# Moduli of diffeomorphisms with homoclinic tangencies 

ホモクリニック接触を持つ微分同相写像の モジュライ（英文）

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## Introduction

This thesis concerns the topological conjugacy problem for diffeomorphisms on a closed manifold $M$. A diffeomorphism $f$ on $M$ is called structurally stable if any diffeomorphism $g$ close to $f$ is topologically conjugate to $f$. The structural stability for diffeomorphisms are well studied by many authors. In particular, R. Mañé (1987) and others proved that, in the $C^{1}$ category, $f$ is structurally stable if and only if $f$ is an Axiom A diffeomorphism with the strong transversality condition. On the other hand, if $f$ has a basic set which has a homoclinic tangency, then it is never structurally stable. So, if $f$ has a homoclinic tangency, then any neighborhood of $f$ in the space of diffeomorphisms contains both diffeomorphisms $g$ which are topologically conjugate and non-conjugate to $f$. Thus, we need topological conjugacy invariants to decide whether a given $g$ is topologically conjugate to $f$ or not.

A modulus $m(f)$ for a diffeomorphism $f$ is a topological conjugacy invariant for $f$, that is, $m(f)=m(g)$ holds for any $g: M \rightarrow M$ which is contained in a certain class of diffeomorphisms on $M$ and topologically conjugate to $f$. The aim of this thesis is to present new moduli for diffeomorphisms of dimensions two and three.

This thesis is organized as follows.
In Chapter 1, we present definitions, notions and concepts needed in this thesis. Besides, we introduce several preceding results on moduli.

In Chapter 2, we study moduli for 2-dimensional diffeomorphisms with cubic homoclinic tangencies (two-sided tangencies of the lowest order) under certain open conditions. Ordinary arguments used in previous studies of conjugacy invariants associated with onesided tangencies do not work in the two-sided case. We present a new method which is applicable to the two-sided case.

In Chapter 3, we investigate moduli of a 3-dimensional diffeomorphism $f$ with a sadldle focus $p$ and a homoclinic quadratic tangency $q$. It is shown there that, for most of such diffeomorphisms, all the eigenvalues of $D f(p)$ are moduli and the restriction of a conjugacy homeomorphism to a local unstable manifold is a uniquely determined linear conformal map.

## Chapter 1

## Basic definitions and concepts

In this chapter, we present some of definitions, notions and concepts needed in this thesis. Refer to [De, Ro1, Ro2] and so on for other standard results on dynamical systems.

### 1.1 Hyperbolic fixed points of diffeomorphisms

Let $M$ be a $C^{r}(1 \leq r \leq \infty)$ manifold and $\operatorname{Diff}^{r}(M)$ the space of $C^{r}$ diffeomorphisms on $M$ with $C^{r}$ topology. Suppose that $f$ is an element of $\operatorname{Diff}^{r}(M)$. For a point $x \in M$, the orbit $\mathcal{O}(x)$ of $x$ for $f$ is defined as $\mathcal{O}(x)=\left\{f^{n}(x) ; n \in \mathbb{Z}\right\}$, where $f^{0}$ is the identity map on $M, f^{n}$ is the composition of $f$ with itself $n$ times if $n>0$ and $f^{n}$ is the composition of $f^{-1}$ with itself $-n$ times if $n<0$. A point $p \in M$ is called a periodic point for $f$ if $p=f^{n}(p)$ holds for some positive integer $n$. The minimum of such an $n$ is called the period of $p$. A point $p \in M$ is called a fixed point for $f$ if $p=f(p)$ holds, that is, a fixed point is a periodic point with period one.

Suppose that $p$ is a fixed point for $f$. Then the derivative $D f(p)$ of $f$ at $p$ is a linear map on the tangent space $T_{p}(M)$ at $p$. By an identification of $T_{p}(M)$ with $\mathbb{R}^{m}$ as vector spaces, one can regard the linear map on $T_{p}(M)$ with that on $\mathbb{R}^{m}$, where $m$ is the dimension of $M$.

Definition 1.1. A fixed point $p$ for $f$ is called hyperbolic if the absolute value $|\lambda|$ of any eigenvalue $\lambda$ of $D f(p)$ is different from one. The hyperbolic fixed point $p$ is called a sink if the absolute value $|\lambda|$ of any eigenvalue $\lambda$ of $D f(p)$ is less than one. The hyperbolic fixed point $p$ is called a source if the absolute value $|\lambda|$ of any eigenvalue $\lambda$ of $D f(p)$ is greater than one. A hyperbolic fixed point which is neither a sink nor a source is said to be a saddle.

Figures 1.1 and 1.2 illustrate hyperbolic fixed points in the case of $\operatorname{dim} M=2$ and $\operatorname{dim} M=3$, respectively, where all the eigenvalues of $D f(p)$ are real.

We also consider the case that some of eigenvalues are non-real. If $D f(p)$ have non-real eigenvalues $r e^{ \pm \sqrt{-1} \theta}$, then $f$ acts on a neighborhood of $p$ as the combination of a rotation and an expansion or contraction. In the case of $\operatorname{dim} M=3$, we have several phase portraits of $f$ near $p$. The hyperbolic fixed point $p$ is a sink if $D f(p)$ has a real eigenvalue $0<\lambda<1$


Figure 1.1: The case of $\operatorname{dim} M=2 . p$ is a sink in (1), a source in (2) and a saddle in (3).


Figure 1.2: The case of $\operatorname{dim} M=3 . p$ is a sink in (1), a source in (2) and a saddle in (3).
and non-real eigenvalues $r e^{ \pm \sqrt{-1} \theta}$ with $r<1$. The hyperbolic fixed point $p$ is a source if $D f(p)$ has a real eigenvalue $\lambda>1$ and non-real eigenvalues $r e^{ \pm \sqrt{-1} \theta}$ with $r>1$. If the hyperbolic fixed point $p$ is neither a sink nor source, then it is called a saddle focus. See Figure 1.3. In Section 3, we study moduli of 3-dimensional diffeomorphisms having saddle foci with a real eigenvalue $0<\lambda<1$ and non-real eigenvalues $r e^{ \pm \sqrt{-1} \theta}$ with $r>1$. See Figure 1.3 (3).

The following linearization theorem is called the Hartman-Grobman theorem. See the Chapter 5 in [Ro1] for the proof.

Theorem 1.2 (Hartman-Grobman Theorem). Let $f: M \rightarrow M$ be a $C^{r}$ diffeomorphism with a hyperbolic fixed point $p$. Then, there exist neighborhoods $U, V$ of $p$ with $U \cup f(U) \subset$ $V$ and a homeomorphism $h: V \rightarrow T_{p}(M)$ with $h(p)=\mathbf{0}$ and such that the following diagram is commutative.


By Theorem 1.2, we can call the linear map $D f(p)$ a linearized map or linearization of $f$ at $p$. Moreover, by Taylor's theorem, we know that the linear map $D f(p)$ approximates


Figure 1.3: The case of $\operatorname{dim} M=3 . p$ is a sink in (1), a source in (2) and a saddle focus in (3).
$f$ near $p$.

### 1.2 Heteroclinic and homoclinic tangencies

Let $f$ be a $C^{r}$ diffeomorphism on $M$ and $p \in M$ a fixed point for $f$. The stable and unstable manifolds $W^{s}(p)$ and $W^{u}(p)$ of $p$ are defined as

$$
\begin{aligned}
& W^{s}(p)=\left\{x \in M ; f^{n}(x) \rightarrow p(n \rightarrow \infty)\right\}, \\
& W^{u}(p)=\left\{x \in M ; f^{-n}(x) \rightarrow p(n \rightarrow \infty)\right\} .
\end{aligned}
$$

Moreover, we define the local stable and local unstable manifolds $W_{\text {loc }}^{s}(p)$ and $W_{\text {loc }}^{u}(p)$ of $p$ as

$$
\begin{aligned}
& W_{\mathrm{loc}}^{s}(p)=\left\{x \in U(p) ; f^{n}(x) \in U(p) \text { for any } n \in \mathbb{N}, \lim _{n \rightarrow \infty} f^{n}(x)=p\right\}, \\
& W_{\mathrm{loc}}^{u}(p)=\left\{x \in U(p) ; f^{-n}(x) \in U(p), \text { for any } n \in \mathbb{N}, \lim _{n \rightarrow \infty} f^{-n}(x)=p\right\},
\end{aligned}
$$

where $U(p)$ is a sufficiently small neighborhood of $p$ in $M$.
The following theorem is called the Stable Manifold Theorem. This theorem shows that the local stable manifold $W_{\text {loc }}^{s}(p)$ and local unstable manifold $W_{\text {loc }}^{u}(p)$ are $C^{r}$ submanifolds of $M$. See the Chapter 5 in [Ro1] for the proof.

Theorem 1.3 (Stable Manifold Theorem). Let $f: M \rightarrow M$ be a diffeomorphism and let $p \in M$ be a saddle fixed point for $f$. Then the local stable manifold $W_{\text {loc }}^{s}(p)$ of $p$ is a $C^{r}$ submanifold of $M$ tangent to the subspace of $T_{p}(M)$ spanned by the eigenvectors with contracting eigenvalues. Similarly, the local unstable manifold $W_{\mathrm{loc}}^{u}(p)$ of $p$ is a $C^{r}$ submanifold of $M$ tangent to the subspace of $T_{p}(M)$ spanned by the eigenvectors with expanding eigenvalues.

We say that the dimension of $W_{\text {loc }}^{s}(p)$ is the stable index of $p$ and denote it by $\operatorname{ind}^{s}(p)$. Then $\operatorname{ind}^{u}(p)=\operatorname{dim} M-\operatorname{ind}^{s}(p)$ is called the unstable index of $p$. For the definitions of stable and unstable manifolds,

$$
W^{s}(p)=\bigcup_{n \geq 1} f^{-n}\left(W_{\mathrm{loc}}^{s}(p)\right), \quad W^{u}(p)=\bigcup_{n \geq 1} f^{n}\left(W_{\mathrm{loc}}^{u}(p)\right) .
$$

This implies that $W^{s}(p)$ and $W^{u}(p)$ are the images of injective $C^{r}$ immersions from $\mathbb{R}^{s}$ and $\mathbb{R}^{u}$ to $M$, respectively, where $s=\operatorname{ind}^{s}(p)$ and $u=\operatorname{ind}^{u}(p)$.

Let $p_{1}$ and $p_{2}$ are two distinct saddle type fixed points of a diffeomorphism $f$ on $M$. A point $q \in M$ is called a heteroclinic point associated with $p_{1}$ and $p_{2}$ if $q \in W^{s}\left(p_{1}\right) \cap W^{u}\left(p_{2}\right)$, i.e., $\lim _{n \rightarrow \infty} f^{n}(q)=p_{1}, \lim _{n \rightarrow \infty} f^{-n}(q)=p_{2}$. We say that the point $q$ is a transverse heteroclinic point if $W^{s}\left(p_{1}\right)$ and $W^{u}\left(p_{2}\right)$ intersect transversely at $q$, namely, $T_{q}(M)=$ $T_{q}\left(W^{s}\left(p_{1}\right)\right) \oplus T_{q}\left(W^{u}\left(p_{2}\right)\right)$ holds. When $q$ is a non-transverse intersection point, $q$ is called a heteroclinic tangency associated with $p_{1}$ and $p_{2}$. See Figure 1.4.


Figure 1.4: $q$ is one of heteroclinic tangencies associated with $p_{1}$ and $p_{2}$.
Let $p$ is a saddle fixed point of a diffeomorphism $f$ on $M$. A point $q \in M$ is called a homoclinic point associated with $p$ if $q \in W^{s}(p) \cap W^{u}(p) \backslash\{p\}$, i.e., $q \neq p$ and $\lim _{n \rightarrow \infty} f^{n}(q)=p$ and $\lim _{n \rightarrow \infty} f^{-n}(q)=p$. We say that the point $q$ is called a transverse homoclinic point if $W^{s}(p)$ and $W^{u}(p)$ intersect transversely at $q$. When $q$ is a non-transverse intersection point, $q$ is called a homoclinic tangency associated with $p$. See Figure 1.5.

Let $f$ be a $C^{r}(n \leq r \leq \infty)$ diffeomorphism with a heteroclinic or homoclinic tangency $q$. We fix a Riemannian metric on $M$ and define the order of tangency as follows. The tangency is of order $n$ if the limit

$$
\lim _{\substack{w \in W_{\text {loc }}^{s}(p) \\ w \rightarrow q}} \frac{d\left(w, W^{u}(p)\right)}{[d(w, q)]^{n}}
$$

exists and is not zero, where $d$ is the distance on $M$ induced from this metric. If $n=2$ (resp. $n=3$ ), then the tangency $q$ is called quadratic (resp. cubic). If $n$ is even, then the tangency
$q$ is said to be one-sided. If $n$ is odd, then the tangency $q$ is two-sided. See Figures 1.5 and 1.6. Homoclinic tangencies have been studied by Newhouse, Palis and Takens and so on since the seventies. For example, see [dM, dMP, dMvS, KS1, KS2, NPT, Ni, Pa, Po, PT].

(1)

(2)

Figure 1.5: The case of $\operatorname{ind}^{s}(p)=\operatorname{ind}^{u}(p)=1 . q$ is a homoclinic quadratic tangency in (1) and a homoclinic cubic tangency in (2).


Figure 1.6: The case of $\operatorname{ind}^{s}(p)=1$ and $\operatorname{ind}^{u}(p)=2 . \quad p$ is a saddle point and $q$ is a homoclinic quadratic tangency associated with $p$.

Now, we define hyperbolic invariant sets for a diffeomotphism $f$. A subset $S$ of $M$ is said to be positively invariant if $f(x) \in S$ for all $x \in S$, i.e., $f(S) \subset S$. On the other hand, a subset $S$ of $M$ is said to be negatively invariant if $f^{-1}(S) \subset S$. Such an $S$ is said to be an invariant set of $f$ if $f(S)=S$. Notice that any periodic orbit and the orbit of a heteroclinic or a homoclinic point are typical examples of invariant sets for $f$. We denote
by $\|\cdot\|_{x}$ the norm on the tangent space $T_{x}(M)$ at $x \in M$ induced from the Riemannian metric on $M$. A closed invariant set $\Lambda$ for $f$ is said to be hyperbolic if it satisfies the following conditions.
(1) At each point $x \in \Lambda$, the tangent space to $M$ splits as the direct sum of subspaces $\mathbb{E}_{x}^{u}$ and $\mathbb{E}_{x}^{s}$, i.e., $T_{x}(M)=\mathbb{E}_{x}^{u} \oplus \mathbb{E}_{x}^{s}$.
(2) The splitting is invariant under the action of the derivative map, i.e., $D f_{x}\left(\mathbb{E}_{x}^{u}\right)=$ $\mathbb{E}_{f(x)}^{u}$ and $D f_{x}\left(\mathbb{E}_{x}^{s}\right)=\mathbb{E}_{f(x)}^{s}$.
(3) There exist $0<\lambda<1$ and $C>0$ independent of $x$ such that, for all $n \geq 0$,

$$
\begin{aligned}
\left\|D f_{x}^{n}\left(\boldsymbol{v}^{s}\right)\right\|_{f^{n}(x)} & \leq C \lambda^{n}\left\|\boldsymbol{v}^{s}\right\|_{x} \text { for } \boldsymbol{v}^{s} \in \mathbb{E}_{x}^{s} \\
\left\|D f_{x}^{-n}\left(\boldsymbol{v}^{u}\right)\right\|_{f^{-n}(x)} & \leq C \lambda^{n}\left\|\boldsymbol{v}^{u}\right\|_{x} \text { for } \boldsymbol{v}^{u} \in \mathbb{E}_{x}^{u}
\end{aligned}
$$

hold.
Notice that the closure of the orbit of a transverse heteroclinic or homoclinic point is a simple example of a hyperbolic invariant set for $f$. For a Morse-Smale diffeomorphism, the set $\operatorname{Per}(f)$ of all periodic points is a finite hyperbolic invariant set. For an Anosov diffeomorphism, e.g. the toral Anosov automorphisms, the ambient manifold $M$ itself is a hyperbolic invariant set. We have many curious examples of hyperbolic invariant sets other than them, e.g. horseshoes, the Plykin attractor, the solenoid, or some invariant sets of Hénon-like maps. For example, see [De, Ro2].

