left(c_{0}, z_{0}, m\right)$ and $\left(c_{0}, z_{0}, n\right)$ satisfy (2.2) . The claim in Case 2 and implicit function theorem imply that $Y_{m}$ and $Y_{n}$ again coincide in a neighborhood of $\left(c_{0}, z_{0}\right)$. For the same reason as above, $\left(c_{0}, z_{0}\right) \notin \overline{X_{n}}$.

Set $\rho:=\left[f_{c_{0}}^{m}\right]^{\prime}\left(z_{0}\right)$. We consider now the only remaining case $\rho \neq 1$ and $\rho^{n / m}=$ $\left[f_{c_{0}}^{n}\right]^{\prime}\left(z_{0}\right)=1$.

Fix any integer $s \geq 2$ such that $\rho^{s}=1$. Let $c_{*}$ be any point outside Mandelbrot set, then each zero point of $f_{c_{*}}^{m}(z)-z$ is simple. It follows that $f_{c_{*}}^{m}(z)-z$ divides $f_{c_{*}}^{\circ m s}(z)-z$. Since $c_{*}$ is any point outside Mandelbrot set, the polynomial $f_{c}^{m}(z)-z$ must divides $f_{c}^{m s}(z)-z$. Let $P(c, z)$ be the polynomial defined by the equation:

$$
\begin{equation*}
f_{c}^{\circ m s}(z)-z=\left(f_{c}^{\circ m}(z)-z\right) \cdot P(c, z) . \tag{2.3}
\end{equation*}
$$

Claim. Let $Z_{s}:=\{(c, z) \mid P(c, z)=0\}$. Then $\left(c_{0}, z_{0}\right) \in Z_{s}$ and there is a neighborhood $V$ of $\left(c_{0}, z_{0}\right)$ in $\mathbb{C}^{2}$ such that

$$
Z_{s} \cap V=\left\{(c(z), z)| | z-z_{0} \mid<\varepsilon_{0}, c(z) \text { is holomorphic with } c\left(z_{0}\right)=c_{0} \text { and } c^{\prime}\left(z_{0}\right)=0\right\} .
$$

Proof. We will prove at first that the map $z \mapsto f_{c_{0}}^{m s}(z)-z$ has a zero of order at least 3 at $z_{0}$. Define $F(z)=f^{m}\left(z+z_{0}\right)-z_{0}$, then it is equivalent to show the function $F^{s}(z):=f_{c_{0}}^{m s}\left(z+z_{0}\right)-z_{0}$ has a local expansion $z+O\left(z^{3}\right)$ at 0 . We have $F(z)=\rho z+a z^{2}+O\left(z^{3}\right)$ in a neighborhood of 0 . One checks by induction

$$
\forall k \geq 1, \quad F^{\circ k}(z)=\rho^{k} z+a \rho^{k-1}\left(1+\rho+\rho^{2}+\cdots+\rho^{k-1}\right) z^{2}+O\left(z^{3}\right)
$$

[^2]But $\rho \neq 1$ and $\rho^{s}=1$, it follows that $1+\rho+\rho^{2}+\cdots+\rho^{s-1}=0$ and $F^{o s}(z)=z+O\left(z^{3}\right)$.
Since $z \mapsto f_{c_{0}}^{\circ m}(z)-z$ has a simple zero, we see from (2.3) that $z \mapsto P\left(c_{0}, z\right)$ has a zero of order at least 2 at $z_{0}$. Therefore $\left(c_{0}, z_{0}\right) \in Z_{s}$ and

$$
\begin{equation*}
\frac{\partial P}{\partial z}\left(c_{0}, z_{0}\right)=0 \tag{2.4}
\end{equation*}
$$

We proceed now to prove

$$
\begin{equation*}
\frac{\partial P}{\partial c}\left(c_{0}, z_{0}\right) \neq 0 \tag{2.5}
\end{equation*}
$$

This will be down in two steps:
Step 1. Let $Q(c, z):=f_{c}^{\circ m}(z)-z$. We have

$$
Q\left(c_{0}, z_{0}\right)=0 \quad \text { and } \quad \frac{\partial Q}{\partial z}\left(c_{0}, z_{0}\right)=\rho-1 \neq 0
$$

According to the implicit function theorem, there is a germ of a holomorphic function $\zeta:\left(\mathbb{C}, c_{0}\right) \rightarrow\left(\mathbb{C}, z_{0}\right)$ with $Q(c, \zeta(c))=0$. In other words, $\zeta(c)$ is a periodic point of period $m$ for $f_{c}$. Let $\rho_{c}$ denote the multiplier of $\zeta(c)$ for $f_{c}$ and set

$$
\dot{\rho}:=\left.\frac{d \rho_{c}}{d c}\right|_{c_{0}} .
$$

Lemma 2.6. We have

$$
\frac{\partial P}{\partial c}\left(c_{0}, z_{0}\right)=\frac{s \cdot \dot{\rho}}{\rho(\rho-1)} .
$$

Proof. Differentiating the equation (2.3) with respect to $z$, and then evaluating at $(c, \zeta(c))$, we get:

$$
\rho_{c}^{s}-1=\left(\rho_{c}-1\right) \cdot P(c, \zeta(c))+\underbrace{\left(f_{c}^{m}(\zeta(c))-\zeta(c)\right)}_{=0} \cdot \frac{\partial P}{\partial z}(c, \zeta(c))=\left(\rho_{c}-1\right) \cdot P(c, \zeta(c)) .
$$

Setting

$$
R(c):=P(c, \zeta(c))=\frac{\rho_{c}^{s}-1}{\rho_{c}-1}
$$

we have

$$
R^{\prime}\left(c_{0}\right)=\frac{\partial P}{\partial c}\left(c_{0}, z_{0}\right)+\underbrace{\frac{\partial P}{\partial z}\left(c_{0}, z_{0}\right)}_{=0} \cdot \zeta^{\prime}\left(c_{0}\right)=\frac{\partial P}{\partial c}\left(c_{0}, z_{0}\right)
$$

Using $\rho^{s}=1$ and $\rho^{s-1}=1 / \rho$, we deduce that

$$
\frac{\partial P}{\partial c}\left(c_{0}, z_{0}\right)=\left.\frac{d}{d c}\left(\frac{\rho_{c}^{s}-1}{\rho_{c}-1}\right)\right|_{c_{0}}=\left.\left(\frac{s \rho^{s-1}}{\rho-1}-\frac{\rho^{s}-1}{(\rho-1)^{2}}\right) \frac{d \rho_{c}}{d c}\right|_{c_{0}}=\frac{s \cdot \dot{\rho}}{\rho(\rho-1)}
$$

Step 2. $\dot{\rho} \neq 0$. The proof of this fact will be postponed to the following section 2.3 using quadratic differential with double poles (see also [DH] for a parabolic implosion approach).

This ends the proof of (2.5), as well as the proof of the claim by combining (2.5) and the implicit function theorem plus the observation that $\left(c_{0}, z_{0}\right) \in Z_{s}$.

Write now $\rho=e^{2 \pi i u / v}$ with $u, v$ co-prime and $v>0$. Then any $s$ satisfying $\rho^{s}=1$ takes the form $s=k v$ for some integer $k \geq 1$. With the same reason as that of existence of polynomial $P(c, z)$ in (2.3), there are polynomials $g$, $h$ such that

$$
f_{c}^{\circ m s}(z)-z=f_{c}^{\circ m k v}(z)-z=\left(f_{c}^{\circ m v}(z)-z\right) g(c, z)=\left(f_{c}^{\circ m}(z)-z\right) h(c, z) g(c, z) .
$$

By definition we have $Z_{s}=\{(c, z) \mid g(c, z) h(c, z)=0\} \supset\{(c, z) \mid h(c, z)=0\}=Z_{v}$. By the claim in Case 3, we conclude that $Z_{s}$ and $Z_{v}$ coincide in a neighborhood of $\left(c_{0}, z_{0}\right)$ as the graph of a single holomorphic function $c(z)$ with vanishing derivative at $z_{0}$.

Remark:(1) If necessary, we can decrease $\varepsilon_{0}$ in claim of case 3 such that $f_{c(z)}^{m}(z)-z \neq 0$ for $0<\left|z-z_{0}\right|<\varepsilon_{0}$. Otherwise, there exist a sequence $\left\{z_{k}\right\}$ with $z_{k} \rightarrow z_{0}$ ( correspondingly, $\left.c_{k}:=c\left(z_{k}\right) \rightarrow c_{0}\right)$ such that $f_{c_{k}}^{m}\left(z_{k}\right)-z_{k}=0$ and $h\left(c_{k}, z_{k}\right) g\left(c_{k}, z_{k}\right)=0$. It follows that $\left[f_{c_{k}}^{\circ m s}\right]^{\prime}\left(z_{k}\right)-1=0$, that is, $\left\{c_{k}\right\}$ is a sequence of parabolic parameter with period of parabolic orbit less than $m$ converging to $c_{0}$. It is impossible.

$$
\begin{equation*}
Z_{s}=\left(Y_{n} \backslash Y_{m}\right) \cup\left\{\left(c_{0}, z_{0}\right)\right\}, X_{n} \subset Y_{n} \backslash Y_{m} \tag{2}
\end{equation*}
$$

Lemma 2.7. There exists $0<\varepsilon_{1}<\varepsilon_{0}$ such that $z$ is mv periodic point of $f_{c(z)}$ with multiplier $\neq 1$ for $0<\left|z-z_{0}\right|<\varepsilon_{1} . c(z)$ is defined in the claim of case 3.

Proof. Note that $P(c(z), z)=0$ implies $z$ is periodic point of $f_{c(z)}$ with period less than $m s$. As $\left(f_{c_{0}}^{m}\right)^{\prime}\left(z_{0}\right)=\rho=e^{2 \pi i u / v}$, by lemma 3.9 below, when $c$ is close enough to $c_{0}$, the orbit of $f_{c_{0}}$ containing $z_{0}$ splits into two periodic orbit of $f_{c}$ with period $m$ and $m v$. Then we can choose $\varepsilon_{1}<\varepsilon_{0}$ such that $z$ belongs to one of the two splitted orbits of $f_{c(z)}$ for $0<\left|z-z_{0}\right|<\varepsilon_{1}$. By remark (1), the period of $z$ under $f_{c(z)}$ must be $m v$. The parabolic parameter in $M_{d}$ with period of parabolic point less than a fixed number are finite, so we can decrease $\varepsilon_{1}$ if necessary, such that $c(z)$ is not parabolic parameter for $0<\left|z-z_{0}\right|<\varepsilon_{1}$.