As in the case of hyperbolic fixed points, we can define the stable and unstable manifolds for a hyperbolic invariant set as follows. Let $\Lambda$ be a hyperbolic invariant set for $f$. The stable and unstable manifolds $W^{s}(x)$ and $W^{u}(x)$ of $x \in \Lambda$ are defined as

$$
\begin{aligned}
W^{s}(x) & =\left\{y \in M ; \lim _{n \rightarrow \infty} d\left(f^{n}(x), f^{n}(y)\right)=0\right\} \\
W^{u}(x) & =\left\{y \in M ; \lim _{n \rightarrow \infty} d\left(f^{-n}(x), f^{-n}(y)\right)=0\right\}
\end{aligned}
$$

The unions

$$
W^{s}(\Lambda)=\bigcup_{x \in \Lambda} W^{s}(x), \quad W^{u}(\Lambda)=\bigcup_{x \in \Lambda} W^{u}(x)
$$

are called the stable and unstable manifolds for $\Lambda$, respectively. For $\varepsilon>0$, we identify the neighborhoods of each point $x \in \Lambda$ in $M$ with $U_{\varepsilon}(x)=\mathbb{E}_{x}^{s}(\varepsilon) \times \mathbb{E}_{x}^{u}(\varepsilon)$, where $\mathbb{E}_{x}^{s}(\varepsilon)=$ $\left\{\boldsymbol{v} \in \mathbb{E}_{x}^{s} ;\|\boldsymbol{v}\|_{x}<\varepsilon\right\}$ and $\mathbb{E}_{x}^{u}(\varepsilon)=\left\{\boldsymbol{v} \in \mathbb{E}_{x}^{u} ;\|\boldsymbol{v}\|_{x}<\varepsilon\right\}$. We define the local stable and local unstable manifolds $W_{\varepsilon}^{s}(x)$ and $W_{\varepsilon}^{u}(x)$ of $x \in \Lambda$ of size $\varepsilon$ as

$$
\begin{aligned}
& W_{\varepsilon}^{s}(x)=\left\{y \in U_{\varepsilon}(x) ; f^{j}(y) \in U_{\varepsilon}\left(f^{j}(x)\right) \text { for } j \geq 0\right\} \\
& W_{\varepsilon}^{u}(x)=\left\{y \in U_{\varepsilon}(x) ; f^{-j}(y) \in U_{\varepsilon}\left(f^{-j}(x)\right) \text { for } j \geq 0\right\}
\end{aligned}
$$

Now, we extend Stable Manifold Theorem to the case of hyperbolic invariant sets. See the Chapter 8 in [Ro1] for the proof.

Theorem 1.4 (Stable Manifold Theorem for hyperbolic invariant sets). Let $f$ be a $C^{r}(1 \leq$ $r \leq \infty)$ diffeomorphism on $M$ and let $\Lambda$ be a compact hyperbolic invariant set for $f$. Then there is an $\varepsilon>0$ such that, for each $x \in \Lambda$, there are two $C^{r}$ embedded disks $W_{\varepsilon}^{s}(x)$ and $W_{\varepsilon}^{u}(x)$ which are tangent to $\mathbb{E}_{x}^{s}$ and $\mathbb{E}_{x}^{u}$, respectively, and satisfy the following conditions.

- $W_{\varepsilon}^{s}(x)$ is represented by the graph of a $C^{r}$ function $\sigma_{x}^{s}: \mathbb{E}_{x}^{s}(\varepsilon) \rightarrow \mathbb{E}_{x}^{u}(\varepsilon)$ with $\sigma_{x}^{s}\left(\mathbf{0}_{x}\right)=$ $\mathbf{0}_{x}$ and $D \sigma_{x}^{s}(\mathbf{0})=\mathbf{0}$ :

$$
W_{\varepsilon}^{s}(x)=\left\{\left(\sigma_{x}^{s}(\boldsymbol{v}), \boldsymbol{v}\right) ; \boldsymbol{v} \in \mathbb{E}_{x}^{s}(\varepsilon)\right\}
$$

Besides, the function $\sigma_{x}^{s}$ and its derivatives vary continuously on x. Similarly, there is a $C^{r}$ function $\sigma_{x}^{u}: \mathbb{E}_{x}^{u}(\varepsilon) \rightarrow \mathbb{E}_{x}^{s}(\varepsilon)$ with $\sigma_{x}^{u}\left(\mathbf{0}_{x}\right)=\mathbf{0}_{x}$ and $D \sigma_{x}^{u}(\mathbf{0})=\mathbf{0}$ :

$$
W_{\varepsilon}^{u}(x)=\left\{\left(\boldsymbol{u}, \sigma_{x}^{u}(\boldsymbol{u})\right) ; \boldsymbol{u} \in \mathbb{E}_{x}^{u}(\varepsilon)\right\} .
$$

The function $\sigma_{x}^{u}$ and its derivatives also vary continuously on $x$.

- There exist $0<\lambda<1$ and $C \geq 1$ such that

$$
\begin{aligned}
& W_{\varepsilon}^{s}(x) \subset\left\{y \in U_{\varepsilon}(x) ; d\left(f^{j}(x), f^{j}(y)\right) \leq C \lambda^{j} d(x, y) \text { for } j \geq 0\right\}, \\
& W_{\varepsilon}^{u}(x) \subset\left\{y \in U_{\varepsilon}(x) ; d\left(f^{-j}(x), f^{-j}(y)\right) \leq C \lambda^{j} d(x, y) \text { for } j \geq 0\right\} .
\end{aligned}
$$

By Theorem 1.4, we have

$$
W^{s}(x)=\bigcup_{n \geq 0} f^{-n}\left(W_{\varepsilon}^{s}\left(f^{n}(x)\right)\right), \quad W^{u}(x)=\bigcup_{n \geq 0} f^{n}\left(W_{\varepsilon}^{s}\left(f^{-n}(x)\right)\right) .
$$

Notice that $W^{s}(x)$ and $W^{u}(x)$ are just the images of injective $C^{r}$ immersions from $\mathbb{R}^{s}$ and $\mathbb{R}^{u}$ to $M$ but not necessarily the images of embeddings, where $s=\operatorname{dim} \mathbb{E}_{x}^{s}$ and $u=\operatorname{dim} \mathbb{E}_{x}^{u}$. Horseshoes or toral Anosov automorphisms are typical examples of such diffeomorphisms. See the Chapter 8 in [Ro1].

### 1.3 Topological conjugacy and structural stability

For two diffeomorphisms $f$ and $g$, if the orbits for $f$ one-to-one correspond to those for $g$ with the same behavior, then we regard that $f$ and $g$ have essentially the same dynamical systems. For example, mutually conjugate linear maps satisfy the property. For classifying such diffeomorphisms, we introduce the notion of topological conjugacy.

Definition 1.5. We say that two diffeomorphisms $f$ and $g$ on a $C^{r}(1 \leq r \leq \infty)$ manifold $M$ are topologically conjugate to each other if there exists a homeomorphism $h: M \rightarrow M$ with $h \circ f=g \circ h$. This homeomorphism $h$ is called a topological conjugacy between $f$ and $g$.

Let $p$ be a periodic point for $f$ with period $n$ and set $p^{\prime}=h(p)$. Then $p^{\prime}$ satisfies $g^{n}\left(p^{\prime}\right)=g^{n}(h(p))=h\left(f^{n}(p)\right)=h(p)=p^{\prime}$. Thus, the point $p^{\prime}$ is also a periodic point for $g$ with the same period $n$.

A subset $D$ of $M$ is called a fundamental domain of $f$ if any non-periodic orbit of $f$ intersects $D$ exactly in one point. Fundamental domains are often used to construct topological conjugacies between diffeomorphisms. For example, let $f$ be a linear map on $\mathbb{R}^{2}$ with real contracting eigenvalues and $g$ another linear map on $\mathbb{R}^{2}$ with non-real contacting eigenvalues. Take a unit circle $C$ on $\mathbb{R}^{2}$, then one can have a pair of annuli $A_{f}$ and $A_{g}$ in $\mathbb{R}^{m}$ bounded by $C \cup f(C)$ and $C \cup g(C)$, respectively. See Figure 1.7. Then $A_{f}^{\prime}=A_{f} \backslash f(C)$ and $A_{g}^{\prime}=A_{g} \backslash g(C)$ are fundamental domains for $f$ and $g$, respectively. There exists a homeomorphism $\tilde{h}: A_{f} \rightarrow A_{g}$ with


Figure 1.7: Fundamental domains $A_{f}^{\prime}$ of $f$ and $A_{g}^{\prime}$ of $g$.

$$
\begin{equation*}
\tilde{h}(f(x))=g(\tilde{h}(x)) \tag{1.1}
\end{equation*}
$$

for any $x \in C$. Extend $\tilde{h}$ to the map $h: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ defined by

$$
h(x)=g^{-n(x)}\left(\tilde{h}\left(f^{n(x)}(x)\right)\right)
$$

for $x \in \mathbb{R}^{2} \backslash\{0\}$ and $h(0)=0$ for $0 \in \mathbb{R}^{2}$, where $n(x)$ is a uniquely determined integer with $f^{n(x)}(x) \in A_{f}^{\prime}$. By (1.1), $h$ is a well defined homeomorphism on $\mathbb{R}^{2}$, which is a topological conjugacy between $f$ and $g$.

For a given diffeomorphism $f$, we are interested in the topological conjugacy class constructed by diffeomorphisms close to $f$. Thus we introduce the notion of the structural stability for diffeomorphisms. A diffeomorphism $f \in \operatorname{Diff}^{r}(M)$ is called structurally stable if there exists a neighborhood $\mathcal{N} \subset \operatorname{Diff}^{r}(M)$ of $f$ such that, for any $g \in \mathcal{N}, f$ and $g$ are topologically conjugate.

Remark 1.6. In this definition of structural stability, the condition that $h$ is a homeomorphism is crucial. We suppose that $h$ is a diffeomorphism. Then $h$ is called a $C^{r}(1 \leq r \leq \infty)$ conjugacy between $f$ and $g$. If $f$ has a fixed point $p$, then, by the chain rule of composition maps, $D h(p) D f(p)=D g(h(p)) D h(p)$ holds. This shows that $D f(p)$ and $D g(h(p))$
are similar matrices via the matrix $D h(p)$. Thus they have the same eigenvalues. On the other hand, for any $f \in \operatorname{Diff}^{r}(M)$ and any fixed point $p$ of $f$, there exists $g \in \operatorname{Diff}^{r}(M)$ arbitrarily $C^{r}$ close to $f$ such that the eigenvalues of $D g\left(p^{\prime}\right)$ are different from those of $D f(p)$, where $p^{\prime}$ is the fixed point of $g$ corresponding to $p$. Namely, any neighborhood of $f \in \operatorname{Diff}^{r}(M)$ contains an element which is not $C^{r}$ conjugate to $f$. Thus, any diffeomorphism with a fixed point is not structurally stable with respect to $C^{r}$ conjugacy.

## 1.4 $C^{r}$ convergence of unstable manifolds

Let $p$ be a hyperbolic fixed point of a diffeomorphism $f$ on $M$ and $U(p)$ a sufficiently small neighborhood of $p$ in $M$. Take a disk $D$ embedded in $M$ of dimension ind ${ }^{u}(p)$ which intersects transversely the local stable manifold $W_{\text {loc }}^{s}(p)$ at a single point $z_{0}$. For any $n \in \mathbb{N}$, let $D_{n}$ be the component of $f^{n}(D) \cap U(p)$ containing $f^{n}\left(z_{0}\right)$. Then, $D_{n}$ uniformly $C^{r}$ converges to $W_{\text {loc }}^{u}(p)$ as $n \rightarrow \infty$. Figure 1.8 illustrates the cases of $\operatorname{ind}^{s}(p)=\operatorname{ind}^{u}(p)=1$ and $\operatorname{ind}^{s}(p)=1, \operatorname{ind}^{u}(p)=2$. More precisely, we have the following theorem called Inclination Lemma. See the Chapter 5 in [Ro1] for the proof.

(1)

(2)

Figure 1.8: (1) The case of $\operatorname{ind}^{s}(p)=\operatorname{ind}^{u}(p)=1$. (2) The case of $\operatorname{ind}^{s}(p)=1$ and $\operatorname{ind}^{u}(p)=2$.

Theorem 1.7 (Inclination Lemma). Let $f$ be $a C^{r}(1 \leq r \leq \infty)$ diffeomorphism on $M$ and $p \in M$ a saddle fixed point. Assume that $M$ has a coordinate neighborhood of $p$ such that $W_{\text {loc }}^{s}(p) \subset \mathbb{R}^{s} \times\{0\}$ and $W_{\text {loc }}^{u}(p) \subset\{0\} \times \mathbb{R}^{u}$, where $s=\operatorname{ind}^{s}(p)$ and $u=$ $\operatorname{ind}^{u}(p)$, if necessary by changing the coordinates suitably. Then, for any $C^{r}$ submanifold $D$ with $\operatorname{dim}(D)=u$ intersecting $W_{\text {loc }}^{s}(p)$ transversely at $z_{0}=\left(x_{0}, 0\right) \in W_{\text {loc }}^{s}(p) \times\{0\}$, the component $D_{n}$ of $f^{n}(D) \cap U(p)$ containing $f^{n}\left(z_{0}\right)$ uniformly $C^{r}$ converges to $W_{\mathrm{loc}}^{u}(p)$ as $n \rightarrow \infty$.

We consider the case that a diffeomorphism $f$ has a homoclinic tangency $r$ associated with a saddle fixed point $p$. First, suppose that $\operatorname{dim} M=2$ and $r$ is either a quadratic or cubic homoclinic tangency. It is not hard to show that $W^{u}(p)$ and $W^{s}(p)$ have a transverse intersection point $z$ in a neighborhood of $r$ under suitable open conditions of $f$. For example, see [GS1, GS2] if $r$ is a quadratic tangency and Lemma 1.2 in [KS1] if $r$ is a cubic tangency. Figure 1.9 illustrates the situations. Take an arc $D^{u}$ in $W^{u}(p)$ such that the interior of $D^{u}$ contains $z$. Then there exists an integer $N$ such that $f^{N}(z) \in U(p)$. Let $D_{0}^{u}$ be the connected component of $f^{N}\left(D^{u}\right) \cap U(p)$ containing $f^{N}(z)$. Let $D_{n}^{u}$ be the component of $f^{N+n}\left(D^{u}\right) \cap U(p)$ containing $f^{N+n}(z)$. By Inclination Lemma (Theorem 1.7), $D_{n}^{u} C^{r}$ converges to $W_{\text {loc }}^{u}(p)$ as $n \rightarrow \infty$.

Next, we consider the case that $\operatorname{dim} M=3$ and $\operatorname{ind}^{s}(p)=1$, $\operatorname{ind}^{u}(p)=2$. By [Ni], under certain open conditions of $f$, there exists a transverse intersection point $z$ of $W^{u}(p)$ and $W^{s}(p)$ near $r$. As in the case of $\operatorname{dim} M=2$, there exists a disk $\widetilde{D}^{u}$ in $W^{u}(p)$ such that the interior of $\widetilde{D}^{u}$ contains $z$. Again by Inclination Lemma, we can take the disk $\widetilde{D}_{n}^{u} C^{r}$ converging to $W_{\text {loc }}^{u}(p)$ as $n \rightarrow \infty$. The sequences $\left\{D_{n}^{u}\right\}$ and $\left\{\widetilde{D}_{n}^{u}\right\}$ are crucial in arguments of Chapters 2 and 3 , respectively.


Figure 1.9: (1) $r$ is a homoclinic quadratic tangency. (2) $r$ is a homoclinic cubic tangency.

### 1.5 Motivation and preceding results

Structurally stable diffeomorphisms have no heteroclinic or homoclinic tangencies. On the other hand, diffeomorphisms with heteroclinic or homoclinic tangencies are typical examples of structurally unstable diffeomorphisms. For such a diffeomorphism $f$, we need topological conjugacy invariants to dicide whether a given diffeomorphism $g$ is topological conjugate to $f$ or not. See Figure 1.10. Such topological conjugacy invariants are called modulus.

Definition 1.8. For a subspace $\mathcal{N}$ of the diffeomorphism space $\operatorname{Diff}^{r}(M)$ with $r \geq 1$, we


Figure 1.10: (1) The case that $f$ is a structurally stable diffeomorphism. (2) The case that $f$ is structurally unstable diffeomorphism.
say that a value $m(f)$ determined by $f \in \operatorname{Diff}^{r}(M)$ is a modulus in $\mathcal{N}$ if $m(g)=m(f)$ holds for any $g \in \mathcal{N}$ topologically conjugate to $f$.

The topological classification of structurally unstable diffeomorphisms on a manifold $M$ is an important subject in the study of dynamical systems. Palis [Pa] suggested that moduli play important roles in such a classification. The research of dynamical systems with moduli have been originated by Palis, de Melo and Takens. Subsequently, Posthumus, van Strien and others have studied enthusiastically this subject. See [dM, dMP, dMvS, GPvS, NPT, Pa, PT, Ta]. Our study in this thesis is based on results of Palis [Pa], de Melo [ dM ] and Posthumus [Po].

We will finish this section by introducing their results. First, we consider the case of $\operatorname{dim} M=2$. Suppose that $f_{i}(i=0,1)$ are elements of $\operatorname{Diff}^{2}(M)$ with two saddle fixed points $p_{i}, q_{i}$ such that $W^{u}\left(p_{i}\right)$ and $W^{s}\left(q_{i}\right)$ have a quadratic heteroclinic tangency $r_{i}$ and there exists a homeomorphism $h: M \rightarrow M$ with $f_{1}=h \circ f_{0} \circ h^{-1}, h\left(p_{0}\right)=p_{1}, h\left(q_{0}\right)=q_{1}$ and $h\left(r_{0}\right)=r_{1}$. See Figure 1.11 (1). Then, under some moderate conditions, Palis [Pa] proved that $\frac{\log \left|\lambda_{0}\right|}{\log \left|\mu_{0}\right|}=\frac{\log \left|\lambda_{1}\right|}{\log \left|\mu_{1}\right|}$, where $\lambda_{i}$ is the contracting eigenvalue of $D f\left(p_{i}\right)$ and $\mu_{i}$ is the expanding eigenvalue of $D f\left(q_{i}\right)$. This means that $m\left(f_{i}\right)=\frac{\log \left|\lambda_{i}\right|}{\log \left|\mu_{i}\right|}$ is one of moduli.

Following his result, de Melo [dM] studied the moduli of the stability of two-dimensional diffeomorphisms $f$, that is, a minimal set of moduli which parametrizes the topological conjugacy classes of $f$ in $\operatorname{Diff}^{r}(M)$. He detected moduli of stability for some classes of two-dimensional diffeomorphisms. In [dM], he also showed that the restrictions of the conjugacy homeomorphism $h$ on each $W^{s}\left(p_{0}\right) \backslash\left\{p_{0}\right\}$ and $W^{u}\left(q_{0}\right) \backslash\left\{q_{0}\right\}$ are local diffeomorphisms if $\frac{\log \left|\lambda_{0}\right|}{\log \left|\mu_{0}\right|}$ is irrational.

Subsequently, Posthumus [Po] proved that the homoclinic version of Palis and de Melo's results. In fact, he proved that, if $f_{i}(i=0,1)$ has a saddle fixed point $p_{i}$ with a homoclinic


Figure 1.11: (1) The situation in Palis' case. (2) The situation in Posthumus' case.
quadratic tangency $r_{i}$, then $\frac{\log \left|\lambda_{0}\right|}{\log \left|\mu_{0}\right|}=\frac{\log \left|\lambda_{1}\right|}{\log \left|\mu_{1}\right|}$ holds, where $\lambda_{i}$ and $\mu_{i}$ are the contracting and expanding eigenvalues of $D f\left(p_{i}\right)$, respectively. See Figure 1.11 (2). Moreover, if $\frac{\log \left|\lambda_{0}\right|}{\log \left|\mu_{0}\right|}$ is irrational, then the eigenvalues are also moduli, that is, $\lambda_{0}=\lambda_{1}$ and $\mu_{0}=\mu_{1}$.

For 2-dimensional diffeomorphisms, various results related to moduli concerning eigenvalues are obtained by some authors; see for example [dMP, dMvS, GPvS, PT]. However, in all of these results, the assumption that the tangency is quadratic or one-sided is crucial. In fact, some of their arguments do not work in the case that $q$ is a two-sided tangency, see Remark 2.9 for the reason.

## Chapter 2

## Moduli of surface diffeomorphisms with cubic tangencies

In this chapter, we study conjugacy invariants for 2-dimensional diffeomorphisms with cubic homoclinic tangencies (two-sided tangencies of the lowest order) under certain open conditions. Some of arguments used in previous works of conjugacy invariants associated with one-sided tangencies do not work in the two-sided case. We present a new method which is applicable to the two-sided case.

### 2.1 Moduli of surface diffeomorphisms with cubic tangencies

The following is the main result in this chapter.
Theorem 2.1. Suppose that $M$ is a closed surface with Riemannian metric. Let $f_{i}(i=$ $0,1)$ be elements of $\operatorname{Diff}^{3}(M)$ each of which has a saddle fixed point $p_{i}$ and a homoclinic cubic tangency $q_{i}$ associated with $p_{i}$ and satisfies the following conditions.
(A1) For $i=0,1$, there exists a neighborhood $U\left(p_{i}\right)$ of $p_{i}$ in $M$ such that $\left.f\right|_{U\left(p_{i}\right)}$ is linear.
(A2) $f_{0}$ is topologically conjugate to $f_{1}$ by a homeomorphism $h: M \rightarrow M$ with $h\left(p_{0}\right)=p_{1}$ and $h\left(q_{0}\right)=q_{1}$.
(A3) Each $f_{i}(i=0,1)$ satisfies the small expanding condition and one of the adaptable conditions with respect to $\left(p_{i}, q_{i}\right)$ in Section 2.8.

Then (M1) and (M2) hold, where $\lambda_{i}, \mu_{i}$ are the eigenvalues of $D f_{0}\left(p_{i}\right)$ with $0<\left|\lambda_{i}\right|<$ $1<\left|\mu_{i}\right|$.
(M1) $\frac{\log \left|\lambda_{0}\right|}{\log \left|\mu_{0}\right|}=\frac{\log \left|\lambda_{1}\right|}{\log \left|\mu_{1}\right|}$.
(M2) Moreover, if $\frac{\log \left|\lambda_{0}\right|}{\log \left|\mu_{0}\right|}$ is irrational, then $\mu_{0}=\mu_{1}$ and $\lambda_{0}=\lambda_{1}$.

Here we say that $f_{0}$ satisfies the small expanding condition at $p_{0}$ if $\left|\mu_{0}\right|=1+\varepsilon$ with $0<\varepsilon<\varepsilon_{0}$ for the constant $\varepsilon_{0}$ given in Lemma 2.6. Note that this condition depends on local expressions of $f_{0}$ such as (2.2) near $p_{0}$ and (2.5) near $f_{0}^{m_{0}}\left(q_{0}\right)$. In Section 2.2, we present a codimension two submanifold $\mathcal{C}$ of $\operatorname{Diff}^{3}(M)$ such that any element of $\mathcal{C}$ sufficiently close to $f_{0}$ also satisfies (A3). In the case that $f$ is of class $C^{\infty}$, we know from Sternberg [St] and Takens [Ta] that (A1) is an open dense condition in Diff ${ }^{\infty}(M)$.

Though we only consider the case of cubic tangencies, we believe that our method still works in the case of two-sided tangencies of higher order. So we propose the following question.

Question 2.2. Is it possible to generalize our theorem to the case where diffeomorphisms have two-sided homoclinic tangencies of higher order?

We will finish the introduction by outlining the proof of the main theorem. Let $f_{0}$ be a diffeomorphism satisfying the conditions of Theorem 2.1. We may assume that $q_{0}$ and $r_{0}=\varphi\left(q_{0}\right)$ are contained in $W_{\text {loc }}^{u}\left(p_{0}\right)$ and $W_{\text {loc }}^{s}\left(p_{0}\right)$ respectively, where $\varphi=f_{0}^{m_{0}}$ for some positive integer $m_{0}$. For the proof of Theorem 2.1, we need to find out a useful connection between the eigenvalues $\mu_{i}$ and $\lambda_{i}$ for $i=0,1$. By applying Inclination Lemma (Lemma 1.7), we have a sequence $\left\{\alpha_{n}^{u}\right\}$ of $\operatorname{arcs}$ in $W^{u}\left(p_{0}\right)$ which meet $W_{\text {loc }}^{s}\left(p_{0}\right)$ transversely at single points $z_{0} \lambda_{0}^{n}$ and $C^{3}$ converge to a sub-arc of $W_{\text {loc }}^{u}\left(p_{0}\right)$. See Figure 2.5. Then $\varphi\left(\alpha_{n}^{u}\right)$ contains an S-shaped arc $\gamma_{0, n}^{\prime}$ framed by the rectangle $S_{n}$ as illustrated in Figure 2.1. We note that


Figure 2.1: $\gamma_{0, n}^{\prime}$ is an S-shaped arc. $\widehat{\gamma}_{1}$ and $\widehat{\gamma}_{k}$ are compressed S-shaped arcs near $q_{0}$ induced from $\gamma_{0, n}^{\prime}$.
such arcs $\gamma_{0, n}^{\prime}$ are subtle and vanish eventually as $n \rightarrow \infty$. See Figures 2.6 and 2.13. Since $h$ is not supposed to be smooth, one can not expect that $h$ sends $\gamma_{0, n}^{\prime}$ to an S-shaped curve in $W^{u}\left(p_{1}\right)$. However Intersection Lemma (Lemma 2.7) shows that it actually holds, which is a key lemma in our argument. For the proof, we send $\gamma_{0, n}^{\prime}$ to a curve $\widehat{\gamma}_{1}$ in a small
neighborhood of $q_{0}$ by $f_{0}^{u_{0}}$ for some $u_{0} \in \mathbb{N}$ and pull it back near $r_{0}$ by $\varphi$. Repeating this process many times, one can amplify $\widehat{\gamma}_{1}$ and finally have a compressed S-shaped curve $\widehat{\gamma}_{k}$ near $q$ the diameter of which is substantial so that it can be distinguished by $h$. From this fact, we know that $h\left(\widehat{\gamma}_{k}\right)$ intersects a compressed $S$-shaped curve $\widehat{\gamma}_{k}^{*}$ in $W^{u}\left(p_{1}\right)$. It follows that there exists a sequence $\left\{r_{n}\right\}$ with $r_{n} \in \gamma_{0, n}^{\prime}$ as illustrated in Figure 2.1 such that $\bar{r}_{n}=h\left(r_{n}\right)$ is contained in the corresponding S-shaped curve $\bar{\gamma}_{0, n}^{\prime}$ in $W^{u}\left(p_{1}\right)$. We note that the images of $r_{n}, \bar{r}_{n}$ by the orthogonal projections to the first coordinates are represented as $a z_{0} \lambda_{0}^{n}+o\left(\lambda_{0}^{n}\right), \bar{a} \bar{z}_{0} \lambda_{1}^{n}+o\left(\lambda_{1}^{n}\right)$ respectively for some non-zero constants $a$, $\bar{a}$. One can take subsequences $\{n(k)\},\{m(k)\}$ of $\mathbb{N}$ such that $f_{0}^{m(k)}\left(r_{n(k)}\right)$ converges to a point $x_{0} \in W_{\text {loc }}^{u}\left(p_{0}\right)$. Then $f_{1}^{m(k)}\left(\bar{r}_{n(k)}\right)$ also converges to $h\left(x_{0}\right) \in W_{\text {loc }}^{u}\left(p_{1}\right)$. By using this fact, we will show that $\lim _{k \rightarrow \infty} \frac{m(k)}{n(k)}=-\frac{\log \lambda_{0}}{\log \mu_{0}}$ and $\lim _{k \rightarrow \infty} \frac{m(k)}{n(k)}=-\frac{\log \lambda_{1}}{\log \mu_{1}}$. This proves the assertion (M1). The assertion (M2) is proved by (M1) together with standard arguments in [dM, Po].

### 2.2 Preliminaries

Let $\left\{a_{n}\right\},\left\{b_{n}\right\}$ be sequences with non-zero entries. Then $a_{n} \approx b_{n}$ means that $\frac{a_{n}}{b_{n}} \rightarrow 1$ as $n \rightarrow \infty$, and $a_{n} \sim b_{n}$ means that there exist constants $C$ and $C^{\prime}$ independent of $n$ with $0<C^{\prime}<1<C$ and satisfying $C^{\prime} \leq \frac{a_{n}}{b_{n}} \leq C$ for any $n$. Suppose next that $\left\{a_{n}\right\},\left\{b_{n}\right\}$ are sequences with non-negative entries. If there exists a constant $C^{\prime}>0$ independent of $n$ and satisfying $a_{n} \leq C^{\prime} b_{n}$ for any $n$, then we denote the property by $a_{n} \precsim b_{n}$.

Throughout the remainder of this chapter, we suppose that $M$ is a closed connected surface and $f: M \rightarrow M$ is a $C^{3}$ diffeomorphism with a saddle fixed point $p$. Let $\mu, \lambda$ be the eigenvalues of $D f(p)$ with

$$
\begin{equation*}
0<|\lambda|<1<|\mu| \tag{2.1}
\end{equation*}
$$

Suppose moreover that $f$ is $C^{3}$ linearizable in a neighborhood $U(p)$ of $p$ in $M$. Then there exists a $C^{3}$ coordinate $(x, y)$ on $U(p)$ satisfying the following condition:

$$
\begin{equation*}
f(x, y)=(\mu x, \lambda y) \tag{2.2}
\end{equation*}
$$

for any $(x, y) \in U(p)$. In particular, this implies that $p=(0,0), W_{\text {loc }}^{u}(p):=\{(x, y) \in$ $U(p) ; y=0\} \subset W^{u}(p)$ and $W_{\text {loc }}^{s}(p):=\{(x, y) \in U(p) ; x=0\} \subset W^{s}(p)$.

Let $\mathcal{C}$ be the subspace of Diff ${ }^{3}(M)$ consisting of elements $f \in \operatorname{Diff}^{3}(M)$ satisfying the following conditions (C1)-(C3).
(C1) $f$ has a saddle periodic point $p$.
(C2) There exists a homoclinic cubic tangency $q$ associated with $p$.
(C3) $f$ satisfies the adaptable conditions in the sense of Section 2.8 with respect to $p, q$.

Note that $\mathcal{C}$ is a codimension two submanifold of $\operatorname{Diff}^{3}(M)$.
Let $q$ be a cubic tangency of $W^{u}(p)$ and $W^{s}(p)$. We assume that $q$ is contained in $W_{\text {loc }}^{u}(p) \subset U(p)$ if necessary replacing $q$ by $f^{-n}(q)$ with sufficiently large $n \in \mathbb{N}$. For the point $q$, there exists $m_{0} \in \mathbb{N}$ such that $r:=f^{m_{0}}(q) \in W_{\text {loc }}^{s}(p) \subset U(p)$. Then one can rearrange the linearizing coordinate on $U(p)$ so that $q=(1,0), r=(0,1)$. Moreover, we may suppose that

$$
U(p)=[-2,2] \times[-2,2], W_{\mathrm{loc}}^{u}(p)=[-2,2] \times\{0\}, W_{\mathrm{loc}}^{s}(p)=\{0\} \times[-2,2] .
$$

Let $U(q), U(r)$ be sufficiently small neighborhoods of $q, r$ in $U(p)$ respectively. Then the component $L^{s}(q)$ of $W^{s}(p) \cap U(q)$ containing $q$ is represented as

$$
L^{s}(q)=\{(x+1, y) \in U(q) ; y=v(x)\}
$$

where $v$ is a $C^{3}$ function satisfying

$$
\begin{equation*}
v(0)=v^{\prime}(0)=v^{\prime \prime}(0)=0 \quad \text { and } \quad v^{\prime \prime \prime}(0) \neq 0 . \tag{2.3}
\end{equation*}
$$

Similarly, the component $L^{u}(r)$ of $W^{u}(p) \cap U(r)$ containing $r$ is represented as

$$
L^{u}(r)=\{(x, y+1) \in U(r) ; x=w(y)\},
$$

where $w$ is a $C^{3}$ function satisfying

$$
\begin{equation*}
w(0)=w^{\prime}(0)=w^{\prime \prime}(0)=0 \quad \text { and } \quad w^{\prime \prime \prime}(0) \neq 0, \tag{2.4}
\end{equation*}
$$

see Figure 2.2


Figure 2.2: $q$ and $r$ are homoclinic cubic tangencies associated with $p$.
Recall that $q=(1,0), r=(0,1)$ are cubic tangencies between $W^{s}(p)$ and $W^{u}(p)$ and $f^{m_{0}}(q)=r$ for some $m_{0} \in \mathbb{N}$. We set $f^{m_{0}}=\varphi$ for short. By (2.3) and (2.4), $\varphi$ is represented in $U(q)$ as follows for some constants $a, b, c, d, e$.

$$
\begin{equation*}
\varphi(x+1, y)=\left(a y+b x y+c x^{3}+H_{1}(x+1, y), 1+d x+e y+H_{2}(x+1, y)\right), \tag{2.5}
\end{equation*}
$$

where $H_{1}, H_{2}$ are $C^{3}$ functions satisfying the following conditions.

$$
\begin{align*}
H_{1}(1,0) & =\partial_{x} H_{1}(1,0)=\partial_{y} H_{1}(1,0)=\partial_{x x} H_{1}(1,0)=\partial_{x y} H_{1}(1,0) \\
& =\partial_{x x x} H_{1}(1,0)=0  \tag{2.6}\\
H_{2}(1,0) & =\partial_{x} H_{2}(1,0)=\partial_{y} H_{2}(1,0)=0
\end{align*}
$$

Since $\varphi$ is a diffeomorphism,

$$
a, d \neq 0
$$

The fact that $q$ is a cubic tangency implies

$$
c \neq 0
$$

Here we put the following extra open condition.

$$
\begin{equation*}
b \neq 0 \tag{2.7}
\end{equation*}
$$

By (2.5) and (2.6), the Jacobian matrix of $\varphi$ at $(x+1, y)$ is given as follows.

$$
\begin{align*}
D \varphi(x+1, y) & =\left[\begin{array}{cc}
b y+3 c x^{2}+\partial_{x} H_{1}(x+1, y) & a+b x+\partial_{y} H_{1}(x+1, y) \\
d+\partial_{x} H_{2}(x+1, y) & e+\partial_{y} H_{2}(x+1, y)
\end{array}\right] \\
& =\left[\begin{array}{cc}
b y+3 c x^{2}+o\left(x^{2}\right)+o(y)+O(x y) & a+b x+o(x)+O(y) \\
d+O(x)+O(y) & e+O(x)+O(y)
\end{array}\right] \tag{2.8}
\end{align*}
$$

Here we only consider the case satisfying the following condition, which belongs to Case $\mathrm{II}_{++}$in Section 2.8.

$$
\begin{equation*}
0<\lambda<1, \mu>1, a>0, b<0, c>0, d<0 \tag{2.9}
\end{equation*}
$$

See Figure 2.3 for the situation of $W_{\text {loc }}^{u}(p)$ and $W_{\text {loc }}^{s}(p)$ in the case of (2.9). Note that (2.9) implies the extra condition (2.7).