Now let $V_{1}$ be a neighborhood of $\left(c_{0}, z_{0}\right)$ in $\mathbb{C}^{2}$ with property that

$$
V_{1} \cap Z_{s}=V_{1} \cap Z_{v}=\left\{(c(z), z) \| z-z_{0} \mid<\varepsilon_{1}\right\}
$$

If $n=m v$, by lemma 2.7 and remark (2), we have

$$
\left(V_{1} \cap Z_{v}\right) \backslash\left\{\left(c_{0}, z_{0}\right)\right\} \subset V_{1} \cap X_{n} \subset V_{1} \cap\left(Y_{n} \backslash Y_{m}\right)=\left(V_{1} \cap Z_{v}\right) \backslash\left\{\left(c_{0}, z_{0}\right)\right\} .
$$

It follows $\left(c_{0}, z_{0}\right) \in \partial X_{n}$ and $\overline{X_{n}}$ coincides with $Z_{v}$ at neighborhood of $\left(c_{0}, z_{0}\right)$. Then $\overline{X_{n}}$ is smooth at point $\left(c_{0}, z_{0}\right)$

If $n=m v k$ for some $k>1$

$$
\left(V_{1} \cap Z_{s}\right) \backslash\left\{\left(c_{0}, z_{0}\right)\right\} \subset V_{1} \cap\left(Y_{n} \backslash Y_{m}\right)=\left(V_{1} \cap Z_{s}\right) \backslash\left\{\left(c_{0}, z_{0}\right)\right\}
$$

Then $V_{1} \cap Z_{s}$ is the neighborhood of $\left(c_{0}, z_{0}\right)$ in $\left(Y_{n} \backslash Y_{m}\right) \cup\left\{\left(c_{0}, z_{0}\right)\right\}$. For $X_{n} \subset Y_{n} \backslash Y_{m}$ and $X_{n} \cap\left(V_{1} \cap Z_{s}\right)=\emptyset$, we have $\left(c_{0}, z_{0}\right) \notin \overline{X_{n}}$

### 2.3 Quadratic differentials with double poles

Set $f:=f_{c_{0}}$,

$$
z_{k}:=f^{k}\left(z_{0}\right), \quad \delta_{k}:=d z_{k}^{d-1}=f^{\prime}\left(z_{k}\right), \quad \zeta_{k}(c):=f_{c}^{\circ k}(\zeta(c)) \quad \text { and } \quad \dot{\zeta}_{k}:=\zeta_{k}^{\prime}\left(c_{0}\right)
$$

Then

$$
\zeta_{k+1}(c)=f_{c}\left(\zeta_{k}(c)\right) \quad \text { and } \quad \zeta_{m}=\zeta_{0}
$$

Since

$$
\delta_{0} \delta_{1} \cdots \delta_{m-1}=\rho \neq 0
$$

there is a unique $m$-tuple $\left(\mu_{0}, \ldots, \mu_{m-1}\right)$ such that

$$
\mu_{k+1}=\frac{\mu_{k}}{d z_{k}^{d-1}}-\frac{d-1}{d z_{k}^{d}}
$$

where the indices are considered to be modulo $m$.
Now consider the quadratic differential $\mathcal{Q}$ (with double poles) defined by

$$
\mathcal{Q}:=\sum_{k=0}^{m-1}\left(\frac{1}{\left(z-z_{k}\right)^{2}}+\frac{\mu_{k}}{z-z_{k}}\right) \mathrm{dz}^{2}
$$

Lemma 2.8 (Compare with [L]). We have

$$
f_{*} \mathcal{Q}=\mathcal{Q}-\frac{\dot{\rho}}{\rho} \cdot \frac{\mathrm{dz}^{2}}{z-c_{0}} .
$$

Proof. By construction of $\mathcal{Q}$ and the calculation of $f_{*} \mathcal{Q}$ in Lemma 2.4, the polar parts of $\mathcal{Q}$ and $f_{*} \mathcal{Q}$ along the cycle of $z_{0}$ are identical. But $f_{*} \mathcal{Q}$ has an extra simple pole at the critical value $c_{0}$ with coefficient

$$
\sum_{k=0}^{m-1}\left(-\frac{\mu_{k}}{d z_{k}^{d-1}}+\frac{d-1}{d z_{k}^{d}}\right)=-\sum_{k=0}^{m-1} \mu_{k+1}
$$

We need to show that this coefficient is equal to $-\frac{\dot{\rho}}{\rho}$.
Using $\zeta_{k+1}(c)=\zeta_{k}(c)^{d}+c$, we get

$$
\dot{\zeta}_{k+1}=d z_{k}^{d-1} \dot{\zeta}_{k}+1
$$

It follows that

$$
\dot{\zeta}_{k+1} \mu_{k+1}-\mu_{k+1}=d z_{k}^{d-1} \dot{\zeta}_{k} \mu_{k+1}=\dot{\zeta}_{k} \mu_{k}-\frac{(d-1) \dot{\zeta}_{k}}{z_{k}}
$$

Therefore

$$
\sum_{k=0}^{m-1} \mu_{k+1}=\sum_{k=0}^{m-1}\left(\dot{\zeta}_{k+1} \mu_{k+1}-\dot{\zeta}_{k} \mu_{k}+\frac{(d-1) \dot{\zeta}_{k}}{z_{k}}\right)=(d-1) \sum_{k=0}^{m-1} \frac{\dot{\zeta}_{k}}{z_{k}}=\frac{\dot{\rho}}{\rho}
$$

where last equality is obtained by evaluating at $c_{0}$ of the logarithmic derivative of

$$
\rho_{c}:=\prod_{k=0}^{m-1} d \zeta_{k}^{d-1}(c)
$$

Lemma 2.9 (Epstein $[\mathrm{E}])$. We have $f_{*} \mathcal{Q} \neq \mathcal{Q}$.
Proof. The proof rests again on the contraction principle, but we can not apply directly Lemma 2.2 since $\mathcal{Q}$ is not integrable near the cycle $\left\langle z_{0}, \ldots, z_{m-1}\right\rangle$. Consider a sufficiently large round disk $V$ so that $U:=f^{-1}(V)$ is relatively compact in $V$. Given $\varepsilon>0$, we set

$$
V_{\varepsilon}:=\bigcup_{k=1}^{m} f^{k}\left(D\left(z_{0}, \varepsilon\right)\right) \quad \text { and } \quad U_{\varepsilon}:=f^{-1}\left(V_{\varepsilon}\right) .
$$

When $\varepsilon$ tends to 0 , we have

$$
\left\|f_{*} \mathcal{Q}\right\|_{V-V_{\varepsilon}} \leq\|\mathcal{Q}\|_{U-U_{\varepsilon}}=\|\mathcal{Q}\|_{V-V_{\varepsilon}}-\|\mathcal{Q}\|_{V-U}+\|\mathcal{Q}\|_{V_{\varepsilon}-U_{\varepsilon}}-\|\mathcal{Q}\|_{U_{\varepsilon}-V_{\varepsilon}}
$$

If we had $f_{*} \mathcal{Q}=\mathcal{Q}$, we would have

$$
0<\|\mathcal{Q}\|_{V-U} \leq\|\mathcal{Q}\|_{V_{\varepsilon}-U_{\varepsilon}}
$$

However, $\|\mathcal{Q}\|_{V_{\varepsilon}-U_{\varepsilon}}$ tends to 0 as $\varepsilon$ tends to 0 , which is a contradiction. Indeed, $\mathcal{Q}=$ $q(z) d z^{2}$, the meromorphic function $q$ is equivalent to $\frac{1}{\left(z-z_{0}\right)}$ as $z$ tends to $z_{0}$. In addition, since the multiplier of $z_{0}$ has modulus 1 ,

$$
D\left(z_{0}, \varepsilon\right) \subset U_{\varepsilon}-V_{\varepsilon} \subset D\left(z_{0}, \varepsilon^{\prime}\right) \quad \text { with } \quad \frac{\varepsilon^{\prime}}{\varepsilon} \xrightarrow{\varepsilon \rightarrow 0} 1 .
$$

Therefor,

$$
\|\mathcal{Q}\|_{V_{\varepsilon}-U_{\varepsilon}} \leq \int_{0}^{2 \pi} \int_{\varepsilon}^{\varepsilon^{\prime}} \frac{1+o(1)}{r^{2}} r d r d \theta=2 \pi(1+o(1)) \log \frac{\varepsilon^{\prime}}{\varepsilon} \xrightarrow{\varepsilon \rightarrow 0} 0
$$

The fact $\dot{\rho} \neq 0$ follows from the above two lemmas.

## 3 The irreducibility of the periodic curves

Recall that $f_{c}$ denote the polynomial $z \mapsto z^{d}+c$, where $d \geq 2$, and we have defined

$$
X_{n}:=\left\{(c, z) \in \mathbb{C}^{2} \mid f_{c}^{n}(z)=z,\left[f_{c}^{n}\right]^{\prime}(z) \neq 1 \text { and for all } 0<m<n, \quad f_{c}^{m}(z) \neq z\right\} .
$$

The objective here is to prove:
Theorem 3.1. For every $n \geq 1$, the set $X_{n}$ is connected.
It follows immediately that the closure of $X_{n}$ in $\mathbb{C}^{2}$ is irreducible.

### 3.1 Kneading sequences

Set $\mathbb{T}=\mathbb{R} / \mathbb{Z}$ and let $\tau: \mathbb{T} \rightarrow \mathbb{T}$ be the angle map

$$
\tau: \mathbb{T} \ni \theta \mapsto d \theta \in \mathbb{T}, d \geq 2
$$

We shall often make the confusion between an angle $\theta \in \mathbb{T}$ and its representative in $[0,1[$. In particular, the angle $\theta / d \in \mathbb{T}$ is the element of $\tau^{-1}(\theta)$ with representative in $[0,1 / d[$ and the angle $(\theta+(d-1)) / d$ is the element of $\tau^{-1}(\theta)$ with representative in $[(d-1) / d, 1[$.

Every angle $\theta \in \mathbb{T}$ has an associated kneading sequence $\nu(\theta)=\nu_{1} \nu_{2} \nu_{3} \ldots$ defined by

For example,

- as $d=3, \nu\left(\frac{1}{7}\right)=\overline{12102 \star} \quad$ and $\quad \nu\left(\frac{27}{28}\right)=\overline{22200 \star}$;

We shall say that an angle $\theta \in \mathbb{T}$, periodic under $\tau$, is maximal in its orbit if its representative in $[0,1)$ is maximal among the representatives of $\tau^{j}(\theta)$ in $[0,1)$ for all $j \geq 1$. If the period is $n$ and the d-expansion $(d \geq 2)$ of $\theta$ is.$\overline{\varepsilon_{1} \ldots \varepsilon_{n}}$, then $\theta$ is maximal


Figure 1: As $d=3$, the kneading sequence of $\theta=1 / 7$ is $\nu(1 / 7)=\overline{12102 \star}$
in its orbit if and only if the periodic sequence $\overline{\varepsilon_{1} \ldots \varepsilon_{n}}$ is maximal (in the lexicographic order) among its shifts. For example, as $d=4, \frac{5}{31}=. \overline{02211}$ is not maximal in its orbit but $\frac{20}{31}=. \overline{22110}$ is maximal in the same orbit.

The following lemma indicates cases where the $d$-expansion $(d \geq 2)$ and the kneading sequence coincide.

Lemma 3.2 (Realization of kneading sequences). Let $\theta \in \mathbb{T}$ be a periodic angle which is maximal in its orbit and let. $\overline{\varepsilon_{1} \ldots \varepsilon_{n}}$ be its d-expansion $(d \geq 2)$. Then, $\varepsilon_{n} \in\{0,1,2, \ldots, d-$ $2\}$ and the kneading sequence $\nu(\theta)$ is equal to $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1} \star}$.

For example,

- as $d=3 \quad \frac{13}{14}=\overline{221001}$ and $\quad \nu(\theta)=\overline{22100 \star}$.
- as $d=4 \quad \frac{28}{31}=\overline{32130} \quad$ and $\quad \nu(\theta)=\overline{3213 \star}$.

Proof. Since $\theta$ is maximal in its orbit under $\tau$, the orbit of $\theta$ is disjoint from $\left.\left.\left.] \frac{\theta}{d}, \frac{1}{d}\right] \cup\right] \frac{\theta+1}{d}, \frac{2}{d}\right]$ $\left.\left.\left.\bigcup \ldots \bigcup] \frac{\theta+(d-2)}{d}, \frac{d-1}{d}\right] \bigcup\right] \theta, 1\right]$. It follows that the orbit $\tau^{j}(\theta), j=0,1, \ldots, n-2$ have the same itinerary relative to the two partitions $\mathbb{T}-\left\{0, \frac{1}{d}, \frac{2}{d}, \ldots, \frac{d-2}{d}, \frac{d-1}{d}\right\}$ and $\mathbb{T}-\left\{\frac{\theta}{d}, \frac{\theta+1}{d}, \ldots, \frac{\theta+(d-2)}{d}, \frac{\theta+(d-1)}{d}\right\}$ (see Figure 2). The first one gives the dexpansion $(d \geq 2)$ whereas the second gives the kneading sequence. Therefore, the kneading sequence of $\theta$ is $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1 \star}}$. Since $\tau^{n-1}(\theta) \in \tau^{-1}(\theta)=\left\{\frac{\theta}{d}, \frac{\theta+1}{d}, \ldots, \frac{\theta+(d-1)}{d}\right\}$ and


Figure 2: As $d=4$, the kneading sequence of $\theta=28 / 31$ is $\nu(28 / 31)=\overline{3213 \star}$
since $\left.\left.\frac{\theta+(d-1)}{d} \in\right] \theta, 1\right]$, we must have $\tau^{n-1}(\theta)=\left\{\frac{\theta}{d}, \frac{\theta+1}{d}, \ldots, \frac{\theta+(d-2)}{d}\right\}<\frac{d-1}{d}$. So $\varepsilon_{n}$, as the first digit of $\tau^{n-1}(\theta)$, must be in $\{0,1,2, \ldots, d-2\}$.

### 3.2 Cyclic expression of kneading sequence

$X=\{0,1, \ldots, d-1\}(d \geq 2)$ is an alphabet. $X^{\star}$ is the set of all sequence of symbols from $X$ with finite length, that is,

$$
X^{\star}=\left\{\nu_{1} \ldots \nu_{t} \mid \nu_{i} \in X, t \in \mathbb{N}^{\star}\right\} .
$$

The element of $X^{\star}$ is called word, its length is denoted by $|\cdot|$. For any $w \in X^{\star}, w$ can be written as $u^{n}:=\underbrace{u \ldots u}_{n}$ with $u \in X^{\star}$ and $n \geq 1$.

For example: $\quad 121212=12^{3}, \quad 1234=1234$.
Definition 3.3. A word is called primitive if it is not the form $u^{n}$ for any $n>1, u \in X^{\star}$.
The following lemma is a basic result about primitive words due to F.W.Levi. One can refer to $[\mathrm{KM}]$ for the proof.
Lemma 3.4 (F.W.Levi). For each $w \in X^{\star}$, there exists an unique primitive word $a(w)$ such that $w=a(w)^{n}$ for some $n \geq 1$.
$a(w)$ is called the primitive root of $w$, this lemma means the primitive root of a word is unique. Let $w$ be a word, we denote by $L_{w}$ the set of all words different from $w$ only at the last digit.

Lemma 3.5. If $w$ is a non-primitive word, then any word in $L_{w}$ is primitive.
Proof. As $w$ is not primitive, then $w=a^{m}$ where $a$ is the primitive root of $w$ and $m>1$. $w^{\prime}$ is any element of $L_{w}$, then $w^{\prime}=a^{m-1} a^{\prime}$ for some $a^{\prime} \in L_{a}$. Now assume $w^{\prime}$ is not primitive, then $w^{\prime}=z^{n}$ where $z$ is the primitive root of $w^{\prime}$ and $n>1$. Obviously $|z| \neq|a|$.

If $|z|<|a|$, then $n>m \geq 2$ and $a=z b$ for some $b \in X^{\star}$.

$$
\begin{aligned}
& a^{m-1} a^{\prime}=z^{n} \Longrightarrow z a^{m-1} a^{\prime}=a^{m-1} a^{\prime} z \Longrightarrow z a^{m-1} a^{\prime}=z b a^{m-2} a^{\prime} z \Longrightarrow \\
& \exists v \in X^{\star}, \text { s.t } a=b v,|v|=|z| \Longrightarrow a^{m-1} b v^{\prime}=b a^{m-2} a^{\prime} z\left(a^{\prime}=b v^{\prime}\right) \Longrightarrow \\
& v^{\prime}=z \text { and } a^{m-1} b=b a^{m-2} a^{\prime} \Longrightarrow a^{m-2} b v b=b a^{m-2} a^{\prime} \Longrightarrow a^{\prime}=v b
\end{aligned}
$$

It is a contradiction to $a=z b$.
If $|z|>|a|$, then there exists $z^{\prime} \in L_{z}$ such that $z^{n-1} z^{\prime}=a^{m}=w$ with $m>n \geq 2$. It reduces to the case above.

Now, let $\theta$ be a periodic angle with period $n \geq 2 . \nu(\theta)$ is the kneading sequence of $\theta$.
Definition 3.6. If there is a word $w=\nu_{1} \ldots \nu_{t}$ such that $\nu(\theta)=\overline{w^{s-1} w_{\star}}:=\underbrace{\overline{w \ldots w} w_{\star}}_{s-1}$, where $w_{\star}=\nu_{1} \ldots \nu_{t-1 \star}$ and $t$ is a proper factor of $n$ with $t s=n$, then $\nu(\theta)$ is called cyclic, otherwise $\nu(\theta)$ is called acyclic.
Definition 3.7. $\nu(\theta)=\overline{w^{s-1} w_{\star}}$ is cyclic. If $w$ is a primitive word, we call $\overline{w^{s-1} w_{\star}} a$ cyclic expression of $\nu(\theta)$.

The following proposition is a corollary of Lemma 3.4 and 3.5.
Proposition 3.8. If $\nu(\theta)$ is cyclic, then its cyclic expression is unique.
Proof. Assume $\overline{w^{s-1} w_{\star}}$ and $\overline{u^{l-1} u_{\star}}$ are two cyclic expression of $\nu(\theta)$ where $w=\nu_{1} \ldots \nu_{t}$ and $u=\epsilon_{1} \ldots \epsilon_{m}$. If $\nu_{t}=\epsilon_{m}$, then $w^{s}=u^{l}$. By Lemma 3.4, we have $w=u$. If $\nu_{t} \neq \epsilon_{m}$, then $w^{s}=u^{l-1} u^{\prime}$ with some $u^{\prime} \in L_{u}$, but this is a contradiction to Lemma 3.5.

### 3.3 Filled-in Julia sets and the Multibrot set

Let us recall some results about filled-in Julia set and Multibrot set that will be used following. These can be found in $[\mathrm{DH}],[\mathrm{Mil}]$ and $[\mathrm{DE}]$.

For $c \in \mathbb{C}$, we denote by $K_{c}$ the filled-in Julia set of $f_{c}$, that is the set of points $z \in \mathbb{C}$ whose orbit under $f_{c}$ is bounded. We denote by $M_{d}$ the Multibrot set for $f_{c}(z)=z^{d}+c$, that is the set of parameters $c \in \mathbb{C}$ for which the critical point 0 belongs to $K_{c}$.