One can set $\mu=1+\varepsilon$ for some $\varepsilon>0$. We only consider the case that $\varepsilon$ is sufficiently small.

Consider the rectangle $R_{\varepsilon}=\left[1+\varepsilon,(1+\varepsilon)^{3}\right] \times\left[0, \varepsilon^{3}\right]$ in $U(q)$. By (2.5),

$$
\begin{align*}
\varphi(1+\varepsilon, 0) & =\left(c \varepsilon^{3}+o\left(\varepsilon^{3}\right), 1+d \varepsilon+o(\varepsilon)\right) \\
\varphi\left(1+\varepsilon, \varepsilon^{3}\right) & =\left((a+c) \varepsilon^{3}+o\left(\varepsilon^{3}\right), 1+d \varepsilon+o(\varepsilon)\right)  \tag{2.10}\\
\varphi\left((1+\varepsilon)^{3}, 0\right) & =\left(27 c \varepsilon^{3}+o\left(\varepsilon^{3}\right), 1+3 d \varepsilon+o(\varepsilon)\right) \\
\varphi\left((1+\varepsilon)^{3}, \varepsilon^{3}\right) & =\left((a+27 c) \varepsilon^{3}+o\left(\varepsilon^{3}\right), 1+3 d \varepsilon+o(\varepsilon)\right)
\end{align*}
$$

Let $\operatorname{pr}_{x}: U(p) \rightarrow W_{\text {loc }}^{u}(p)$ and $\operatorname{pr}_{y}: U(p) \rightarrow W_{\text {loc }}^{s}(p)$ be the orthogonal projections with respect to the linearizing coordinate on $U(p)$. Then there exist constants $\tau_{0}, \tau_{1}$ with $0<\tau_{0}<\tau_{1}$ independent of $\varepsilon$ and satisfying

$$
\begin{equation*}
\operatorname{pr}_{x}\left(\varphi\left(R_{\varepsilon}\right)\right) \subset\left[\tau_{0} \varepsilon^{3}, \tau_{1} \varepsilon^{3}\right] \tag{2.11}
\end{equation*}
$$



Figure 2.3: The case of $\mathrm{II}_{++}$.

Since $d<0$ by (2.9), it follows from (2.10) that

$$
\begin{equation*}
\operatorname{pr}_{y}\left(\varphi\left(R_{\varepsilon}\right)\right) \subset[1+3.5 d \varepsilon, 1+0.5 d \varepsilon] \subset[1+4 d \varepsilon, 1] . \tag{2.12}
\end{equation*}
$$

For any $\boldsymbol{x} \in R_{\varepsilon}$, let $u_{0}=u_{0}(\boldsymbol{x})$ be a uniquely determined positive integer such that $f^{i}(\varphi(\boldsymbol{x})) \in U(p)$ for $i=1, \ldots, u_{0}$ and $\operatorname{pr}_{x}\left(f^{u_{0}}(\varphi(\boldsymbol{x}))\right) \subset\left((1+\varepsilon)^{2},(1+\varepsilon)^{3}\right]$. Since $\operatorname{pr}_{x}\left(f^{u_{0}}(\varphi(\boldsymbol{x}))\right)=\mu^{u_{0}} \operatorname{pr}_{x}(\varphi(\boldsymbol{x}))$,

$$
1<(1+\varepsilon)^{2}<\mu^{u_{0}} \operatorname{pr}_{x}(\varphi(\boldsymbol{x}))<\tau_{1} \mu^{u_{0}} \varepsilon^{3} .
$$

Since $\operatorname{pr}_{y}(\varphi(\boldsymbol{x}))<1$ by (2.12), it follows that

$$
\operatorname{pr}_{y}\left(f^{u_{0}}(\varphi(\boldsymbol{x}))\right)=\lambda^{u_{0}} \operatorname{pr}_{y}(\varphi(\boldsymbol{x}))<\lambda^{u_{0}} .
$$

Consider the following conditions for $\varepsilon>0$ :

$$
\begin{equation*}
\tau_{1}<\varepsilon^{-1} \quad \text { and } \quad(1+\varepsilon)^{\frac{3}{2}}=\mu^{\frac{3}{2}}<\lambda^{-1} . \tag{2.13}
\end{equation*}
$$

If these conditions are satisfied, then the following inequalities

$$
\begin{equation*}
1<1+\varepsilon<\mu^{u_{0}} \operatorname{pr}_{x}(\varphi(\boldsymbol{x}))<\tau_{1} \mu^{u_{0}} \varepsilon^{3}<\mu^{u_{0}} \varepsilon^{2} \tag{2.14}
\end{equation*}
$$

hold. This implies that

$$
\operatorname{pr}_{y}\left(f^{u_{0}}(\varphi(\boldsymbol{x}))\right)=\lambda^{u_{0}} \operatorname{pr}_{y}(\varphi(\boldsymbol{x}))<\lambda^{u_{0}}<\mu^{-\frac{3}{2} u_{0}}<\varepsilon^{3} .
$$

Thus the positive integer $u_{0}(\boldsymbol{x})$ satisfies

$$
\begin{equation*}
f^{u_{0}(\boldsymbol{x})}(\varphi(\boldsymbol{x})) \in R_{\varepsilon} \tag{2.15}
\end{equation*}
$$

for all $\boldsymbol{x} \in R_{\varepsilon}$. See Figure 2.4.


Figure 2.4: The rectangles $R_{\varepsilon}$ and $f^{u_{0}(\boldsymbol{x})}\left(\varphi\left(R_{\varepsilon}\right)\right)$ for $\boldsymbol{x} \in R_{\varepsilon}$.

### 2.3 Sequence of Rectangles

Let $f: M \rightarrow M$ be a $C^{3}$ diffeomorphism given in Section 2.2. In particular, $f$ satisfies the linearizing condition (2.2) on $U(p)$. As is seen in Subsection 1.4, $W^{u}(p)$ and $W^{s}(p)$ have a transverse intersection point other than $p$. Let $\delta^{u}$ be a segment in $W_{\text {loc }}^{u}(p)$ with $\operatorname{Int} \delta^{u} \supset\{p, q\}$. Then, by Inclination Lemma (Theorem 1.7), there exists a sequence $\left\{\alpha_{n}^{u}\right\}_{n=0}^{\infty}$ of arcs in $W^{u}(p) C^{3}$ converging to $\delta^{u}$ and satisfying the following conditions:

- $\alpha_{0}^{u}$ meets $W_{\text {loc }}^{s}(p)$ transversely in a single point $\boldsymbol{z}_{0}=\left(0, z_{0}\right)$.
- Each $\alpha_{n}^{u}$ contains $f^{n}\left(\boldsymbol{z}_{0}\right)=\left(0, z_{0} \lambda^{n}\right)$, and the intersection $\tilde{\alpha}_{n}^{u}=\alpha_{n}^{u} \cap U(q)$ is an arc meeting $L^{s}(q)$ transversely in a single point $c_{n}$ for any sufficiently large $n>0$.

See Figure 2.5. Note that $\alpha_{0}^{u}$ is represented by the graph of a $C^{3}$-function $y_{0}: \delta^{u} \rightarrow \mathbb{R}_{+}$,


Figure 2.5: A sequence $\left\{\alpha_{n}^{u}\right\}_{n=0}^{\infty} C^{3}$-converging to $\delta^{u}$.
that is, $\alpha_{0}^{u}=\left\{\left(x, y_{0}(x)\right) ; x \in \delta^{u}\right\}$. Then each $\alpha_{n}^{u}$ is represented by the graph of the function $y_{n}: \delta^{u} \rightarrow \mathbb{R}_{+}$with

$$
\begin{equation*}
y_{n}(x)=\lambda^{n} y_{0}\left(\mu^{-n} x\right) \quad \text { for } \quad x \in \delta^{u} . \tag{2.16}
\end{equation*}
$$

We parametrise $\tilde{\alpha}_{n}^{u}$ in $\left[(1+\varepsilon)^{-3},(1+\varepsilon)^{3}\right]$ by $\alpha_{n}(t)=\left(t+1, \tilde{y}_{n}(t)\right)$ with $(1+\varepsilon)^{-3}-1 \leq$ $t \leq(1+\varepsilon)^{3}-1$, where $\tilde{y}_{n}(t)=y_{n}(t+1)$. By (2.5) and (2.8),

$$
\begin{array}{r}
\varphi\left(\alpha_{n}(t)\right)=\left(a \tilde{y}_{n}(t)+b t \tilde{y}_{n}(t)+c t^{3}+\text { h.o.t., } 1+d t+e \tilde{y}_{n}(t)+\text { h.o.t. }\right) \\
D \varphi\left(\alpha_{n}(t)\right)\left(\alpha_{n}^{\prime}(t)\right)=\left(a \tilde{y}_{n}^{\prime}(t)+b \tilde{y}_{n}(t)+b t \tilde{y}_{n}^{\prime}(t)+3 c t^{2}+\right.\text { h.o.t., }  \tag{2.18}\\
\left.d+e \tilde{y}_{n}^{\prime}(t)+\text { h.o.t. }\right)
\end{array}
$$

where the primes represent the derivative on $t$ and 'h.o.t.' denotes the sum of the higher order terms on $t$. By (2.16),

$$
\left|\tilde{y}_{n}^{\prime}(t)\right|=\left|y_{n}^{\prime}(t+1)\right|=\lambda^{n} \mu^{-n}\left|y_{0}^{\prime}\left(\mu^{-n}(t+1)\right)\right| .
$$

Suppose that $\sigma$ is the maximum of $\left|y_{0}^{\prime}(x)\right|$ on $\delta^{u}$. Then

$$
\left|\tilde{y}_{n}^{\prime}(t)\right|=\left|y_{n}^{\prime}(t+1)\right|=\lambda^{n} \mu^{-n}\left|y_{0}^{\prime}\left(\mu^{-n}(t+1)\right)\right| \leq \lambda^{n} \mu^{-n} \sigma
$$

for any $n \in \mathbb{N}$. This implies that

$$
\begin{equation*}
\left|\tilde{y}_{n}^{\prime}(t)\right| \precsim \lambda^{n} \mu^{-n} \tag{2.19}
\end{equation*}
$$

Suppose that $d \varphi_{\alpha_{n}(t)}\left(\alpha_{n}^{\prime}(t)\right)$ is vertical at $t=t_{n}$. Then $\lim _{n \rightarrow \infty} t_{n}=0$ and, by (2.18),

$$
b \tilde{y}_{n}\left(t_{n}\right)+\left(a+b t_{n}\right) \tilde{y}_{n}^{\prime}\left(t_{n}\right) \approx-3 c t_{n}^{2}
$$

Since $\tilde{y}_{n}(t) \approx \lambda^{n} z_{0}$ and $\left|\tilde{y}_{n}^{\prime}(t)\right| \precsim \lambda^{n} \mu^{-n}$, this condition is equivalent to

$$
\begin{equation*}
3 c t_{n}^{2} \approx-b \tilde{y}_{n}\left(t_{n}\right) \approx-b \lambda^{n} z_{0} \tag{2.20}
\end{equation*}
$$

It follows that, for all sufficiently large $n, d \varphi_{\alpha_{n}(t)}\left(\alpha_{n}^{\prime}(t)\right)$ is vertical at two points $t_{n, \pm}$ with

$$
\begin{equation*}
t_{n, \pm} \approx \pm \sqrt{\frac{-b z_{0}}{3 c}} \lambda^{\frac{n}{2}} \tag{2.21}
\end{equation*}
$$

Let $\tilde{t}_{n, \pm}$ be the elements of $\left[(1+\varepsilon)^{-3}-1,(1+\varepsilon)^{3}-1\right]$ with $\tilde{t}_{n,-}<t_{n,-}, t_{n,+}<\tilde{t}_{n,+}$ such that $\varphi\left(\alpha_{n}\left(\tilde{t}_{n, \pm}\right)\right)$ is the intersection point of $\varphi\left(\alpha_{n}(t)\right)$ and the vertical line $L_{n, \pm}$ tangent to $\varphi\left(\alpha_{n}(t)\right)$ at $\varphi\left(\alpha_{n}\left(t_{n, \mp}\right)\right)$. Let $S_{n}$ be the smallest orthogonal rectangle in $U(r)$ containing the four points $\varphi\left(\alpha_{n}\left(\tilde{t}_{n,-}\right)\right), \varphi\left(\alpha_{n}\left(t_{n,-}\right)\right), \varphi\left(\alpha_{n}\left(t_{n,+}\right)\right), \varphi\left(\alpha_{n}\left(\tilde{t}_{n,+}\right)\right)$. See Figure 2.6.

Now we will estimate the size of $S_{n}$. Let $D_{n}$ be the distance between $S_{n}$ and $W_{\text {loc }}^{s}(p)$. Then

$$
\begin{align*}
D_{n} & \approx a \tilde{y}_{n}\left(t_{n,+}\right)+b t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right)+c t_{n,+}^{3} \\
& \approx a z_{0} \lambda^{n}+b z_{0} \sqrt{\frac{-b z_{0}}{3 c}} \lambda^{\frac{3}{2} n}-\frac{b z_{0}}{3} \sqrt{\frac{-b z_{0}}{3 c}} \lambda^{\frac{3}{2} n} \sim \lambda^{n} \tag{2.22}
\end{align*}
$$

By (2.5), the width $W_{0, n}$ of $S_{n}$ is represented as

$$
\begin{aligned}
W_{0, n} & \approx\left(a \tilde{y}_{n}\left(t_{n,-}\right)+b t_{n,-} \tilde{y}_{n}\left(t_{n,-}\right)+c t_{n,-}^{3}\right)-\left(a \tilde{y}_{n}\left(t_{n,+}\right)+b t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right)+c t_{n,+}^{3}\right) \\
& =a\left(\tilde{y}_{n}\left(t_{n,-}\right)-\tilde{y}_{n}\left(t_{n,+}\right)\right)+b\left(t_{n,-} \tilde{y}_{n}\left(t_{n,-}\right)-t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right)\right)+c\left(t_{n,-}^{3}-t_{n,+}^{3}\right)
\end{aligned}
$$



Figure 2.6: The smallest orthogonal rectangle $S_{n}$.