If $c \in M_{d}$, then $K_{c}$ is connected. There is a conformal isomorphism $\phi_{c}: \mathbb{C} \backslash \bar{K}_{c} \rightarrow \mathbb{C} \backslash \overline{\mathbb{D}}$ which satisfies $\phi_{c} \circ f_{c}=\left(\phi_{c}\right)^{d}$ and $\phi_{c}^{\prime}(\infty)=1$. The dynamical ray of angle $\theta \in \mathbb{T}$ is

$$
R_{c}(\theta):=\left\{z \in \mathbb{C} \backslash K_{c} \mid \arg \left(\phi_{c}(z)\right)=2 \pi \theta\right\} .
$$

If $\theta$ is rational, then as $r$ tends to 1 from above, $\phi_{c}^{-1}\left(r \mathrm{e}^{2 \pi i \theta}\right)$ converges to a point $\gamma_{c}(\theta) \in K_{c}$. We say that $R_{c}(\theta)$ lands at $\gamma_{c}(\theta)$. We have $f_{c} \circ \gamma_{c}=\gamma_{c} \circ \tau$ on $\mathbb{Q} / \mathbb{Z}$. In particular, if $\theta$ is periodic under $\tau$, then $\gamma_{c}(\theta)$ is periodic under $f_{c}$. In addition, $\gamma_{c}(\theta)$ is either repelling (its multiplier has modulus $>1$ ) or parabolic (its multiplier is a root of unity).

If $c \notin M_{d}$, then $K_{c}$ is a Cantor set. There is a conformal isomorphism $\phi_{c}: U_{c} \rightarrow V_{c}$ between neighborhoods of $\infty$ in $\mathbb{C}$, which satisfies $\phi_{c} \circ f_{c}=\left(\phi_{c}\right)^{d}$ on $U_{c}$. We may choose $U_{c}$ so that $U_{c}$ contains the critical value $c$ and $V_{c}$ is the complement of a closed disk. For each $\theta \in \mathbb{T}$, there is an infimum $r_{c}(\theta) \geq 1$ such that $\phi_{c}^{-1}$ extends analytically along $R_{0}(\theta) \cap\left\{z \in \mathbb{C}\left|r_{c}(\theta)<|z|\right\}\right.$. We denote by $\psi_{c}$ this extension and by $R_{c}(\theta)$ the dynamical ray

$$
R_{c}(\theta):=\psi_{c}\left(R_{0}(\theta) \cap\left\{z \in \mathbb{C}\left|r_{c}(\theta)<|z|\right\}\right) .\right.
$$

As $r$ tends to $r_{c}(\theta)$ from above, $\psi_{c}\left(r \mathrm{e}^{2 \pi i \theta}\right)$ converges to a point $x \in \mathbb{C}$. If $r_{c}(\theta)>1$, then $x \in \mathbb{C} \backslash K_{c}$ is an iterated preimage of 0 and we say that $R_{c}(\theta)$ bifucates at $x$. If $r_{c}(\theta)=1$, then $\gamma_{c}(\theta):=x$ belongs to $K_{c}$ and we say that $R_{c}(\theta)$ lands at $\gamma_{c}(\theta)$. Again, $f_{c} \circ \gamma_{c}=\gamma_{c} \circ \tau$ on the set of $\theta$ such that $R_{c}(\theta)$ does not bifurcate. In particular, if $\theta$ is periodic under $\tau$ and $R_{c}(\theta)$ does not bifurcate, then $\gamma_{c}(\theta)$ is periodic under $f_{c}$.

The Multibrot set is connected. The map

$$
\phi_{M_{d}}: \mathbb{C} \backslash M_{d} \ni c \mapsto \phi_{c}(c) \in \mathbb{C} \backslash \overline{\mathbb{D}}
$$

is a conformal isomorphism. For $\theta \in \mathbb{T}$, the parameter ray $R_{M_{d}}(\theta)$ is

$$
R_{M_{d}}(\theta):=\left\{c \in \mathbb{C} \backslash M_{d} \mid \arg \left(\phi_{M_{d}}(c)\right)=2 \pi \theta\right\} .
$$

It is known that if $\theta$ is rational, then as $r$ tends to 1 from above, $\phi_{M_{d}}^{-1}\left(r \mathrm{e}^{2 \pi i \theta}\right)$ converges to a point $\gamma_{M_{d}}(\theta) \in M_{d}$. We say that $R_{M_{d}}(\theta)$ lands at $\gamma_{M_{d}}(\theta)$.

If $\theta$ is periodic for $\tau$ of exact period $n$ and if $c_{0}:=\gamma_{M_{d}}(\theta)$, then the point $\gamma_{c_{0}}(\theta)$ is periodic for $f_{c_{0}}$ with period $p$ dividing $n(p s=n, s \geq 1)$ and multiplier a $s$-th root of unity. If the period of $\gamma_{c_{0}}(\theta)$ for $f_{c_{0}}$ is exactly $n$ then the multiplier is $1, c_{0}$ is called primitive parabolic parameter, otherwise $c_{0}$ is called satellite parabolic parameter.
Lemma 3.9 (near parabolic map). $c_{0}$ is defined as above. When we make a small perturbation to $c_{0}$ in parameter space, If $c_{0}$ is a primitive parabolic parameter, then the parabolic orbit of $f_{c_{0}}$ is splitted into a pair of nearby periodic orbits of $f_{c}$, both have length $n$; If $c_{0}$ is a satellite parabolic parameter, then the parabolic orbit of $f_{c_{0}}$ is splitted into a pair of nearby periodic orbits of $f_{c}$, one has length $p$ and the other has length $s p=n$.

This lemma was proved by Milnor in [Mil] lemma 4.2 for the case $d=2$, but we can translate the proof word by word to the general case.


Figure 3: The parameter rays $R_{M_{3}}(7 / 26)$ and $R_{M_{3}}(9 / 26)$ land on a common root of a primitive hyperbolic component while $R_{M_{3}}(19 / 80)$ and $R_{M_{3}}(11 / 80)$ land on a common root of a satellite hyperbolic component. Only angles of rays are labelled in the graph.


Figure 4: The dynamical plane of $f_{c_{0}} . c_{0}:=\gamma_{M_{3}}(7 / 26)=\gamma_{M_{3}}(9 / 26)$ is the root of some primitive hyperbolic component as illustrated in Figure 3. The dynamical rays $R_{c_{0}}(7 / 26)$ and $R_{c_{0}}(9 / 26)$ land on a common parabolic point of $f_{c_{0}}$ with period 3 .

Let $H$ be periodic $n(n>1)$ hyperbolic component of $M_{d}$. For every parameter $c \in$ $H, f_{c}$ has an attracting periodic orbit $\left\{z(c), \ldots, f_{c}^{n-1}(z(c))\right\}$. Its multiplier define a map

$$
\mu_{H}: H \rightarrow \mathbb{D}, c \mapsto \frac{\partial}{\partial z} f_{c}^{n}(z(c))
$$

then $\mu_{H}: H \rightarrow \mathbb{D}$ is $d-1$ covering map with only one branched point .It extends continuously to a neighborhood of $\bar{H}$. Considering parameter $c \in \partial H$ such that $\mu_{H}(c)=1$, Eberlein proved that among these points, there is exactly one $c$ which is the landing point of two parameter rays of period $n$, this point is called root of $H$ (see Figure 3); the other $d-2$ points are landing points of only one parameter ray of period $n$ each, they are called co-root of $H$ (see Figure 6). $H$ is called primitive or satellite hyperbolic component according to whether its root is primitive or satellite parabolic parameter.

If $c$ is the root of some hyperbolic component and $c \neq \gamma_{M_{d}}(0)$, then two periodic parameter rays $R_{M_{d}}(\theta)$ and $R_{M_{d}}(\eta)$ land on $c$, we say $\theta$ and $\eta$ are companion angles, and $\theta, \eta$ have the same period under $\tau . c$ is primitive if and only if the orbit of $R_{M_{d}}(\theta)$ and $R_{M_{d}}(\eta)$ under $\tau$ are distinct. In dynamical plane, the dynamic rays $R_{c}(\theta)$ and $R_{c}(\eta)$ land at a common point $x_{1}:=\gamma_{c}(\theta)=\gamma_{c}(\eta)$. This point is on the parabolic orbit of $f_{c}$ with its immediate basin containing the critical value. $R_{c}(\theta)$ and $R_{c}(\eta)$ are adjacent to the Fatou component containing $c$ and the curve $R_{c}(\theta) \cup R_{c}(\eta) \cup\left\{x_{1}\right\}$ is a Jordan curve that cuts the plane into two connected components: one component, denoted by $V_{1}$, contains the critical value $c$; the other component, denoted by $V_{0}$, contains $R_{c}(0)$ and all points of parabolic cycle except $x_{1}$. Since $V_{1}$ contains the critical value, its preimage $U_{\star}=f_{c}^{-1}\left(V_{1}\right)$ is connected and contains the critical point 0 . It is bounded by the dynamical rays $R_{c}(\theta / d), \ldots, R_{c}((\theta+d-1) / d) ; R_{c}(\eta / d), \ldots, R_{c}((\eta+d-1) / d)$. Suppose $\theta>\eta$, and since each component of $\mathbb{C} \backslash \overline{U_{\star}}$ is conformally mapped to $V_{0}$ which is bounded by $R_{c}(\theta)$ and $R_{c}(\eta)$, it is easy to see that $R_{c}((\theta+k-1) / d)$ and $R_{c}((\eta+k) / d)$ land on a common point which is one of the preimage of $x_{1}$ for $k \in \mathbb{Z}_{d}$. Denote $U_{k}$ the component of $\mathbb{C} \backslash R_{c}((\theta+k-1) / d) \cup\left\{\gamma_{c}((\eta+k) / d)\right\} \cup R_{c}((\eta+k) / d)$ disjoint with $U_{\star}$. See Figure 4 (primitive case) and Figure 5 (satellite case). Note that $f_{c}: U_{k} \rightarrow V_{0}$ is conformal.

If $c$ is a co-root of some hyperbolic component, then exactly one period parameter ray $R_{M_{d}}(\beta)$ land on it (see Figure 6). In dynamical plane, $R_{c}(\beta)$ is the unique dynamical ray landing on a parabolic periodic point $\gamma_{c}(\beta):=x_{1}$, whose immediate basin contains the critical value $c$. The parameter $c$ is a primitive parabolic parameter. Denote $V_{1}$ the union of Fatou component containing $c$ and external ray $R_{c}(\beta), V_{0}=\mathbb{C} \backslash \overline{V_{1}}, U_{\star}=f_{c}^{-1}\left(V_{1}\right) . U_{k}$ is the component of $f_{c}^{-1}\left(V_{0}\right)$ adjacent with $R_{c}((\beta+k-1) / d)$ and $R_{c}((\beta+k) / d), k \in \mathbb{Z}_{d}$. (see Figure 7).

Remark: in our paper, if $c$ is a parabolic parameter, then $f_{c}$ has unique parabolic orbit, denoted by $\left\{x_{0}, x_{1}, \ldots, x_{p-1}\right\}$. $x_{1}$ is the point whose immediate basin contains critical value $c$.

The following lemma provides a criterion for $\theta$ such that $\gamma_{M_{d}}(\theta)$ is a primitive parabolic parameter.
Definition 3.10. Let $\theta$ be a periodic angle of period $n$ and the d-expansion of $\theta$ be.$\overline{\epsilon_{1} \ldots \epsilon_{n}}$. We call $\epsilon_{1} \ldots \epsilon_{n}$ the periodic part of the $d$-expansion of $\theta$.


Figure 5: The dynamical plane of $f_{c_{1}} \cdot c_{1}:=\gamma_{M_{3}}(11 / 80)=\gamma_{M_{3}}(19 / 80)$ is the root of some satellite hyperbolic component as illustrated in Figure 3. The dynamical rays $R_{c_{1}}(11 / 80)$ and $R_{c_{1}}(19 / 80)$ land on a common parabolic point of $f_{c_{1}}$ with period 2.


Figure 6: Multibrot set $M_{4}$. The parameter rays $R_{M_{4}}(1 / 15)$ and $R_{M_{4}}(4 / 15)$ land on the root of some hyperbolic component. $R_{M_{4}}(2 / 15)$ and $R_{M_{4}}(1 / 5)$ land on two co-root of this hyperbolic component respectively.

Lemma 3.11. $\theta$ is periodic under $\tau$ with period $n \geq 2$. If $c_{0}:=\gamma_{M_{d}}(\theta)$ is the root of some satellite hyperbolic component, then $\theta$ satisfies the following properties:
(1) $\nu(\theta)$ is cyclic.
(2) Denote by $\overline{w^{s-1} w_{\star}}$ the cyclic expression of $\nu(\theta)$ where $w=\nu_{1} \ldots \nu_{t}, t$ is a proper factor of $n$ and $t s=n$. Then the last digit of the period part of the $d$-expansion of $\theta$ is $\nu_{t}$ or $\nu_{t}-1$.

Moreover, if $\theta$ is maximal in its orbit, then $\nu(\theta)$ also satisfies
(3) $t$ is the length of parabolic orbit and the last digit of the period part of the d-expansion of $\theta$ must be $\nu_{t}-1 \in[0, d-2]$.

Proof. Let $\eta$ be the companion angle of $\theta$, then in dynamical plane of $f_{c_{0}}, R_{c_{0}}(\theta)$ and $R_{c_{0}}(\eta)$ land on $x_{1}$ (see Figure 5). As $V_{1}$ contains no points and external rays of the parabolic orbit, then $\left\{x_{0}, x_{1}, \ldots, x_{p-1}\right\}$ together with their external rays belong to $\bigcup_{k=0}^{d-1} \bar{U}_{k}$.

For $c_{0}$ is satellite parabolic parameter, the length $p$ of parabolic orbit is a proper factor of $n$ and $f_{c_{0}}$ acts on the rays of the orbit transitively. Then we have, in $\nu(\theta)=\overline{\nu_{1} \ldots \nu_{n-1} \star}$, $\nu_{j}=\nu_{j(\bmod ) p}$ for $1 \leq j \leq n-1$, that is, $\nu(\theta)=\overline{u^{l-1} u_{\star}}$ where $u=\nu_{1} \ldots \nu_{p}$. By definition of kneading sequence, we can see $\tau^{\circ(p-1)}(\theta) \in\left(\left(\theta+\nu_{p}-1\right) / d,\left(\theta+\nu_{p}\right) / d\right)$. It follows $x_{0}$


Figure 7: The dynamical plane of $f_{c_{0}} . \quad c_{0}:=\gamma_{M_{4}}(1 / 5)$ is a co-root of the hyperbolic component illustrated in Figure 6. $R_{c_{0}}(1 / 5)$ is the unique dynamical ray landing on $\gamma_{c_{0}}(1 / 5)$ which is the parabolic point of $f_{c_{0}}$ with period 2 .
together with its external rays belong to $\bar{U}_{\nu_{p}}$. Then $\tau^{n-1}(\theta)$ is either $\left(\theta+\nu_{p}-1\right) / d(\theta>\eta)$ or $\left(\theta+\nu_{p}\right) / d(\theta<\eta)$ (see Figure 8). So the last digit of $d$-expansion of $\theta$ is either $\nu_{p}-1(\theta>\eta)$ or $\nu_{p}(\theta<\eta)$. Let $w=\nu_{1} \ldots \nu_{t}$ be the primitive root of $u$, then $u=w^{p / t}$. We have $\overline{w^{s-1} w_{\star}}$ is the cyclic expression of $\nu(\theta)$ (proposition 3.8) and $\nu_{t}=\nu_{p}$, so $\theta$ satisfies property (1) and (2).


Figure 8:
Furthermore, if $\theta$ is maximal in its orbit, then $\theta>\eta$, so the last digit of the period part of the $d$-expansion of $\theta$ must be $\nu_{t}-1$. By lemma $3.2, \theta=\overline{w^{s-1} \nu_{1} \ldots \nu_{t-1}\left(\nu_{t}-1\right)}$ and $0 \leq \nu_{t}-1 \leq d-2$. Note that the angles of external rays belonging to $x_{1}$ are $\theta, \tau^{p}(\theta), \ldots, \tau^{(s-1) p}(\theta)$ with the order $\theta>\tau^{p}(\theta)>\cdots>\tau^{(s-1) p}(\theta)$. The maximum of $\theta$ implies $\eta$ is the second largest angle in orbit of $\theta$, then $\eta=\tau^{p}(\theta)=\overline{u^{l-2} \nu_{1} \ldots \nu_{p-1}\left(\nu_{p}-1\right) u}$. If $u$ is not primitive, then $p / t>1$. It follows $\tau^{t}(\theta)>\tau^{p}(\theta)=\eta$, a contradiction to that $\eta$ is the second largest angle in orbit of $\theta$. So $u$ is a primitive word and hence $t=p$ is length of parabolic orbit.

Then once $\theta$ doesn't satisfy the property in this lemma, we have $\gamma_{M_{d}}(\theta)$ is a primitive parabolic parameter. The lemma below can be seen as a application of lemma 3.11.
Lemma 3.12. Assume $\theta=\overline{w^{s-1} \nu_{1} \ldots \nu_{t-1}\left(\nu_{t}-1\right)}$ is maximal in its orbit, where $w=$ $\nu_{1} \ldots \nu_{t}$ is primitive with $\nu_{t} \in[1, d-1]$ and $t$ is a proper factor of $n$ with $t s=n$. Let

$$
\begin{aligned}
& \beta_{v_{t}-i}=\overline{w^{s-1} \nu_{1} \ldots \nu_{t-1}\left(\nu_{t}-i\right)} \text { for } 2 \leq i \leq \nu_{t} \\
& \beta_{-1}= \begin{cases}\overline{w^{s-1} \nu_{1} \ldots\left(\nu_{t-1}-1\right)(d-1)} & \text { as } t \geq 2 \\
. \overline{k \ldots k(k-1)(d-1)} & \text { as } t=1\end{cases}
\end{aligned}
$$

Then $\gamma_{M_{d}}\left(\beta_{\nu_{t}-i}\right)$ is a primitive parabolic parameter for any $2 \leq i \leq \nu_{t} . \gamma_{M_{d}}\left(\beta_{-1}\right)$ is a satellite parabolic parameter for $\theta=\overline{(d-1) \cdots(d-1)(d-2)}$ and a primitive parabolic parameter for any other case.

Proof. Let $\beta=\overline{w^{s-1} \nu_{1} \ldots \nu_{t-1} j}$ be any angle among $\left\{\beta_{\nu_{t}-i}\right\}_{2 \leq i \leq \nu_{t}}$, then $0 \leq j \leq \nu_{t}-2$. The maximum of $\theta$ implies the maximum of $\beta$ in its orbit. Since $w$ is primitive, by lemma 3.2, we have $\overline{w^{s-1} w_{\star}}$ is the cyclic expression of $\nu(\beta)$. As $j \leq \nu_{t}-2<\nu_{t}-1$, with the maximum of $\beta$, the property (3) in lemma 3.11 is not satisfied. So $\gamma_{M_{d}}(\beta)$ is a primitive parabolic parameter.

For $\beta_{-1}$, the maximum of $\theta$ implies $\beta_{-1}$ is greater than $\tau\left(\beta_{-1}\right), \tau^{2}\left(\beta_{-1}\right), \ldots, \tau^{n-2}\left(\beta_{-1}\right)$ but less than $\tau^{n-1}\left(\beta_{-1}\right)$. It follows $\nu(\beta)=\left\{\begin{array}{ll}\overline{w^{s-1} \nu_{1} \ldots \nu_{t-1} \star} & \text { as } t \geq 2 \\ \overline{k \ldots k k \star} & \text { as } t=1\end{array} \overline{w^{s-1} w_{\star}}\right.$. It is the cyclic expression of $\nu(\beta)$, then if $\beta$ satisfies the property in lemma 3.11, $\nu_{t}$ is either 0 or $d-1$. Since $1 \leq \nu_{t} \leq d-1$, we have $\nu_{t}$ must be $d-1$, then the maximum of $\theta$ implies $\theta=. \overline{(d-1) \cdots(d-1)(d-2)}$. So $\gamma_{M_{d}}\left(\beta_{-1}\right)$ is a primitive parabolic parameter as long as $\theta \neq . \overline{(d-1) \cdots(d-1)(d-2)}$. In the case of $\theta=. \overline{(d-1) \cdots(d-1)(d-2)}$, we will see in lemma 3.14 that $\gamma_{M_{d}}(\theta)$ is the root of a hyperbolic component attached to the main cardioid and $\beta_{-1}$ is the companion angle of $\theta$. In this case, $\gamma_{M_{d}}\left(\beta_{-1}\right)$ is a satellite parabolic parameter.