It follows from Mean Value Theorem together with (2.19) that

$$
\left|\tilde{y}_{n}\left(t_{n,-}\right)-\tilde{y}_{n}\left(t_{n,+}\right)\right| \precsim \lambda^{n} \mu^{-n}\left|t_{n,-}-t_{n,+}\right| \sim \lambda^{\frac{3}{2} n} \mu^{-n} .
$$

Moreover, by (2.21), we have

$$
\begin{aligned}
c\left(t_{n,-}^{3}-t_{n,+}^{3}\right) & \approx c\left(t_{n,-}\left(\frac{-b \tilde{y}_{n}\left(t_{n,-}\right)}{3 c}\right)-t_{n,+}\left(\frac{-b \tilde{y}_{n}\left(t_{n,+}\right)}{3 c}\right)\right) \\
& =-\frac{b}{3}\left(t_{n,-} \tilde{y}_{n}\left(t_{n,-}\right)-t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right)\right) .
\end{aligned}
$$

Since

$$
\begin{aligned}
t_{n,-} \tilde{y}_{n}\left(t_{n,-}\right)-t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right) & =\left(t_{n,-}-t_{n,+}\right) \tilde{y}_{n}\left(t_{n,-}\right)+t_{n,+}\left(\tilde{y}_{n}\left(t_{n,-}\right)-y_{n}\left(t_{n,+}\right)\right) \\
& \approx-\sqrt{\frac{-b z_{0}}{3 c}} \lambda^{\frac{n}{2}} \cdot z_{0} \lambda^{n}+O\left(\lambda^{\frac{n}{2}} \cdot \lambda^{\frac{3}{2} n} \mu^{-n}\right) \sim-\lambda^{\frac{3}{2} n},
\end{aligned}
$$

we have

$$
\begin{equation*}
W_{0, n} \approx O\left(\lambda^{\frac{3}{2} n} \mu^{-n}\right)+\frac{2 b}{3}\left(t_{n,-} \tilde{y}_{n}\left(t_{n,-}\right)-t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right)\right) \sim \lambda^{\frac{3}{2} n} . \tag{2.23}
\end{equation*}
$$

Next we estimate the height $H_{0, n}$ of $S_{n}$. For that, we estimate $W_{0, n}$ again by using $\tilde{t}_{n,+}$ and $t_{n,+}$ instead of $t_{n,-}$ and $t_{n,+}$. Since $\tilde{t}_{n,+}>t_{n,+}$, one can set $\tilde{t}_{n,+}=t_{n,+}+\rho_{n} \lambda^{\frac{n}{2}}$ for some $\rho_{n}>0$.

$$
\begin{aligned}
& W_{0, n} \approx a\left(\tilde{y}_{n}\left(\tilde{t}_{n,+}\right)-\tilde{y}_{n}\left(t_{n,+}\right)\right)+b\left(\tilde{t}_{n,+} \tilde{y}_{n}\left(\tilde{t}_{n,+}\right)-t_{n,+} \tilde{y}_{n}\left(t_{n,+}\right)\right)+c\left(\tilde{t}_{n,+}^{3}-\tilde{t}_{n,-}^{3}\right) \\
& \quad=\left(a+b \tilde{t}_{n,+}\right)\left(\tilde{y}_{n}\left(\tilde{t}_{n,+}\right)-\tilde{y}_{n}\left(t_{n,+}\right)\right)+b\left(\tilde{t}_{n,+}-t_{n,+}\right) \tilde{y}_{n}\left(t_{n,+}\right)+c\left(\tilde{t}_{n,+}^{3}-t_{n,-}^{3}\right) .
\end{aligned}
$$

Again by Mean Value Theorem together with (2.19),

$$
\left|\tilde{y}_{n}\left(\tilde{t}_{n,+}\right)-\tilde{y}_{n}\left(t_{n,+}\right)\right| \precsim \lambda^{n} \mu^{-n} \cdot \rho_{n} \lambda^{\frac{n}{2}}=\rho_{n} \lambda^{\frac{3}{2} n} \mu^{-n} .
$$

Moreover, we have

$$
\left(\tilde{t}_{n,+}-t_{n,+}\right) \tilde{y}_{n}\left(t_{n,+}\right) \sim \rho_{n} \lambda^{\frac{n}{2}} \cdot \lambda^{n}=\rho_{n} \lambda^{\frac{3}{2} n}
$$

and

$$
\begin{aligned}
\tilde{t}_{n,+}^{3}-t_{n,+}^{3} & =3 \rho_{n}^{2} \lambda^{n} t_{n,+}+3 \rho_{n} \lambda^{\frac{n}{2}} t_{n,+}^{2}+\rho_{n}^{3} \lambda^{\frac{3}{2} n} \\
& \approx\left(3 \rho_{n} \sqrt{\frac{-b z_{0}}{3 c}}-\frac{3 b z_{0}}{c}+\rho_{n}^{2}\right) \rho_{n} \lambda^{\frac{3}{2} n}
\end{aligned}
$$

This shows that

$$
W_{0, n} \sim\left(a \mu^{-n}+b z_{0}+3 \rho_{n} \sqrt{\frac{-b z_{0}}{3 c}}-\frac{3 b z_{0}}{c}+\rho_{n}^{2}\right) \rho_{n} \lambda^{\frac{3}{2} n}
$$

Since $W_{0, n} \sim \lambda^{\frac{3}{2} n}$, it follows that $\rho_{n} \sim 1$ and hence $\tilde{t}_{n,+} \sim \lambda^{\frac{n}{2}}$. Similarly $-\tilde{t}_{n,-} \sim \lambda^{\frac{n}{2}}$. This implies that

$$
\begin{equation*}
\left|\tilde{t}_{n, \pm}\right| \sim \lambda^{\frac{n}{2}} \tag{2.24}
\end{equation*}
$$

Therefore we have

$$
\begin{align*}
H_{0, n} & =\left(1+d \tilde{t}_{n,-}+e \tilde{y}_{n}\left(\tilde{t}_{n,-}\right)\right)-\left(1+d \tilde{t}_{n,+}+e \tilde{y}_{n}\left(\tilde{t}_{n,+}\right)\right) \\
& =d\left(\tilde{t}_{n,-}-\tilde{t}_{n,+}\right)+e\left(\tilde{y}_{n}\left(\tilde{t}_{n,-}\right)-\tilde{y}_{n}\left(\tilde{t}_{n,+}\right)\right) \sim \lambda^{\frac{n}{2}}+O\left(\lambda^{\frac{3}{2} n} \mu^{-n}\right) \sim \lambda^{\frac{n}{2}} \tag{2.25}
\end{align*}
$$

In particular, $\left\{S_{n}\right\}$ is a sequence of rectangles converging to the cubic tangency $r$.

### 2.4 Slope Lemma

Let $\boldsymbol{v}=\left[\begin{array}{l}u \\ v\end{array}\right] \in T_{\boldsymbol{x}}(M)$ be a tangent vector at $\boldsymbol{x} \in U(p)$ with $u \neq 0$. Then we say that $\left|v u^{-1}\right|$ is the (absolute) slope of $\boldsymbol{v}$ and denote it by Slope $(\boldsymbol{v})$.

Consider any tangent vector $\boldsymbol{v}_{0}=\left[\begin{array}{l}1 \\ \delta\end{array}\right] \in T_{\boldsymbol{x}}(M)$ at $\boldsymbol{x}=(x+1, y) \in R_{\varepsilon}$ with $|\delta| \leq \varepsilon^{\frac{5}{2}}$. We set $\boldsymbol{v}_{0}^{\prime}=D \varphi(x+1, y)\left(\boldsymbol{v}_{0}\right)$ and $\boldsymbol{v}_{1}=D f^{u_{0}}(\varphi(x+1, y))\left(\boldsymbol{v}_{0}^{\prime}\right)$. By (2.8),

$$
\operatorname{Slope}\left(\boldsymbol{v}_{0}^{\prime}\right) \approx \frac{|d+e \delta|}{\left|3 c x^{2}+a \delta\right|}
$$

Since $\varepsilon \leq x$ and $|\delta| \leq \varepsilon^{\frac{5}{2}}$,

$$
\begin{aligned}
\operatorname{Slope}\left(\boldsymbol{v}_{0}^{\prime}\right) & \approx \frac{|d+e \delta|}{\left|3 c x^{2}+a \delta\right|} \leq \frac{|d|+|e \delta|}{\left|3 c x^{2}\right|-|a \delta|} \leq \frac{|d|+\left|e \varepsilon^{\frac{5}{2}}\right|}{\left|3 c \varepsilon^{2}\right|-\left|a \varepsilon^{\frac{5}{2}}\right|} \\
& =\frac{|d|+\left|e \varepsilon^{\frac{5}{2}}\right|}{|3 c|-\left|a \varepsilon^{\frac{1}{2}}\right|} \varepsilon^{-2}=\frac{|d|+\left|e \varepsilon^{\frac{5}{2}}\right|}{|3 c|-\left|a \varepsilon^{\frac{1}{2}}\right|} \varepsilon^{\frac{1}{2}} \cdot \varepsilon^{-\frac{5}{2}}
\end{aligned}
$$

By taking $\varepsilon_{1}>0$ sufficiently small, for any $0<\varepsilon \leq \varepsilon_{1}$, we have

$$
\text { Slope }\left(\boldsymbol{v}_{0}^{\prime}\right) \leq 2 \frac{|d|+\left|e \varepsilon^{\frac{5}{2}}\right|}{|3 c|-\left|a \varepsilon^{\frac{1}{2}}\right|} \varepsilon^{\frac{1}{2}} \cdot \varepsilon^{-\frac{5}{2}} \leq \frac{|3 d|}{|2 c|} \varepsilon^{\frac{1}{2}} \cdot \varepsilon^{-\frac{5}{2}} \leq 1 \cdot \varepsilon^{-\frac{5}{2}}
$$

Then, by (2.13) and (2.14), we have

$$
\operatorname{Slope}\left(\boldsymbol{v}_{1}\right)=\operatorname{Slope}\left(\boldsymbol{v}_{0}^{\prime}\right) \lambda^{u_{0}} \mu^{-u_{0}} \leq \varepsilon^{-\frac{5}{2}} \lambda^{u_{0}} \mu^{-u_{0}} \leq \varepsilon^{-\frac{5}{2}} \mu^{-\frac{5}{2} u_{0}} \leq \varepsilon^{-\frac{5}{2}} \varepsilon^{5}=\varepsilon^{\frac{5}{2}}
$$

Thus we get the following lemma. See Figure 2.7.


Figure 2.7: The tangent vectors $\boldsymbol{v}_{0}, \boldsymbol{v}_{0}^{\prime}$ and $\boldsymbol{v}_{1}$.

Lemma 2.3 (Slope Lemma I). Suppose that $f$ satisfies the conditions (2.13). Then there exists a constant $\varepsilon_{1}>0$ such that, if $\varepsilon \in\left(0, \varepsilon_{1}\right]$, then

$$
\begin{equation*}
\text { Slope }\left(\boldsymbol{v}_{0}^{\prime}\right) \leq \varepsilon^{-\frac{5}{2}} \quad \text { and } \quad \text { Slope }\left(\boldsymbol{v}_{1}\right) \leq \varepsilon^{\frac{5}{2}} \tag{2.26}
\end{equation*}
$$

for any tangent vector $\boldsymbol{v}_{0} \in T_{\boldsymbol{x}}(M)$ at $\boldsymbol{x}=(x+1, y) \in R_{\varepsilon}$ with $\operatorname{Slope}\left(\boldsymbol{v}_{0}\right) \leq \varepsilon^{\frac{5}{2}}$.
Fix a sufficiently small $s>0$ and set $\operatorname{pr}_{x}\left(S_{n}\right)=\left[s_{n}^{-}, s_{n}^{+}\right]$for $n \in \mathbb{N}$. If $n$ is sufficiently large, then $\left[s_{n}^{-}, s_{n}^{+}\right] \subset(0, s]$. Let $\beta_{n}^{u}(s)$ be the component of $\varphi\left(\alpha_{n}^{u}\right) \cap \operatorname{pr}_{x}^{-1}((0, s])$ containing $\varphi\left(\alpha_{n}\left(\left[\tilde{t}_{n,-}, \tilde{t}_{n,+}\right]\right)\right)$. For any $\boldsymbol{x} \in \beta_{n}^{u}(s)$, let $j_{n}(\boldsymbol{x})$ be a positive integer such that $f^{j}(\boldsymbol{x}) \in$ $U(p)$ for $j=1, \ldots, j_{n}(\boldsymbol{x})$ and $\operatorname{pr}_{x}\left(f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})\right) \in\left[1+\varepsilon,(1+\varepsilon)^{3}\right]$. For any $\varepsilon>0$, one can take $s$ so that $\operatorname{pr}_{x}\left(f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})\right) \in R_{\varepsilon}$ for any $\boldsymbol{x} \in \beta_{n}^{u}(s)$. Let $\boldsymbol{v}(\boldsymbol{x})$ be a unit vector tangent to $\beta_{n}^{u}(s)$ at $\boldsymbol{x}$.

The following result is applied to $f_{1}$ in the proof of Theorem 2.1.
Lemma 2.4 (Slope Lemma II). Let $\varepsilon_{1}$ be the constant given in Lemma 2.3. For any $\varepsilon \in\left(0, \varepsilon_{1}\right]$, there exist $s>0$ and $n_{0} \in \mathbb{N}$ such that

$$
\operatorname{Slope}\left(D f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})(\boldsymbol{v}(\boldsymbol{x}))\right)<\varepsilon^{\frac{5}{2}}
$$

if $n \geq n_{0}$ and $\boldsymbol{x} \in \beta_{n}^{u}(s) \backslash S_{n}$.

Proof. We only consider the case where $\boldsymbol{x}$ is an element of $\beta_{n}^{u}(s) \backslash S_{n}$ with $\mathrm{pr}_{x}(\boldsymbol{x}) \geq s_{n}^{+}$. Then $t \geq \tilde{t}_{n,+}$ holds if $\varphi\left(\alpha_{n}(t)\right)=\boldsymbol{x}$. The proof in the case of $\operatorname{pr}_{x}(\boldsymbol{x}) \leq s_{n}^{-}$is done quite similarly. Since $\rho_{n} \sim 1$ and $\tilde{t}_{n,+}=t_{n,+}+\rho_{n} \lambda^{\frac{n}{2}}, t-t_{n,+} \geq \tilde{t}_{n,+}-t_{n,+} \sim \lambda^{\frac{n}{2}}$. This implies that

$$
\begin{equation*}
t^{2}-t_{n,+}^{2} \sim t^{2} \succsim \lambda^{n} \tag{2.27}
\end{equation*}
$$

In fact, if $t-t_{n,+} \geq \frac{t}{2}$, then $t^{2}-t_{n,+}^{2}=\left(t-t_{n,+}\right)\left(t+t_{n,+}\right)>\frac{t^{2}}{2}$ and hence (2.27) holds. On the other hand, if $t-t_{n,+} \leq \frac{t}{2}$, then $t \leq 2 t_{n,+}$ and so $t \sim \lambda^{\frac{n}{2}}$. It follows that $t+t_{n,+} \sim \lambda^{\frac{n}{2}}$ and $t-t_{n,+} \sim \lambda^{\frac{n}{2}}$. Then $t^{2}-t_{n,+}^{2} \sim \lambda^{n} \sim t^{2}$. Thus (2.27) holds.

We set $\xi_{n}(t)=\operatorname{pr}_{x}(\boldsymbol{x})=\operatorname{pr}_{x}\left(\varphi\left(\alpha_{n}(t)\right)\right)$. By (2.17),

$$
\begin{align*}
& \xi_{n}(t)=a \tilde{y}_{n}(t)+b t \tilde{y}_{n}(t)+c t^{3}+\text { h.o.t. }  \tag{2.28}\\
& \xi_{n}^{\prime}(t)=a \tilde{y}_{n}^{\prime}(t)+b \tilde{y}_{n}(t)+b t \tilde{y}_{n}^{\prime}(t)+3 c t^{2}+\text { h.o.t.. }
\end{align*}
$$

From the definition of $j_{n}(\boldsymbol{x})$,

$$
\mu^{j_{n}(\boldsymbol{x})} \xi_{n}(t)=\mu^{j_{n}(\boldsymbol{x})} \operatorname{pr}_{x}(\boldsymbol{x})=\operatorname{pr}_{x}\left(f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})\right) \in\left[1+\varepsilon,(1+\varepsilon)^{3}\right] .
$$

This implies that $\mu^{j_{n}(\boldsymbol{x})} \xi_{n}(t) \sim 1$. We note that $\xi_{n}^{\prime}\left(t_{n,+}\right)=0$. By Mean Value Theorem, $\tilde{y}_{n}(t)-\tilde{y}_{n}\left(t_{n,+}\right)=\tilde{y}_{n}^{\prime}(c)\left(t-t_{n,+}\right)$ for some $t_{n,+}<c<t$. From this fact together with (2.16), (2.19), (2.27) and (2.28), we know that

$$
\xi_{n}^{\prime}(t)=\xi_{n}^{\prime}(t)-\xi_{n}^{\prime}\left(t_{n,+}\right) \sim t^{2}-t_{n,+}^{2} \sim t^{2} .
$$

By (2.18), Slope $(\boldsymbol{v}(\boldsymbol{x})) \sim t^{-2}$. Hence we have

$$
\begin{equation*}
\operatorname{Slope}\left(D f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})(\boldsymbol{v}(\boldsymbol{x}))\right)=\operatorname{Slope}(\boldsymbol{v}(\boldsymbol{x})) \cdot \frac{\lambda^{j_{n}(\boldsymbol{x})}}{\mu^{j_{n}(\boldsymbol{x})}} \sim t^{-2} \lambda^{j_{n}(\boldsymbol{x})} \xi_{n}(t) . \tag{2.29}
\end{equation*}
$$

Now we need to consider the following two cases.
Case 1. $c t^{3} \leq a \tilde{y}_{n}(t)$. By (2.28), $\xi_{n}(t) \sim \lambda^{n}$. Since $t^{-2} \precsim \lambda^{-n}$ by $t \succsim \lambda^{\frac{n}{2}}$, it follows from (2.29) that

$$
\operatorname{Slope}\left(D f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})(\boldsymbol{v}(\boldsymbol{x}))\right) \precsim \lambda^{-n} \lambda^{j_{n}(\boldsymbol{x})} \lambda^{n}=\lambda^{j_{n}(\boldsymbol{x})} .
$$

Case 2. $c^{3} \geq a \tilde{y}_{n}(t)$. Again by (2.28), we have $\xi_{n}(t) \sim t^{3}$. Then, by (2.29),

$$
\operatorname{Slope}\left(D f^{j_{n}(\boldsymbol{x})}(\boldsymbol{x})(\boldsymbol{v}(\boldsymbol{x}))\right) \sim t^{-2} \lambda^{j_{n}(\boldsymbol{x})} t^{3}=t \lambda^{j_{n}(\boldsymbol{x})} \precsim \lambda^{j_{n}(\boldsymbol{x})} .
$$

Let $n_{0}(s)$ be the minimum positive integer with $s_{n_{0}(s)}^{+}<s$. Since $n_{0}(s)$ goes to infinity as $s \rightarrow+0$, one can take $s=s(\varepsilon)>0$ such that our desired inequality holds for any $\boldsymbol{x} \in \beta_{n}^{u}(s) \backslash S_{n}$.

### 2.5 Sequence of rectangle-like boxes

Now we will define a sequence $\left\{B_{k, n}\right\}_{k=1}^{\infty}$ of rectangle-like boxes and estimate the sizes of them.

Recall that $\operatorname{pr}_{x}\left(S_{n}\right)=\left[s_{n}^{-}, s_{n}^{+}\right]$. Let $i_{n}$ be the positive integer with $(1+\varepsilon)^{2}<\mu^{i_{n}} s_{n}^{+} \leq$ $(1+\varepsilon)^{3}$. By (2.15), $f^{i_{n}}\left(S_{n}\right)$ is contained in $R_{\varepsilon}$ for any sufficiently large $n$. We set $f^{i_{n}}\left(S_{n}\right)=B_{1, n}=B_{1}$ for short. Since $s_{n}^{+} \sim \lambda^{n}$ by (2.22) and (2.23), we have

$$
\begin{equation*}
\mu^{i_{n}} \lambda^{n} \sim 1 \tag{2.30}
\end{equation*}
$$

We denote the width and height of $B_{1}$ and the distance between $B_{1}$ and $W_{\text {loc }}^{u}(p)$ by $W_{1, n}=W_{1}, H_{1, n}=H_{1}$ and $L_{1, n}=L_{1}$ respectively. It follows from (2.22), (2.23) and (2.25) that

$$
\begin{equation*}
W_{1, n} \sim \lambda^{\frac{3}{2} n} \mu^{i_{n}} \sim \lambda^{\frac{n}{2}}, \quad H_{1, n} \sim \lambda^{\frac{n}{2}+i_{n}}, \quad L_{1, n} \sim \lambda^{i_{n}} . \tag{2.31}
\end{equation*}
$$

Note that, for any sufficiently large $n, H_{1} \ll L_{1} \ll W_{1}$. Consider a closed interval $\delta_{1}$ in $W_{\text {loc }}^{u}(p)$ which is a small neighborhood of $\operatorname{pr}_{x}\left(B_{1}\right)$.

Let $v_{i}^{(1)}, e_{i}^{(1)}(i=0,1,2,3)$ be the vertices and edges of $B_{1}$ as illustrated in Figure 2.8 (a). We consider the image $\varphi\left(B_{1}\right)$. By Lemma 2.3 , for $i=0,2$,


Figure 2.8: The rectangle $B_{1}$ and the parallelogram-like box $B_{1}^{\prime}$.

$$
\operatorname{diam}\left(\operatorname{pr}_{x}\left(\varphi\left(e_{i}^{(1)}\right)\right)\right) \succsim \varepsilon^{\frac{5}{2}} W_{1} \sim \varepsilon^{\frac{5}{2}} \lambda^{\frac{n}{2}}
$$

On the other hand, for $i=1,3$,

$$
\operatorname{diam}\left(\operatorname{pr}_{x}\left(\varphi\left(e_{i}^{(1)}\right)\right)\right) \precsim H_{1} \sim \lambda^{\frac{n}{2}+i_{n}}
$$

Since $\lambda^{i_{n}} \varepsilon^{-\frac{5}{2}}$ can be supposed to be arbitrarily small for all sufficiently large $n$,

$$
\begin{align*}
x_{+}^{(1)}-x_{-}^{(1)} & \succsim \varepsilon^{\frac{5}{2}} W_{1}-O\left(\lambda^{\frac{n}{2}+i_{n}}\right) \sim \varepsilon^{\frac{5}{2}} \lambda^{\frac{n}{2}}-O\left(\lambda^{\frac{n}{2}+i_{n}}\right) \\
& =\varepsilon^{\frac{5}{2}} \lambda^{\frac{n}{2}}\left(1-\frac{O\left(\lambda^{i_{n}}\right)}{\varepsilon^{\frac{5}{2}}}\right) \sim \varepsilon^{\frac{5}{2}} \lambda^{\frac{n}{2}} . \tag{2.32}
\end{align*}
$$

where $x_{+}^{(1)}=\operatorname{pr}_{x}\left(\varphi\left(v_{1}^{(1)}\right)\right)$ and $x_{-}^{(1)}=\operatorname{pr}_{x}\left(\varphi\left(v_{3}^{(1)}\right)\right)$, see Figure $2.8(\mathrm{~b})$. Let $B_{1}^{\prime}$ be the intersection $\operatorname{pr}_{x}^{-1}\left(\left[x_{-}^{(1)}, x_{+}^{(1)}\right]\right) \cap \varphi\left(B_{1}\right)$. Any compact region in $U(p)$ like $B_{1}^{\prime}$ is called a parallelogram-like box.

Let $u_{1}$ be the positive integer with $(1+\varepsilon)^{2}<\mu^{u_{1}} x_{+}^{(1)} \leq(1+\varepsilon)^{3}$. By $(2.15), f^{u_{1}}\left(B_{1}^{\prime}\right)$ is contained in $R_{\varepsilon}$ for any sufficiently large $n$. We denote $f^{u_{1}}\left(B_{1}^{\prime}\right)$ by $B_{2}$. We call that any compact region in $U(p)$ like $B_{2}$ is a rectangle-like box.

Let $B$ be either a parallelogram-like or rectangle-like box. The horizontal width of $B$ is the diameter of the interval $\operatorname{pr}_{x}(B)$. The vertical height of $B$ is the maximum of the lengths of $\eta\left(x_{0}\right)$ with $x_{0} \in \operatorname{pr}_{x}(B)$, where $\eta\left(x_{0}\right)$ is the intersection of $B$ and the vertical line $x=x_{0}$. See Figure 2.9 in the case of $B=B_{1}^{\prime}$. Suppose that $B$ is a rectangle-like


Figure 2.9: A vertical segment $\eta\left(x_{0}\right)$ connecting the opposite pair of edges of the parallelogram-like box $B_{1}^{\prime}$.
box and $\delta$ is an almost horizontal arc in $U(q)$ with $B \cap \delta=\emptyset$ and $\operatorname{pr}_{x}(B) \subset \operatorname{pr}_{x}(\delta)$. Then the vertical distance between $B$ and $\delta$ is the maximum of $\sigma\left(x_{1}\right)$ with $x_{1} \in \operatorname{pr}_{x}(B)$, where $\sigma\left(x_{1}\right)$ is the length of the shortest segment in the vertical line $x=x_{1}$ connecting $B$ with $\delta$.

Let $\delta_{2}$ be a sub-arc of $f^{u_{1}}\left(\varphi\left(\delta_{1}\right)\right) \subset W^{u}(p)$ such that $\mathrm{pr}_{x}\left(\delta_{2}\right)$ is a small neighborhood of $\operatorname{pr}_{x}\left(B_{2}\right)$ in $W_{\text {loc }}^{u}(p)$. See Figure 2.10. We denote the horizontal width and vertical height of $B_{2}$ and the vertical distance between $B_{2}$ and $\delta_{2}$ by $W_{2}, H_{2}$ and $L_{2}$ respectively. By (2.14), (2.31) and (2.32),

$$
\begin{equation*}
W_{2}=\left(x_{+}^{(1)}-x_{-}^{(1)}\right) \mu^{u_{1}} \succsim \varepsilon^{\frac{5}{2}} \lambda^{\frac{n}{2}} \mu^{u_{1}} \geq \varepsilon^{-\frac{1}{2}} \tau_{1}^{-1} \lambda^{\frac{n}{2}} \sim \varepsilon^{-\frac{1}{2}} W_{1} . \tag{2.33}
\end{equation*}
$$

For any $x_{0}$ with $x_{-}^{(1)} \leq x_{0} \leq x_{+}^{(1)}, \eta\left(x_{0}\right)$ is a vertical segment connecting $\varphi\left(e_{0}^{(1)}\right)$ with $\varphi\left(e_{2}^{(1)}\right)$. By this fact together with (2.26), one can show the vertical height $H_{1}^{\prime}$ of $B_{1}^{\prime}$ satisfies $H_{1}^{\prime} \precsim H_{1} \varepsilon^{-\frac{5}{2}}$. It follows from (2.13) and (2.14) that

$$
\begin{equation*}
H_{2}=\lambda^{u_{1}} H_{1}^{\prime} \precsim \mu^{-\frac{3}{2} u_{1}} \varepsilon^{-\frac{5}{2}} H_{1}<\left(\varepsilon^{2}\right)^{\frac{3}{2}} \varepsilon^{-\frac{5}{2}} H_{1}=\varepsilon^{\frac{1}{2}} H_{1} . \tag{2.34}
\end{equation*}
$$

Let $L_{2}$ be the vertical distance between $B_{2}$ and $\delta_{2}$. By using an argument similar to that for the estimation (2.34), we have

$$
\begin{equation*}
L_{2} \precsim \varepsilon^{\frac{1}{2}} L_{1} . \tag{2.35}
\end{equation*}
$$

The following lemma is obtained immediately from (2.33), (2.34) and (2.35).
Lemma 2.5. Let $\varepsilon_{1}>0$ be the constant given in Lemma 2.3. Then there exists a constant $\varepsilon_{0} \in\left(0, \varepsilon_{1}\right]$ such that, for any $\varepsilon \in\left(0, \varepsilon_{0}\right)$, the inequalities

$$
W_{2} \geq 10 W_{1}, \quad H_{2} \leq 10^{-1} H_{1} \quad \text { and } \quad L_{2} \leq 10^{-1} L_{1}
$$

hold.
If $\mu=1+\varepsilon$ for an $\varepsilon \in\left(0, \varepsilon_{0}\right)$, then we say that $f$ satisfies the small expanding conditions at $p$.

We repeat the process as above. Let $B_{2}^{\prime}$ be the subset of $\varphi\left(B_{2}\right)$ cobounded by the vertical lines $x=x_{-}^{(2)}$ and $x=x_{+}^{(2)}$ passing through two of the four vertices of $\varphi\left(B_{2}\right)$ and satisfying $\left[x_{-}^{(2)}, x_{+}^{(2)}\right] \subset \operatorname{Int}\left(\operatorname{pr}_{x}\left(\varphi\left(B_{2}\right)\right)\right)$. Let $u_{2}$ be the positive integer with $(1+\varepsilon)^{2}<$ $\mu^{u_{2}} x_{+}^{(2)} \leq(1+\varepsilon)^{3}$ and $f^{u_{2}}\left(B_{2}^{\prime}\right) \subset R_{\varepsilon}$ for sufficient large $n \in \mathbb{N}$. Set $B_{3}=f^{u_{2}}\left(B_{2}^{\prime}\right)$. Let $\delta_{3}$ be a sub-arc of $f^{u_{2}}\left(\varphi\left(\delta_{2}\right)\right)$ such that $\operatorname{pr}_{x}\left(\delta_{3}\right)$ is a small neighborhood of $\operatorname{pr}_{x}\left(B_{3}\right)$ in $W_{\text {loc }}^{u}(p)$. We denote the horizontal width and vertical height of $B_{3}$ and the vertical distance between $B_{3}$ and $\delta_{3}$ by $W_{3}, H_{3}$ and $L_{3}$ respectively.

The objects $B_{k}^{\prime}, u_{k}, B_{k+1}, \delta_{k}, W_{k+1}, H_{k+1}, L_{k+1}(k=3,4,5, \ldots)$ are defined inductively if

$$
\begin{equation*}
B_{j} \subset R_{\varepsilon} \tag{2.36}
\end{equation*}
$$

for $j=1,2, \ldots, k$.
The top and bottom sides of the rectangle $B_{1}$ are horizontal and $\gamma_{1}=B_{1} \cap W^{u}(p)$ consists of three proper arcs in $B_{1}$. By Slope Lemma I (Lemma 2.3), for $k=2,3, \cdots$, the top and bottom sides of the rectangle-like box $B_{k}$ are almost horizontal and $\gamma_{k}=$ $B_{k} \cap W^{u}(p)$ consists of three proper arcs in $B_{k}$. See Figure 2.12. Thus we have the following lemma.
Lemma 2.6. Let $\varepsilon_{0}>0$ be the constant given in Lemma 2.5. For any $\varepsilon \in\left(0, \varepsilon_{0}\right]$, there exists the maximum integer $k_{0}=k_{0}(\varepsilon, n)$ satisfying (2.36). Moreover,

$$
\begin{equation*}
W_{k+1} \geq 10 W_{k}, \quad H_{k+1} \leq 10^{-1} H_{k} \quad \text { and } \quad L_{k+1} \leq 10^{-1} L_{k} \tag{2.37}
\end{equation*}
$$

hold for any $k=1,2, \ldots, k_{0}$.
See Figure 2.10 for the situation of Lemma 2.6. We note that, since $W_{1}=W_{1, n} \sim \lambda^{\frac{n}{2}}$ by (2.31), $\lim _{n \rightarrow \infty} k_{0}(\varepsilon, n)=\infty$ for a fixed $\varepsilon$ with $0<\varepsilon \leq \varepsilon_{0}$.


Figure 2.10: The pairs of the rectangle-like box $B_{k}$ and the sub-arc $\delta_{k}$ of $W^{u}(p)$ for $k=1,2, \ldots, k_{0}+1$.

### 2.6 Intersection Lemma

Recall that $\mathcal{C}$ is the codimension two submanifold of $\operatorname{Diff}^{3}(M)$ defined in Section 2.2. Let $f_{0}, f_{1}$ be elements of $\mathcal{C}$ satisfying the conditions (A1)-(A3) in Theorem 2.1. In particular, $\varepsilon>0$ is taken so that Slope Lemmas I and II (Lemmas 2.3 and 2.4) hold. Moreover, we suppose that $f_{0}, f_{1}$ satisfy the condition (2.9), which is one of the adaptable cases given in Section 2.8.

From now on, we set $f_{0}=f$ and use the notations in Sections 2.2-2.5. Here the subscription ' 0 ' is omitted from the notations. For example, $\lambda_{0}=\underline{\lambda}, \mu_{0}=\mu, p_{0}=p$, $q_{0}=q$ and so on. We also set $f_{1}=\bar{f}$ and represent the notations for $\bar{f}$ by adding bars to the corresponding notations for $f$, e.g. $\bar{\lambda}, \bar{\mu}, \bar{p}, \bar{q}, \bar{m}_{0}, \bar{S}_{n}, \bar{B}_{k}, \bar{W}_{k}$ and so on.

Let $h: M \rightarrow M$ be a homeomorphism with $\bar{f}=h \circ f \circ h^{-1}$. Here we note that $h(r)$ is not necessarily equal to $\bar{r}$. In fact, $h(r)=\bar{r}$ if and only if $m_{0}=\bar{m}_{0}$ or equivalently $\bar{\varphi}=h \circ \varphi \circ h^{-1}$. We may assume that $m_{0} \leq \bar{m}_{0}$ if necessary replacing $f$ and $\bar{f}$. Then $h\left(f^{m_{0}-m_{0}}(r)\right)=\bar{r}$. Since the constants appeared in (2.5) depend on the coordinate on $U(p)$, one can not replace the coordinates on $U(p)$ or $U(\bar{p})$ so as to satisfy $h(r)=\bar{r}$.

For any $C^{1}$ arc $\alpha$ in $U(p)$, the union of the end points of $\alpha$ is denoted by $\partial \alpha$. When any vector tangent to $\alpha$ is not vertical, the maximum $\operatorname{Slope}(\alpha)$ of $\operatorname{Slope}(\boldsymbol{v}(\boldsymbol{x}))$ for vectors $\boldsymbol{v}(\boldsymbol{x})$ tangent to $\alpha$ at $\boldsymbol{x} \in \alpha$ is well defined.

If $s>0$ is small enough, then $\gamma_{0, n}^{\prime}=\beta_{n}^{u}(s) \cap S_{n}$ is equal to $\alpha_{n}^{u} \cap S_{n}$ for any sufficiently large $n \in \mathbb{N}$.

The following is a key lemma for the proof of Theorem 2.1.
Lemma 2.7 (Intersection Lemma). Let $\gamma_{0, n}^{\prime}=\beta_{n}^{u}(s) \cap S_{n}$ and $\bar{\gamma}_{0, n}^{\prime}=\bar{\beta}_{n}^{u}(\bar{s}) \cap \bar{S}_{n}$. Then there exists an $n_{0} \in \mathbb{N}$ such that, for any $n \geq n_{0}$,

$$
\begin{equation*}
h\left(f^{\bar{m}_{0}-m_{0}}\left(\gamma_{0, n}^{\prime}\right)\right) \cap \bar{\gamma}_{0, n}^{\prime} \neq \emptyset \tag{2.38}
\end{equation*}
$$

Proof. We suppose that, for any $n_{0} \in \mathbb{N}$, there would exist $n>n_{0}$ such that

$$
\gamma_{0, n}^{\prime *} \cap \bar{\gamma}_{0, n}^{\prime}=\emptyset
$$

where $\gamma_{0, n}^{\prime *}=h \circ f^{\bar{m}_{0}-m_{0}}\left(\gamma_{0, n}^{\prime}\right)$, and introduce a contradiction.