Remark. In this lemma, we distinguish $\beta_{-1}$ according to whether $t \geq 2$ or $t=1$. It is because that we don't find a uniform expression of $\beta_{-1}$ for the two cases rather than the case of $t=1$ is special.

### 3.4 Itineraries outside the Multibrot set

If $c \in \mathbb{C} \backslash M_{d}$, the Julia set of $f_{c}$ is a Cantor set. If $c \in R_{M_{d}}(\theta)$ with $\theta \neq 0$ not necessarily periodic, then the dynamical rays $R_{c}(\theta / d) \ldots R_{c}((\theta+d-1) / d)$ bifurcate on the critical point. The set $R_{c}(\theta / d) \cup \ldots \cup R_{c}((\theta+d-1) / d) \cup\{0\}$ separates the complex plane in $d$ connected components. We denote by $U_{0}$ the component containing the dynamical ray $R_{c}(0)$ and by $U_{1}, \ldots, U_{d-1}$ the other component in counterclockwise (see Figure 9).

The orbit of a point $x \in K_{c}$ has an itinerary with respect to this partition. In other words, to each $x \in K_{c}$, we can associate a sequence $\iota_{c}(x) \in\{0,1, \ldots d-1\}{ }^{\mathbb{N}}$ whose $j$-th term is equal to $k$ if $f_{c}^{\circ j-1}(x) \in U_{k}$. A point $x \in K_{c}$ is periodic for $f_{c}$ if and only if the itinerary $\iota_{c}(x)$ is periodic for the shift with the same period.

The map $\iota_{c}: K_{c} \rightarrow\{0,1, \ldots d-1\}^{\mathbb{N}}$ is a bijection. In particular, for each itinerary $\iota \in\{0, \ldots, d-1\}^{\mathbb{N}}$ and each $c \in \mathbb{C} \backslash\left(M_{d} \cup R_{M_{d}}(0)\right)$, there is a unique point $x(\iota, c) \in K_{c}$ whose itinerary is $\iota$. For a given $\iota \in\{0, \ldots, d-1\}^{\mathbb{N}}$, the map $\mathbb{C} \backslash\left(M_{d} \cup R_{M_{d}}(0)\right) \longrightarrow$ $\mathbb{C} \quad c \mapsto x(\iota, c) \in \mathbb{C}$ is continuous, and even holomorphic (as can be seen by applying the Implicit Function Theorem).
Proposition 3.13. Let $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1 \star}}$ be the kneading sequence of a periodic angle $\theta$ with period $n \geq 2$. If $c_{0}:=\gamma_{M_{d}}(\theta)$ is a primitive parabolic parameter and if one follows continuously the periodic points of period $n$ of $f_{c}$ as c makes a small turn around $c_{0}$, then


Figure 9: The regions $U_{0}, U_{1}, U_{2}, U_{3}$ for a parameter $c$ belonging to $R_{M_{4}}(1 / 15)$.
the periodic points with itineraries $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1} k}$ and $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1}(k+1)}$ get exchanged where $k \in \mathbb{Z}_{d}$ is the last digit of the period part of the d-expansion of $\theta$.

Proof. Since $c_{0}$ is a primitive parabolic parameter, then the periodic point $x_{1}:=\gamma_{c_{0}}(\theta)$ has period $n$ and multiplier 1. According to Case 2 in the proof of smoothness and lemma 3.9, the projection from a small neighborhood of $\left(c_{0}, x_{1}\right)$ in $X_{n}$ to the first coordinate is a degree 2 covering. So the neighborhood of $\left(c_{0}, x_{1}\right)$ in $\overline{X_{n}}$ can be written as

$$
\left\{\left(c_{0}+\delta^{2}, x(\delta)\right),\left(c_{0}+\delta^{2}, x(-\delta)\right)| | \delta \mid<\varepsilon\right\}
$$

where $x:(\mathbb{C}, 0) \rightarrow\left(\mathbb{C}, x_{1}\right)$ is a holomorphic germ with $x^{\prime}(0) \neq 0$. In particular, the pair of periodic points for $f_{c}$ which are splitted from $x_{1}$ get exchanged when $c$ makes a small turn around $c_{0}$. So, using analytic continuation on $\mathbb{C} \backslash\left(M_{d} \cup R_{M_{d}}(0)\right)$, it is enough to show that there exists a $c \in \mathbb{C} \backslash M_{d}$ close to $c_{0}$ such that $x\left( \pm \sqrt{c-c_{0}}\right)$ have itineraries $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1} k}$ and $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1}(k+1)}$ where $k \in \mathbb{Z}_{d}$ is the last digit of the period part of the $d$-expansion of $\theta$.

Let us denote by $V_{0}\left(c_{0}\right), V_{1}\left(c_{0}\right), U_{0}\left(c_{0}\right), \ldots, U_{d-1}\left(c_{0}\right)$ and $U_{\star}\left(c_{0}\right)$ the sets defined in the previous section. For $j \geq 0$, set $x_{j}:=f_{c_{0}}^{j}\left(x_{0}\right)$ and observe that for $j \in[1, n-1]$, we have $x_{j} \in U_{\varepsilon_{j}}\left(c_{0}\right)$.

For $c \in R_{M_{d}}(\theta)$, consider the following compact subsets of the Riemann sphere :

$$
R(c):=R_{c}(\theta) \cup\{c, \infty\} \quad \text { and } \quad S(c):=R_{c}(\theta / d) \cup \ldots \cup R_{c}((\theta+d-1) / d) \cup\{0, \infty\} .
$$

Denote by $U_{0}(c)$ the component of $\mathbb{C} \backslash S(c)$ containing $R_{c}(0)$ and by $U_{1}(c), \ldots, U_{d-1}(c)$ the other component in counterclockwise. From any sequence $\left\{c_{m}\right\} \subset R_{M_{d}}(\theta)$ converging to $c_{0}$, by extracting a subsequence if necessary, we can assume $R\left(c_{m}\right)$ and $S\left(c_{m}\right)$ converge respectively, for the Hausdorff topology on compact subsets of $\mathbb{C} \cup\{\infty\}$, to connected compact sets $R$ and $S$. Since $S(c)=f_{c}^{-1}(R(c))$, we have $S=f_{c_{0}}^{-1}(R)$. According to [PR, Section 2 and 3], $R \cap\left(\mathbb{C} \backslash K_{c_{0}}\right)=R_{c_{0}}(\theta)$, the intersection of $R$ with the boundary of $K_{c_{0}}$ is reduced to $\left\{x_{1}\right\}$ and the intersection of $R$ with the interior of $K_{c_{0}}$ is contained in the immediate basin of $x_{1}$, whence in $V_{1}$. It follows $R \subset \overline{V_{1}}\left(c_{0}\right)$ and $S \subset \bar{U}_{\star}\left(c_{0}\right)$, that means any compact subset of $\mathbb{C} \backslash \bar{U}_{\star}\left(c_{0}\right)$ is contained in $\mathbb{C} \backslash S\left(c_{m}\right)$ for $m$ sufficiently large .

For $j \in[1, n-1]$ and let $D_{j}$ be a sufficiently small disk around $x_{j}$ so that

$$
\bar{D}_{j} \subset U_{\varepsilon_{j}}\left(c_{0}\right) \subset \mathbb{C} \backslash \bar{U}_{\star}\left(c_{0}\right) .
$$

According to the previous discussion, if $m$ is sufficiently large, we have

$$
f_{c_{m}}^{j-1}\left(x\left( \pm \sqrt{c_{m}-c_{0}}\right)\right) \subset D_{j} \subset U_{\varepsilon_{j}}\left(c_{m}\right)
$$

So the first $n-1$ symbols of the itineraries of $x\left( \pm \sqrt{c_{m}-c_{0}}\right)$ are all $\epsilon_{1}, \ldots, \epsilon_{n-1}$. As $x\left(\sqrt{c_{m}-c_{0}}\right)$ and $x\left(-\sqrt{c_{m}-c_{0}}\right)$ are different $n$ periodic points of $f_{c_{m}}$, their itineraries must be different. It follows $f_{c_{m}}^{n-1}\left(x\left( \pm \sqrt{c_{m}-c_{0}}\right)\right)$, which are splitted from $x_{0}$, lie in different component of $\mathbb{C} \backslash S\left(c_{m}\right)$. Combining with the fact that $R_{c_{0}}((\theta+k) / d)$ lands on $x_{0}(k$ is the last digit of the period part of the $d$-expansion of $\theta$ ), we have $f_{c_{m}}^{n-1}\left(x\left( \pm \sqrt{c_{m}-c_{0}}\right)\right)$ belong to $U_{k}\left(c_{m}\right)$ and $U_{k+1}\left(c_{m}\right)$ respectively, then $x\left( \pm \sqrt{c_{m}-c_{0}}\right)$ have itineraries $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1} k}$ and $\overline{\varepsilon_{1} \ldots \varepsilon_{n-1}(k+1)}$ respectively.

Lemma 3.14. For $\theta=1-1 /\left(d^{n}-1\right)=. \overline{(d-1) \cdots(d-1)(d-2)}(n \geq 2)$, we have $\gamma_{M_{d}}(\theta)$ is the root of some periodic n hyperbolic component attached to the main cardioid. If $\eta$ is denoted the companion angle of $\theta$, then $\eta=d \theta-d+1$.