Recall that $i_{n} \in \mathbb{N}$ satisfies $(1+\varepsilon)^{2}<f^{i_{n}}\left(s_{n}^{+}\right) \leq(1+\varepsilon)^{3}$ and $f^{i_{n}}\left(S_{n}\right) \subset R_{\varepsilon}$. For short, we set

$$
\gamma_{1}=\gamma_{1, n}:=f^{i_{n}}\left(\gamma_{0, n}^{\prime}\right) \quad \text { and } \quad \gamma_{1}^{*}=\gamma_{1, n}^{*}:=h\left(\gamma_{1}\right) .
$$

Then $\gamma_{1}^{*}=\bar{f}^{i_{n}-\left(\bar{m}_{0}-m_{0}\right)}\left(\gamma_{0, n}^{*}\right)$. Since $h(q)=\bar{q}$ and $\bar{f}=h \circ f \circ h^{-1}$, we have $h(1)=1, h(1+$ $\varepsilon)=1+\bar{\varepsilon}, h\left((1+\varepsilon)^{2}\right)=(1+\bar{\varepsilon})^{2}, h\left((1+\varepsilon)^{3}\right)=(1+\bar{\varepsilon})^{3}$ and $1+\bar{\varepsilon}<\operatorname{pr}_{x}\left(\gamma_{1}^{*}\right) \leq(1+\bar{\varepsilon})^{3}$. Strictly, $\gamma_{1}^{*}$ may slightly exceed $R_{\bar{\varepsilon}}$. Then we may rearrange our argument so that Lemmas 2.3 and 2.4 for $\bar{f}$ still hold if $\gamma_{1}^{*}$ is contained in a sufficiently small neighborhood of $R_{\bar{\varepsilon}}$. Then, by applying Lemma 2.4 to $\bar{f}$, one can show that $\gamma_{1}^{*}$ is a sub-arc almost parallel to $\bar{\delta}_{1} \subset W_{\text {loc }}^{u}(\bar{p})$ and Slope $\left(\gamma_{1}^{*}\right)<\bar{\varepsilon}^{\frac{5}{2}}$ for any sufficiently large $n$.

The intersection $\gamma_{1}^{\prime}=\varphi\left(\gamma_{1}\right) \cap B_{1}^{\prime}$ consists of mutually disjoint three arcs connecting the vertical sides of $B_{1}^{\prime}$. See Figure 2.11. We set $\gamma_{2}=f^{u_{1}}\left(\gamma_{1}^{\prime}\right)$ and $\gamma_{2}^{*}=h\left(\gamma_{2}\right)$. Note that


Figure 2.11: $\gamma_{1}^{\prime}$ is a disjoint union of arcs connecting the vertical sides of $B_{1}^{\prime}$ and $\gamma_{2}$ is a disjoint union of three proper arcs connecting the vartical sides of $B_{2}$.
$\gamma_{2}$ is a disjoint union of three proper arcs in $B_{2}$ connecting the vertical sides of $B_{2}$. Let $\widehat{\gamma}_{2}^{*}$ be the smallest arc in $W^{u}(\bar{p})$ containing $\gamma_{2}^{*}$. By applying Lemma 2.4 to $\bar{f}$, we have Slope $\left(\widehat{\gamma}_{2}^{*}\right)<\bar{\varepsilon}^{\frac{5}{2}}$. In particular, $\widehat{\gamma}_{2}^{*}$ is almost parallel to $\bar{\delta}_{2}$. Repeating the same argument, one can have sequences $\left\{\gamma_{k}\right\}$ satisfying the following conditions.

- Each $\gamma_{k}$ is a disjoint union of three proper arcs in $B_{k}$ connecting the vertical sides of $B_{k}$.
- For each $\gamma_{k}^{*}=h\left(\gamma_{k}\right)$, the smallest arc $\widehat{\gamma}_{k}^{*}$ in $W^{u}(\bar{p})$ containing $\gamma_{k}^{*}$ is almost parallel to $\bar{\delta}_{k}$.

See Figure 2.12.
Take $\boldsymbol{x} \in \gamma_{k}$ arbitrarily and set $\boldsymbol{x}^{*}=h(\boldsymbol{x}) \in \gamma_{k}^{*}$. Since $h$ is uniformly continuous on $R_{\varepsilon}$, for any $\bar{l}>0$, there exists $l>0$ independent of $\boldsymbol{x}$ such that $h \circ f^{\overline{m_{0}}-m_{0}}\left(N_{l}(\boldsymbol{x})\right) \subset N_{\bar{l}}\left(\boldsymbol{x}^{*}\right)$, where $N_{l}(\boldsymbol{x})$ is the $l$-neighborhood of $\boldsymbol{x}$ and $N_{\bar{l}}\left(\boldsymbol{x}^{*}\right)$ is the $\bar{l}$-neighborhood of $\boldsymbol{x}^{*}$ in $M$. If $n$ is sufficiently large, then $N_{l}(\boldsymbol{x})$ must intersect the three arcs of $\gamma_{k}$. However, $N_{\bar{l}}\left(\boldsymbol{x}^{*}\right)$ intersects only one arc of $\gamma_{k}^{*}$. This gives a contradiction. Thus (2.38) holds for all sufficiently large $n$.


Figure 2.12: $N_{l}(\boldsymbol{x})$ intersects the three $\operatorname{arcs}$ of $\gamma_{k}$, but $N_{\bar{l}}\left(\boldsymbol{x}^{*}\right)$ intersects only one arc of $\gamma_{k}^{*}$.

### 2.7 Proof of Theorem 2.1

Now we are ready to prove Theorem 2.1. The proof is done by using our Intersection Lemma (Lemma 2.7) together with arguments in [dM, Pa, Po] and so on. We only consider the case where both $f$ and $\bar{f}$ satisfy the condition (2.9), which belongs to Case $\mathrm{II}_{++}$in Section 2.8, and the small expanding condition at $p$ and $\bar{p}$ respectively. The proof of any other adaptable case is done similarly.

Proof of (M1) of Theorem 2.1. By Intersection Lemma (Lemma 2.7), one can take $\bar{r}_{n} \in$ $\bar{\gamma}_{0, n}^{\prime} \cap h \circ f^{\bar{m}_{0}-m_{0}}\left(\gamma_{0, n}^{\prime}\right)$. Since $\bar{r}_{n}$ converges to $\bar{r}$ as $n \rightarrow \infty, r_{n}=\left(h \circ f^{\bar{m}_{0}-m_{0}}\right)^{-1}\left(\bar{r}_{n}\right) \in \gamma_{0, n}^{\prime}$ converges to $r$ as $n \rightarrow \infty$. See Figure 2.13.


Figure 2.13: For $\bar{r}_{n} \in \bar{\gamma}_{0, n}^{\prime} \cap h \circ f^{\bar{m}_{0}-m_{0}}\left(\gamma_{0, n}^{\prime}\right), r_{n}=\left(h \circ f^{\bar{m}_{0}-m_{0}}\right)^{-1}\left(\bar{r}_{n}\right) \in \gamma_{0, n}^{\prime}$ converges to $r$ as $n \rightarrow \infty$.

Let $W_{\text {loc, },}^{u}(p)$ be the component of $W_{\text {loc }}^{u}(p) \backslash\{p\}$ containing $q$. Take a fundamental
domain $D$ for $f$ in $W_{\text {loc },+}^{u}(p)$. Then there exist subsequences $\left\{r_{n(k)}\right\} \subset\left\{r_{n}\right\},\{m(k)\}$ of $\mathbb{N}$ and $\boldsymbol{x}_{0} \in D$ satisfying the following conditions.

- $r_{n(k)}$ converges to $r$ as $k \rightarrow \infty$.
- $\boldsymbol{x}_{n(k)}:=f^{m(k)}\left(r_{n(k)}\right)$ converges to $\boldsymbol{x}_{0}=\left(x_{0}, 0\right)$ as $k \rightarrow \infty$.
- $q_{n(k)}:=\varphi^{-1}\left(r_{n(k)}\right)$ converges to $q$ as $k \rightarrow \infty$.

Then

$$
\begin{aligned}
x_{0} & =\lim _{k \rightarrow \infty} \operatorname{pr}_{x}\left(\boldsymbol{x}_{n(k)}\right)=\lim _{k \rightarrow \infty} \operatorname{pr}_{x}\left(r_{n(k)}\right) \mu^{m(k)} \\
& =\lim _{k \rightarrow \infty} a z_{0}\left(\lambda^{n(k)}+O\left(\lambda^{\frac{3}{2} n(k)}\right)\right) \mu^{m(k)} \\
& =\lim _{k \rightarrow \infty} a z_{0} \lambda^{n(k)} \mu^{m(k)}
\end{aligned}
$$

It follows that $\lim _{k \rightarrow \infty} \lambda^{n(k)} \mu^{m(k)}=\frac{x_{0}}{a z_{0}}$. Then there exist constants $C_{0}$ and $C_{1}$ with $0<C_{0}<C_{1}$ and such that

$$
C_{0}<\lambda^{n(k)} \mu^{m(k)}<C_{1}
$$

for any $k$. Taking the logarithms of this inequalities, we have

$$
\frac{\log C_{0}}{n(k) \log \mu}<\frac{\log \lambda}{\log \mu}+\frac{m(k)}{n(k)}<\frac{\log C_{1}}{n(k) \log \mu}
$$

This shows that $\lim _{k \rightarrow \infty} \frac{m(k)}{n(k)}=-\frac{\log \lambda}{\log \mu}$. By applying a similar argument to $\bar{f}$, one can prove

$$
\lim _{k \rightarrow \infty} \bar{\lambda}^{n(k)} \bar{\mu}^{m(k)-\left(\bar{m}_{0}-m_{0}\right)}=\frac{h\left(x_{0}\right)}{\bar{a} \bar{z}_{0}}
$$

and hence $\lim _{k \rightarrow \infty} \frac{m(k)}{n(k)}=\lim _{k \rightarrow \infty} \frac{m(k)-\left(\bar{m}_{0}-m_{0}\right)}{n(k)}=-\frac{\log \bar{\lambda}}{\log \bar{\mu}}$. Consequently, $\frac{\log \lambda}{\log \mu}=\frac{\log \bar{\lambda}}{\log \bar{\mu}}$ holds.

Lemma 2.8. If $\frac{\log \lambda}{\log \mu}$ is irrational, then the restriction $\left.h\right|_{W_{+}^{u}(p)}$ is locally $C^{1}$ diffeomorphic, where $W_{+}^{u}(p)$ is the component of $W^{u}(p) \backslash\{p\}$ containing $q$.

Proof. Let $s_{n}$ be the real number with $\operatorname{pr}_{x}\left(r_{n}\right)=\mu^{-s_{n}}$. Since $\operatorname{pr}_{x}\left(r_{n}\right) \approx a z_{0}\left(\lambda^{n}+O\left(\lambda^{\frac{3}{2} n}\right)\right)$ by (2.17) and (2.21), we have

$$
1=\operatorname{pr}_{x}\left(r_{n}\right) \mu^{s_{n}} \approx a z_{0}\left(\lambda^{n}+O\left(\lambda^{\frac{3}{2} n}\right)\right) \mu^{s_{n}} \approx a z_{0} \lambda^{n} \mu^{s_{n}}
$$

Thus $c_{n}=a z_{0} \lambda^{n} \mu^{s_{n}}$ satisfies $\lim _{n \rightarrow \infty} c_{n}=1$. Moreover,

$$
\begin{equation*}
s_{n}=\frac{\log c_{n}}{\log \mu}-\frac{\log \left(a z_{0}\right)}{\log \mu}-n \frac{\log \lambda}{\log \mu} \tag{2.39}
\end{equation*}
$$

Since $-\frac{\log \lambda}{\log \mu}$ is irrational, the set

$$
\left\{-\frac{\log \left(a z_{0}\right)}{\log \mu}-n \frac{\log \lambda}{\log \mu} \bmod 1 ; n=1,2, \ldots\right\}
$$

is dense in the interval $[0,1]$. Since $\lim _{n \rightarrow \infty} \log c_{n}=0$, the set $S=\left\{s_{n} \bmod 1 ; n=\right.$ $1,2, \ldots\}$ is also dense in $[0,1]$.

Take a point $x_{0}$ of $\left[\mu^{-1}, 1\right]$ arbitrarily, and let $\sigma \in[0,1]$ be the real number with $\mu^{-\sigma}=$ $x_{0}$. Since $\left[\mu^{-1}, 1\right]$ is a fundamental domain for $f$ in $W_{\text {loc, }+}^{u}(p)$, it follows from the density of $S$ that there exist subsequences $\{n(k)\},\{m(k)\}$ of $\mathbb{N}$ such that $\lim _{k \rightarrow \infty}\left(s_{n(k)}-m(k)\right)=\sigma$. Then

$$
\begin{aligned}
x_{0} & =\mu^{-\sigma}=\lim _{k \rightarrow \infty} \mu^{-s_{n(k)}+m(k)}=\lim _{k \rightarrow \infty} \operatorname{pr}_{x}\left(r_{n(k)}\right) \mu^{m(k)} \\
& =\lim _{k \rightarrow \infty} a z_{0}\left(\lambda^{n(k)}+O\left(\lambda^{\frac{3}{2} n(k)}\right)\right) \mu^{m(k)}=\lim _{k \rightarrow \infty} a z_{0} \lambda^{n(k)} \mu^{m(k)} .
\end{aligned}
$$

Thus we have $\lim _{k \rightarrow \infty} \lambda^{n(k)} \mu^{m(k)}=\frac{x_{0}}{a z_{0}}$.
Since $\bar{f}$ is conjugate to $f$ via $h, \operatorname{pr}_{x}\left(\bar{r}_{n(k)}\right) \bar{\mu}^{m(k)-\left(\bar{m}_{0}-m_{0}\right)}$ converges to $h\left(x_{0}\right)$. As above, we have

$$
\lim _{k \rightarrow \infty} \bar{\lambda}^{n(k)} \bar{\mu}^{m(k)-\left(\bar{m}_{0}-m_{0}\right)}=\frac{h\left(x_{0}\right)}{\bar{a} \bar{z}_{0}} .
$$

If we set $\tau=\frac{\log \bar{\mu}}{\log \mu}=\frac{\log \bar{\lambda}}{\log \lambda}$, then $\bar{\mu}=\mu^{\tau}$ and $\bar{\lambda}=\lambda^{\tau}$. It follows that

$$
\frac{x_{0}^{\tau}}{a^{\tau} z_{0}^{\tau}}=\frac{h\left(x_{0}\right)}{\bar{a} \bar{z}_{0}} \bar{\mu}^{\bar{m}_{0}-m_{0}} .
$$

Thus $\left.h\right|_{W_{\text {loc },+}^{u}(p)}$ is a $C^{1}$ diffeomorphism represented as

$$
h(x)=\frac{\bar{a} \bar{z}_{0}}{a^{\tau} z_{0}^{\tau} \bar{\mu}^{\bar{m}_{0}-m_{0}}} x^{\tau},
$$

where $W_{\text {loc },+}^{u}(p)$ is the component of $W_{\text {loc }}^{u}(p) \backslash\{p\}$ containing $q$. Since $W_{+}^{u}(p)=\bigcup_{n=0}^{\infty} f^{n}\left(W_{\text {loc, },+}^{u}(p)\right)$ and both $f$ and $\bar{f}$ are $C^{3}$ diffeomorphisms, $\left.h\right|_{W_{+}^{u}(p)}$ is locally $C^{1}$ diffeomorphic. This completes the proof.