Proof. Let $c_{0}:=\gamma_{M_{d}}(\theta)$, then $x_{1}:=\gamma_{c_{0}}(\theta)$ is the parabolic periodic point of $f_{c_{0}}$ as previous. By lemma 3.2, $\nu(\theta)=\overline{(d-1) \cdots(d-1) \star}$, so $\overline{(d-1) \cdots(d-1) \star}$ is the cyclic expression of $\nu(\theta)$. If $x_{0} \neq x_{1}$, then the length of parabolic orbit is greater than 1 . It implies the property (3) in lemma 3.11 is not satisfied, so $c_{0}$ is a primitive parabolic parameter. According to proposition3.13, when $c \in \mathbb{C} \backslash M_{d}$ is close to $c_{0}, x_{1}$ splits into two $n$ periodic point $y, z$ of $f_{c}$ with itineraries $\overline{(d-1) \cdots(d-1)(d-2)}$ and $\overline{(d-1) \cdots(d-1)(d-1)}$. It leads to a contradiction to the period $n$ of $y$ and $z$. So $x_{0}=x_{1}$ and then $c_{0}$ is the root of some periodic $n$ satellite hyperbolic component attached to the main cardioid.

By the maximum of $\theta$, we have $U_{d-1}$ is bounded by $R_{c_{0}}((\theta+d-2) / d)$ and $R_{c_{0}}((\eta+$ $d-1) / d) . \nu(\theta)=\overline{(d-1) \cdots(d-1) \star}$ implies $R_{c_{0}}(\theta) \subset \bar{U}_{d-1}$, then $\theta \leq(\eta+d-1) / d$ and $x_{0}$ is on the boundary of $U_{d-1}$. On the other hand, $(\eta+d-1) / d$ is in the orbit of $\theta$, so $\theta \geq(\eta+d-1) / d$. Then we have $\eta=d \theta-d+1$.

Remark. The dynamical rays $R_{c_{0}}(\theta)$ and $R_{c_{0}}(\eta)$ are consecutive among the rays landing at $x_{0}$. Lemma 3.14 implies $R_{c_{0}}(\theta)$ is mapped to $R_{c_{0}}(\eta)$. It follows that each dynamical ray landing at $x_{0}$ is mapped to the one which is once further clockwise.

Proposition 3.15. Let $\theta=1-1 /\left(d^{n}-1\right)=. \overline{(d-1) \cdots(d-1)(d-2)}$ be periodic with period $n \geq 2$. If one follows continuously the periodic points of period $n$ of $f_{c}$ as $c$ makes a small turn around $\gamma_{M_{d}}(\theta)$, then the periodic points in the cycle of $\iota_{c}^{-1}(\overline{(d-1) \cdots(d-1)(d-2)})$ get permuted cyclically.


Figure 10: The dynamical plane of $f_{c_{0}} . c_{0}:=\gamma_{M_{3}}(\theta)$ with $\theta=. \overline{2221}$

Proof. Set $c_{0}:=\gamma_{M_{d}}(\theta)$. By Lemma 3.14, all the dynamical rays $R_{c_{0}}\left(\tau^{j}(\theta)\right)$ land on a common fixed point $x_{0}$. This fixed point is parabolic and the companion angle of $\theta$, denoted by $\eta$, equals to $d \theta-(d-1) \equiv d \theta(\bmod \mathbb{Z})$. $V_{1}\left(c_{0}\right) \subset U_{d-1}\left(c_{0}\right)$ which is bounded by $R_{c_{0}}((\theta+d-2) / d)$ and $R_{c_{0}}(\theta)$.

According to Case 3 in the proof of smoothness and lemma 3.9, we have the projection from a small neighborhood of $\left(c_{0}, x_{0}\right)$ in $X_{n}$ to the parameter plane is a degree $n$ covering. Then the neighborhood of $\left(c_{0}, x_{0}\right)$ in $\overline{X_{n}}$ can be written as

$$
\left\{\left(c_{0}+\delta^{n}, x(\delta)\right),\left(c_{0}+\delta^{n}, x(\omega \delta)\right), \ldots,\left(c_{0}+\delta^{n}, x\left(\omega^{n-1} \delta\right)\right)| | \delta \mid<\varepsilon\right\}
$$

where $x:(\mathbb{C}, 0) \rightarrow\left(\mathbb{C}, x_{0}\right)$ is a holomorphic germ satisfying $x^{\prime}(0) \neq 0$. So, for $c$ close to $c_{0}$, the set $\left.x\left\{\sqrt[n]{c-c_{0}}\right)\right\}$ is a cycle of period $n$ of $f_{c}$, and when $c$ makes a small turn around $c_{0}$, the periodic points in the cycle $\left.x\left\{\sqrt[n]{c-c_{0}}\right)\right\}$ get permuted cyclically. So, combining with analytic continuation on $\mathbb{C} \backslash\left(M_{d} \cup R_{M_{d}}(0)\right)$, it is enough to show there exists a $c \in \mathbb{C} \backslash M_{d}$ close enough to $c_{0}$ such that the point $\iota_{c}^{-1}(\overline{(d-1) \cdots(d-1)(d-2)})$ belongs to $x\left\{\sqrt[n]{c-c_{0}}\right\}$. Equivalently, we must show that there is a sequence $\left\{c_{j}\right\} \subset$ $\mathbb{C} \backslash M_{d}$ converging to $c_{0}$, such that the periodic point $y_{j}:=\iota_{c_{j}}^{-1}(\overline{(d-1) \cdots(d-1)(d-2)})$ converges to $x_{0}$.

Let $\left\{c_{j}\right\} \subset R_{M_{d}}(\theta)$ converge to $c_{0}$ as $j \rightarrow \infty$. Without loss of generality, we may assume that the sequence $y_{j}$ converges to a point $z, R\left(c_{j}\right)$ converges to $R$ and $S\left(c_{j}\right)$ converges to $S$ in Hausdoff topology. The definition of $R(c), S(c), U_{0}(c), \ldots, U_{d-1}(c)$ are in the proof of proposition 3.13. As $\left(c_{0}, z\right)$ is on $\overline{X_{n}}$, then $z$ is either the parabolic fixed point or repelling $n$ periodic point of $f_{c_{0}}$.

Suppose $z$ is a repelling $n$ periodic point, set $z_{i}:=f_{c_{0}}^{i}(z)$. Now we will define a new sequence of open domain $\left\{W_{k}\left(c_{0}\right)\right\} . W_{k}\left(c_{0}\right)$ is the connected component of $U_{\star}\left(c_{0}\right) \backslash$ the closure of Fatou component containing 0 , adjacent with $U_{k}\left(c_{0}\right), U_{k+1}\left(c_{0}\right)$ (see Figure 10). According to [PR, Section 2 and 3], $R \cap\left(\mathbb{C} \backslash K_{c_{0}}\right)=R_{c_{0}}(\theta)$, the intersection of $R$ with the boundary of $K_{c_{0}}$ is reduced to $\left\{x_{0}\right\}$ and the intersection of $R$ with the interior of $K_{c_{0}}$ is contained in the immediate basin of $x_{0}$. It follows $\left\{z_{0}, \ldots, z_{n-1}\right\} \bigcap S=\emptyset$. Then for $j$ sufficiently large, $\left\{z_{0}, \ldots, z_{n-1}\right\} \subset \mathbb{C} \backslash S_{c_{j}}$. As $y_{j}$ has itineraries $\frac{(d-1) \cdots(d-1)(d-2)}{}$, we have $\left\{z_{0}, \ldots z_{n-2}\right\} \subset U_{d-1}\left(c_{0}\right) \bigcup \bar{W}_{d-1}\left(c_{0}\right), z_{n-1} \in U_{d-2}\left(c_{0}\right) \bigcup \bar{W}_{d-2}\left(c_{0}\right)$.
Claim 1. $z_{n-1} \notin \bar{W}_{d-2}\left(c_{0}\right)$.
Proof. In $J\left(f_{c_{0}}\right), x_{0}$ is the unique periodic point with more than one external rays landing on it (refer to [Poi, proposition 3.3]). So there is exactly one external ray landing on $z_{n-1}$ with period $n$. Its angle is denoted by $\frac{a}{d^{n}-1}, a$ is a integer. If $z_{n-1} \in \bar{W}_{d-2}$, the angle of external ray belonging to $z_{n-1}$ satisfy

$$
\frac{\eta+d-2}{d}<\frac{a}{d^{n}-1}<\frac{\theta+d-2}{d}\left(\theta=1-\frac{1}{d^{n}-1}, \eta=d \theta-d+1\right)
$$

by simple computation, we have

$$
\frac{k\left(d^{n}-1\right)}{d-1}-d^{n-1}-1+\frac{1}{d}<a<\frac{k\left(d^{n}-1\right)}{d-1}-d^{n-1}
$$

a contradiction to $a$ is an integer. This ends the proof of claim 1.
Claim 2. $z_{n-1} \notin U_{d-2}\left(c_{0}\right)$.
Proof. If $z_{n-1} \in U_{d-2}\left(c_{0}\right)$, we label the sectors at $x_{0}$ by $S_{i}(0 \leq i \leq n-1)$ clockwise with $S_{0}=V_{1}\left(c_{0}\right)$. The dynamics between these sectors satisfy

$$
V_{1}\left(c_{0}\right)=S_{0} \xrightarrow{f_{c_{0}}} S_{1} \xrightarrow{f_{c_{0}}} \cdots \xrightarrow{f_{c_{0}}} S_{n-2} \xrightarrow{f_{c_{0}}} S_{n-1}=\mathbb{C} \backslash \bar{U}_{d-1}\left(c_{0}\right)
$$

As $\left\{z_{0}, \ldots z_{n-2}\right\} \subset U_{d-1}\left(c_{0}\right) \bigcup \bar{W}_{d-1}\left(c_{0}\right)$, we have $z_{0}=f_{c_{0}}\left(z_{n-1}\right)$ belongs to the union of $\overline{W_{d-1}}\left(c_{0}\right)$ and $\bigcup_{i=1}^{n-2} S_{i}$. If $z_{0} \in S_{i_{0}}\left(1 \leq i_{0} \leq n-2\right)$, then $f_{c_{0}}^{\left(n-2-i_{0}\right)}\left(z_{0}\right)=z_{n-2-i_{0}} \in$
$f_{c_{0}}^{\left(n-2-i_{0}\right)}\left(S_{i_{0}}\right)=S_{n-2}$. It follows $f_{c_{0}}\left(z_{n-2-i_{0}}\right)=z_{n-1-i_{0}}$ must belong to $\bar{W}_{d-1}\left(c_{0}\right)$. So $z_{n-i_{0}} \in S_{0}$ and $f_{c_{0}}^{\left(i_{0}-1\right)}\left(z_{n-i_{0}}\right)=z_{n-1} \in f_{c_{0}}^{\left(i_{0}-1\right)}\left(S_{0}\right)=S_{i_{0}-1}$, contradiction to $z_{n-1} \in U_{d-2}$. If $z_{0} \in \bar{W}_{d-1}\left(c_{0}\right)$, then $z_{1} \in S_{0}$. We have $f_{c_{0}}^{(n-2)}\left(z_{1}\right)=z_{n-1} \in f_{c_{0}}^{(n-2)}\left(S_{0}\right)=S_{n-2}$, also a contradiction to $z_{n-1} \in U_{d-2}\left(c_{0}\right)$. This ends the proof of claim 2.

The two claim imply the assumption that $z$ is repelling $n$ periodic point is false and then $z$ must be a parabolic fixed point of $f_{c_{0}}$, that is $z=x_{0}$.

### 3.5 Proof of Theorem 3.1

Fix $n>1$ (the case $n=1$ has been treated directly at the beginning). We proceed to show that $X_{n}$ is connected.

Set $X:=\mathbb{C} \backslash\left(M_{d} \cup R_{M_{d}}(0)\right)$ and $F_{n}:=\mathbb{C} \backslash$ all the landing points of periodic $n$ parameter rays. Take any pair of points $(a, w),\left(a^{\prime}, w^{\prime}\right)$ in $X_{n}$. By analytic continuation, we may assume $a, a^{\prime} \in X$. Again by analytic continuation on simply connected open set $X$, we may assume $a=a^{\prime}$. Thus it is enough to show that there exists a loop in $F_{n}$ based on $a$ such that the analytic continuation along the loop connects $w$ and $w^{\prime}$. We will give a algorithm to find such a loop.

Let z be any $n$ periodic point of $f_{a}$.
step 1 In the orbit of $z$, there is a point with maximal itineraries among the shift of $\iota_{a}(z)$ in the lexicograph order, denoted by $\overline{\epsilon_{1} \ldots \epsilon_{n}}$. Set $\theta=. \overline{\epsilon_{1} \ldots \epsilon_{n}}(\theta$ is maximal in its orbit). If $\theta$ satisfies the properties in lemma 3.11, do step 2 below. Otherwise, $\gamma_{M_{d}}(\theta)$ is a primitive parabolic parameter. According to lemma 3.2 and proposition 3.13, when $a$ makes a turn around $\gamma_{M_{d}}(\theta)$, the periodic point of $f_{a}$ with itineraries $\overline{\epsilon_{1} \ldots \epsilon_{n}}$ and $\overline{\epsilon_{1} \ldots\left(\epsilon_{n}+1\right)}$ get changed. Then $z$ is connected to a new orbit containing $\iota_{a}^{-1}\left(\epsilon_{1} \ldots\left(\epsilon_{n}+1\right)\right)$. For this new orbit, repeat doing step 1 .
step $2 \theta=. \overline{\epsilon_{1} \ldots \epsilon_{n}}$ is maximal in its orbit and satisfies the properties in lemma 3.11. If $\theta=\overline{(d-1) \cdots(d-1)(d-2)}$, step 2 ends. Otherwise, let $\overline{w^{s-1} w_{\star}}$ be the cyclic expression of $\nu(\theta)$ where $w=\nu_{1} \ldots \nu_{t}, \nu_{t} \in[1, d-1]$. As in lemma 3.12, we obtain a sequence of angles $\left\{\beta_{\nu_{t}-2}, \ldots, \beta_{0}, \beta_{-1}\right\}$ and know that $\gamma_{M_{d}}\left(\beta_{\nu_{t}-i}\right)$ is a primitive parabolic parameter with $\nu(\theta)=\overline{\epsilon_{1} \ldots \epsilon_{n-1 \star}}$ for any $i \in\left[2, \nu_{t}+1\right]$. Then by proposition 3.13 again, as $a$ makes a turn around $\gamma_{M_{d}}\left(\beta_{\nu_{t}-i}\right)\left(2 \leq i \leq \nu_{t}+1\right)$, the periodic points of $f_{a}$ with itineraries $\overline{\epsilon_{1} \ldots \epsilon_{n-1}\left(\nu_{t}-i\right)}$ and $\overline{\epsilon_{1} \ldots \epsilon_{n-1}\left(\nu_{t}-i+1\right)}$ get changed. Then let $a$ makes turns around from $\gamma_{M_{d}}\left(\beta_{\nu_{t}-2}\right)$ to $\gamma_{M_{d}}\left(\beta_{-1}\right)$ one by one, we have $\iota_{a}^{-1}\left(\overline{\epsilon_{1} \ldots \epsilon_{n-1} \epsilon_{n}}\right)$ are connected with $\iota_{a}^{-1}\left(\overline{\epsilon_{1} \ldots \epsilon_{n-1}(d-1)}\right)$ by analytic continuation through the points $\iota_{a}^{-1}\left(\overline{\epsilon_{1} \ldots \epsilon_{n-1}\left(\epsilon_{n}-1\right)}\right), \ldots, \iota_{a}^{-1}\left(\overline{\epsilon_{1} \ldots \epsilon_{n-1} 0}\right)$. For the new periodic point $\iota_{a}^{-1}\left(\overline{\epsilon_{1} \ldots \epsilon_{n-1}(d-1)}\right)$, do step 1 .

Every time a $n$ periodic point of $f_{a}$ passes though step 1 or step 2 , the sum of all digits in the itineraries of the output periodic point is greater than that of the input one. For fixed $n$, this sum is bounded (the bound is $(d-1) n-1$ ), then each $n$ periodic point $z$ can be connected to the orbit containing $\iota_{a}^{-1}(\overline{(d-1) \cdots(d-1)(d-2)})$.

In our case, applying the procedure above to $w$ and $w^{\prime}$, we have $w$ and $w^{\prime}$ are connected to two points of the periodic orbit containing $\iota_{a}^{-1}(\overline{(d-1) \cdots(d-1)(d-2)})$. Proposition 3.15 tells us, by analytic continuation, any two point in this orbit can be connected as long as $a$ makes the appropriate number of turns around $\gamma_{M_{d}}\left(1-\frac{1}{d^{n}-1}\right)$. Thus $w$ and $w^{\prime}$ are connected.

## References

[B] T. Bousch, Sur quelques problèmes de dynamique holomorphe, Ph.D. thesis, Université de Paris-Sud, Orsay, 1992.
[BT] X. Buff and Tan Lei, The quadratic dynatomic curve are smooth and irreducible. Volume special en l'honneur des 80 ans de Milnor a paraitre.
[DE] Dominik Eberlein: Rational Parameter Rays of Multibrot Sets. Diploma Thesis, Technische Universit at Munchen, in preparation.
[DH] A. Douady and J.H. Hubbard, Etude dynamique des polynômes complexes (Deuxième partie), 1985.
[E] AL Epstein,Infinitesimal Thurston rigidity and the fatou-Shishikura inequality, Stony Brook IMS Preprint 1999-1.
[KM] K.C.Lyndon, M.P.Schutzenberger, On the equation $a^{M}=b^{N} c^{P}$ is a free group, Michigan Math. J, (1962) 289-298
[LS] E. Lau and D. Schleicher, Internal addresses in the Mandelbrot set and irreducibility of polynomials, Stony Brook Preprint 19, 1994.
[L] G. Levin, On explicit connections between dynamical and parameter spaces, Journal d'Analyse Mathematique, 91 (2003), 297-327.
[Mil] J. Milnor, Periodic orbits, extenal rays and the Mantelbrot set: an expository account.
[Mo] P. Morton, On certain algebraic curves related to polynomial maps. Compositio Math. 103 (1996), no. 3, 319-350.
[Poi] Alfredo Poirier, On Post Critically Finite Polynomials Part Two: Hubbard Trees.
[PR] C. L. Petersen and G. Ryd, Convergence of rational rays in parameter- spaces, in 'The Mandelbrotset, Theme and Variations. edited by Tan Lei, London Mathematical Society, Lecture Note Series 274. Cambridge University Press 2000.


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[^1]:    ${ }^{1}$ use the formula $0=f_{c}^{\circ n}(z)-z=f_{c}^{\circ k m+\ell}(z)-z=f_{c}^{\circ \ell}\left(f_{c}^{\circ k m}(z)\right)-z=f_{c}^{\circ \ell}(z)-z$ and the minimality of $m$ to conclude that $m$ divides $n$.

[^2]:    ${ }^{2}$ One can prove that $D:=\left\{(c, z) \mid f_{c}^{\circ n}(z)=z,\left[f_{c}^{n}\right]^{\prime}(z)=1\right\}$ is finite as follows: Denote by $X(c)$, resp. $Y(z)$ the resultant of the two polynomials $f_{c}^{\circ n}(z)-z$ and $\left[f_{c}^{n}\right]^{\prime}(z)-1$ considered as polynomials of $z$, resp. of $c$. Then $X(c)$ is a polynomial of $c$, resp. $Y(z)$ is a polynomial of $z$. The projection of $D$ to each coordinate equals the zeros of $X$, resp. of $Y$. As no point of the form $(0, z),(c, 0)$ is in $D$, we have $X(0) \neq 0 \neq Y(0)$ so $X, Y$ each has finite many roots. As $D \subset\left(X^{-1}(0) \times \mathbb{C}\right) \cap\left(\mathbb{C} \times Y^{-1}(0)\right)$ we know that $D$ is finite.