Proof of (M2) of Theorem 2.1. Take a sequence $\left\{q_{j}\right\}$ on $W_{\text {loc, }+}^{u}(p)$ converging to $q$ and set $t_{j}=\varphi\left(q_{j}\right)$. See Figure 2.14. Let $t_{j}^{\prime}$ be the image of $t_{j}$ by the horizontal projection to $W_{\text {loc }}^{s}(p)$. Obviously, both $t_{j}$ and $t_{j}^{\prime}$ converge to $r$ as $k \rightarrow \infty$. There exist subsequences $\left\{t_{j(k)}\right\}$ of $\left\{t_{j}\right\},\{l(k)\}$ of $\mathbb{N}$ and a point $x_{1}$ of $W_{\text {loc }}^{u}(p)$ with $\lim _{k \rightarrow \infty} f^{l(k)}\left(t_{j(k)}\right)=x_{1}$. Then the following approximations

$$
x_{1} \sim \operatorname{pr}_{x}\left(t_{j(k)}\right) \mu^{l(k)} \sim\left[d\left(t_{j(k)}^{\prime}, r\right)\right]^{3} \mu^{l(k)} \sim\left[d\left(q_{j(k)}, q\right)\right]^{3} \mu^{l(k)}
$$



Figure 2.14: The case of $\mathrm{II}_{++}$.
hold. It follows that $\mu^{-l(k)} \sim\left[d\left(q_{j(k)}, q\right)\right]^{3}$. Similarly, we have $\bar{\mu}^{-l(k)+\left(\bar{m}_{0}-m_{0}\right)} \sim\left[d\left(\bar{q}_{j(k)}, \bar{q}\right)\right]^{3}$, where $\bar{q}_{j(k)}=h\left(q_{j(k)}\right)$. Since $\left.h\right|_{W_{+}^{u}(p)}$ is locally $C^{1}$-diffeomorphic by Lemma 2.8,

$$
d\left(\bar{q}_{j(k)}, \bar{q}\right) \sim d\left(q_{j(k)}, q\right) .
$$

Thus

$$
\left(\frac{\bar{\mu}}{\mu}\right)^{-l(k)} \sim\left(\frac{d\left(\bar{q}_{j(k)}, \bar{q}\right)}{d\left(q_{j(k)}, q\right)}\right)^{3} \bar{\mu}^{-\left(\bar{m}_{0}-m_{0}\right)} \sim 1
$$

This implies that $\mu=\bar{\mu}$. By (M1), we also have $\lambda=\bar{\lambda}$. This completes the proof of the part (M2).

Remark 2.9. Some arguments used in the case that the tangency between $W^{s}(p)$ and $W^{u}(p)$ is one-sided (for example [ $\left.\mathrm{dM}, \mathrm{Pa}, \mathrm{Po}\right]$ ) can not be applicable to the two-sided case. Here we explain the reason.

Suppose that a homoclinic tangency $q_{0}$ is one-sided, say a quadratic tangency. Take an arc $\gamma$ in $U\left(q_{0}\right)$ meeting $W_{\text {loc }}^{u}\left(p_{0}\right)$ orthogonally at $q_{0}$. Let $\left\{w_{i}\right\}$ be a sequence in $\gamma$ converging to $q_{0}$ from above. Then

$$
\begin{equation*}
d\left(w_{i}, W^{s}\left(p_{0}\right)\right) \approx d\left(w_{i}, W_{\mathrm{loc}}^{u}\left(p_{0}\right)\right) \tag{2.40}
\end{equation*}
$$

holds. On the other hand, their images by the conjugacy homeomorphism $h$ satisfy

$$
\begin{equation*}
d\left(h\left(w_{i}\right), W^{s}\left(p_{1}\right)\right) \leq d\left(h\left(w_{i}\right), W_{\mathrm{loc}}^{u}\left(p_{1}\right)\right) . \tag{2.41}
\end{equation*}
$$

See Figure 2.15 (a). By using (2.40) and (2.41), one can show that $\frac{\log \lambda_{1}}{\log \mu_{1}} \leq \frac{\log \lambda_{0}}{\log \mu_{0}}$. By

(a)


Figure 2.15: (a) The case of quadratic tangencies. (b) The case of cubic tangencies.
applying the same argument to $h^{-1}$, we also have $\frac{\log \lambda_{1}}{\log \mu_{1}} \geq \frac{\log \lambda_{0}}{\log \mu_{0}}$, and hence $\frac{\log \lambda_{1}}{\log \mu_{1}}=$ $\frac{\log \lambda_{0}}{\log \mu_{0}}$.

Now we consider the case of two-sided tangencies, say cubic tangencies, and $\left\{w_{i}\right\}$ is a sequence as above. Then the approximation (2.40) still holds. However, the inequality (2.41) would not hold as is suggested in Figure 2.15 (b). So it might be difficult to get the inequality $\frac{\log \lambda_{1}}{\log \mu_{1}} \leq \frac{\log \lambda_{0}}{\log \mu_{0}}$ only by arguments in [dM, Pa, Po]. Thus we need another idea in the study of moduli associated with two-sided homoclinic tangencies.

### 2.8 Adaptable conditions

In this section, we will present conditions on the signs of $a, b c, \lambda$ and $\mu$ under which any arguments presented throughout the previous sections are valid.

Recall that we have set

$$
U(p)=[-2,2] \times[-2,2], W_{\mathrm{loc}}^{u}(p)=[-2,2] \times\{0\}, W_{\mathrm{loc}}^{s}(p)=\{0\} \times[-2,2] .
$$

The union $W_{\text {loc }}^{u}(p) \cup W_{\text {loc }}^{s}(p)$ divides $U(p)$ to four components. The closures of these components containing $(1,1),(-1,1),(-1,-1)$ and $(1,-1)$ are called the first, second, third and fourth quadrants of $U(p)$ and denoted by $Q_{1}, Q_{2}, Q_{3}$ and $Q_{4}$, respectively. In our argument it is required that $\varphi\left(R_{\varepsilon}\right)$ or some substitution is in $Q_{1}$. If $R_{\varepsilon}$ lies in $Q_{2}$, then we may use

$$
R_{\varepsilon}^{-}=\left[(1+\varepsilon)^{-3},(1+\varepsilon)^{-1}\right] \times\left[0, \varepsilon^{3}\right]
$$

instead of $R_{\varepsilon}$. Then $\varphi\left(R_{\varepsilon}^{-}\right)$is in $Q_{1}$. See Figure 2.16. Thus one can arrange the placement of $\varphi\left(R_{\varepsilon}\right)$ suitably under any conditions on the signs of $a, b c, \lambda$ and $\mu$.


Figure 2.16: (1) $\varphi\left(R_{\varepsilon}\right)$ is in $Q_{1}$. (2) $\varphi\left(R_{\varepsilon}^{-}\right)$is in $Q_{1}$.

Definition 2.10 (Adaptable condition). $f$ satisfies the adaptable condition with respect to $(p, q)$ if, for all sufficiently large positive integers $n$ (or positive even or odd integers), there exists a rectangle $S_{n}$ defined as in Section 2.3 and either $S_{n}$ or its image $f\left(S_{n}\right)$ lies in $Q_{1}$.

As was seen in Section $2.3, S_{n}$ exists if and only if there exists $t_{n}$ satisfying the condition

$$
3 c t_{n}^{2} \approx-b \lambda^{n} z_{0}
$$

which corresponds to (2.20). Here $z_{0}$ is the positive constant as illustrated in Figure 2.5.
Now we will see that the existence of $S_{n}$ and the placements of $S_{n}$ and $f\left(S_{n}\right)$ are strictly determined by the signs of $a, b c, \lambda$ and $\mu$, which are classified to the sixteen cases as in Table 2.1

First we suppose that $\lambda>0$. Then there exists $t_{n}$ satisfying (2.20') if and only if $b c<0$. Moreover, if $a>0$, then $S_{n}$ is in $Q_{1}$, which belongs to Case $\mathrm{II}_{+}$. See Figure 2.17 (1). If $a<0$, then $S_{n}$ is in $Q_{2}$. Hence $f\left(S_{n}\right)$ is $Q_{1}$ if $\mu<0$, which is in Case IV ${ }_{+-}$. See Figure 2.17 (2).

(1)


Figure 2.17: (1) The case of $\mathrm{II}_{+}$or $\mathrm{II}_{-}$. (2) The case of $\mathrm{IV}_{+-}$or $\mathrm{IV}_{-}$.

| Case |  |  | $a$ | $b c$ | $\lambda$ |  | $\mu$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| I | $\mathrm{I}_{+}$ | $\mathrm{I}_{++}$ | + | + |  |  | + |
|  |  | $\mathrm{I}_{+-}$ |  |  |  |  | - |
|  | I_ | $\mathrm{I}_{-+}$ |  |  | - |  | + |
|  |  | I |  |  |  |  | - |
| II | $\mathrm{II}_{+}$ | II | + | - | + |  | + |
|  |  | H. ${ }_{\text {. }}$. |  |  |  |  | - |
|  | II_ | 11. |  |  | - |  | + |
|  |  | II. ${ }^{\text {a }}$. |  |  |  |  | - |
| III | $\mathrm{III}_{+}$ | $\mathrm{III}_{++}$ | - | + | + |  | + |
|  |  | $\mathrm{III}_{+-}$ |  |  |  |  | - |
|  | III- | III |  |  | - |  | + |
|  |  | III |  |  |  |  | - |
| IV | $\mathrm{IV}_{+}$ | $\mathrm{IV}_{++}$ | - | - | $+$ |  | + |
|  |  | IV |  |  |  |  | - |
|  | IV_ | IV ${ }_{-+}$ |  |  | - |  | + |
|  |  | IV |  |  |  |  | - |

Table 2.1: The shaded cells are the cases in which $f$ satisfies the adaptable conditions.

Next we suppose that $\lambda<0$. Then there exists $t_{n}$ satisfying (2.20') if and only if either (i) $b c<0$ and $n$ is even or (ii) $b c>0$ and $n$ is odd. In the case (i), $S_{n}$ is in $Q_{1}$ if $a>0$, which belongs to Case II_. See Figure 2.17 (1). If $a<0$ and $\mu<0$, then $f\left(S_{n}\right)$ is in $Q_{1}$, which belongs to Case IV__. See Figure 2.17 (2). On the other hand, in the case (ii), $S_{n}$ is in $Q_{1}$ if $a<0$, which belongs to Case III_. See Figure 2.18(1). If $a>0$ and $\mu<0$, then $f\left(S_{n}\right)$ is in $Q_{1}$, which belongs to Case I__. See Figure 2.18 (2).


Figure 2.18: (1) The case of III_. (2) The case of I_-.
Thus we have the following proposition.
Proposition 2.11. If one of Cases I__ , II, $^{\text {III_, }}$ IV $_{+-}$and IV_- holds, then $f$ satisfies the adaptable condition with respect to $(p, q)$.

It follows from the proposition that $f$ satisfies the adaptable condition in nine of the sixteen cases in Table 2.1.

## Chapter 3

## Moduli of 3-dimensional diffeomorphisms with saddle foci

In this chapter, we investigate moduli of a 3 -dimensional diffeomorphism $f$ with a sadldle focus $p$ and a homoclinic quadratic tangency $q$ associated with $p$. We show that, for most of such diffeomorphisms, all the eigenvalues of $D f(p)$ are moduli and the restriction of a conjugacy homeomorphism to a local unstable manifold is a uniquely determined linear conformal map.

### 3.1 Moduli of 3-dimensional diffeomorphisms with saddle foci

First, we prove the following theorem.
Theorem 3.1. Let $M$ be a 3-manifold and $f_{j}(j=0,1)$ elements of $\operatorname{Diff}^{r}(M)$ for some $r \geq$ 3 which have hyperbolic fixed points $p_{j}$ and homoclinic quadratic tangencies $q_{j}$ positively associated with $p_{j}$ and satisfy the following conditions.

- For $j=0,1$, there exists a neighborhood $U\left(p_{j}\right)$ of $p_{j}$ in $M$ such that $\left.f_{j}\right|_{U\left(p_{j}\right)}$ is linear and $D f_{j}\left(p_{j}\right)$ has non-real eigenvalues $r_{j} e^{ \pm \sqrt{-1} \theta_{j}}$ and a real eigenvalue $\lambda_{j}$ with $r_{j}>1$, $\theta_{j} \neq 0 \bmod \pi$ and $0<\lambda_{j}<1$.
- $f_{0}$ is topologically conjugate to $f_{1}$ by a homeomorphism $h: M \rightarrow M$ with $h\left(p_{0}\right)=p_{1}$ and $h\left(q_{0}\right)=q_{1}$.

Then the following (D1) and (D2) hold.
(D1) $\frac{\log \lambda_{0}}{\log r_{0}}=\frac{\log \lambda_{1}}{\log r_{1}}$.
(D2) Either $\theta_{0}=\theta_{1}$ or $\theta_{0}=-\theta_{1} \bmod 2 \pi$.
Here we say that a homoclinic quadratic tangency $q_{0}$ is positively associated with $p_{0}$ if both $f_{0}^{n}\left(q_{0}\right)$ and $f_{0}^{-n}(\alpha)$ lie in the same component of $U\left(p_{0}\right) \backslash W_{\text {loc }}^{u}\left(p_{0}\right)$ for a sufficiently
large $n \in \mathbb{N}$ and any small curve $\alpha$ in $W^{s}\left(p_{0}\right)$ containing $q_{0}$. Theorem 3.1 holds also in the case when $\theta_{0}=0 \bmod \pi$ or $-1<\lambda_{j}<0$ except for some rare case, see Remark 3.4 for details.

Remark 3.2. Assertion (D1) of Theorem 3.1 is implied in the case (D) of Theorem 1.1 in [NPT, Chapter III]. Assertion (D2) is also proved by Dufraine [Du2] under weaker assumptions. The author used non-spiral curves in $W_{\text {loc }}^{u}(p)$ emanating from $p$. On the other hand, we employ unstable bent disks defined in Section 3.2 which are originally introduced by Nishizawa [Ni]. By using such disks, we construct a convergent sequence of mutually parallel straight segments in $W_{\text {loc }}^{u}(p)$ which are mapped to straight segments in $W_{\text {loc }}^{u}(h(p))$ by $h$, see Figure 3.9. An advantage of our proof is that these sequences are applicable to prove our main theorem, Theorem 3.3 below.

Results corresponding to Theorem 3.1 for 3-dimensional flows with Shilnikov cycles are obtained by Togawa [To], Carvalho-Rodrigues [CR] and for those with connections of saddle-foci by Bonatti-Dufraine [BD], Dufraine [Du1], Rodrigues [Rod] and so on. See the Section 2 in [Rod] for details. Moreover Carvalho-Rodrigues $[\mathrm{CR}]$ present results on moduli of 3 -dimensional flows with Bykov cycles.

The following theorem is the main theorem in this chapter.
Theorem 3.3. Under the assumptions in Theorem 3.1, suppose moreover that $\theta_{0} / 2 \pi$ is irrational. Then the following conditions hold.
(E1) $\lambda_{0}=\lambda_{1}$ and $r_{0}=r_{1}$.
(E2) The restriction $\left.h\right|_{W_{\mathrm{loc}}\left(p_{0}\right)}: W_{\mathrm{loc}}^{u}\left(p_{0}\right) \rightarrow W_{\mathrm{loc}}^{u}\left(p_{1}\right)$ is a uniquely determined linear conformal map.

In contrast to Posthumus' results for 2-dimensional diffeomorphisms, the eigenvalues $\lambda_{0}$ and $r_{0}$ are proved to be moduli without the assumption that $\frac{\log \lambda_{0}}{\log r_{0}}$ is irrational.

The restriction $\left.h\right|_{W_{\text {loc }}^{u}\left(p_{0}\right)}$ is said to be a linear conformal map if $\left.h\right|_{W_{\text {loc }}\left(p_{0}\right)}$ is represented as $\left.h\right|_{W_{\text {loc }}^{u}\left(p_{0}\right)}(z)=\rho e^{\sqrt{-1} \omega} z\left(z \in W_{\text {loc }}^{u}\left(p_{0}\right)\right)$ for some $\rho \in \mathbb{R} \backslash\{0\}$ and $\omega \in \mathbb{R}$ under the natural identification of $W_{\text {loc }}^{u}\left(p_{0}\right), W_{\text {loc }}^{u}\left(p_{1}\right)$ with neighborhoods of the origin in $\mathbb{C}$ via their linearizing coordinates.

For any $r_{j}>1$ and $\theta_{j} \in \mathbb{R}(j=0,1)$, let $\varphi_{j}: \mathbb{C} \rightarrow \mathbb{C}$ be the map defined by $\varphi_{j}(z)=$ $r_{j} e^{\sqrt{-1} \theta_{j}} z$. Then there are many choices of conjugacy homeomorphisms on $\mathbb{C}$ for $\varphi_{0}$ and $\varphi_{1}$. For example, we take two-sided Jordan curves $\Gamma_{j}$ in $\mathbb{C}$ with $\varphi_{j}\left(\Gamma_{j}\right) \cap \Gamma_{j}=\emptyset$ and bounding disks in $\mathbb{C}$ containing the origin arbitrarily. Then there exists a conjugacy homeomorphism $h: \mathbb{C} \rightarrow \mathbb{C}$ for $\varphi_{0}$ and $\varphi_{1}$ with $h\left(\Gamma_{0}\right)=\Gamma_{1}$. On the other hand, Theorem 3.3 (E2) implies that we have severe constraints in the choice of conjugacy homeomorphisms for 3 -dimensional diffeomorphisms as above. Intuitively, it says that only a homeomorphism $h$ with $\left.h\right|_{W_{\text {loc }}^{u}}(p)$ linear and conformal can be a candidate for a conjugacy between $f_{0}$ and $f_{1}$. As an application of the linearity and conformality of $\left.h\right|_{W_{\text {loc }}^{u}} ^{u}(p)$, we will present a new modulus for $f_{0}$ other than $\theta_{0}, \lambda_{0}, r_{0}$, see Corollary 3.9 in Section 3.5.

### 3.2 Front curves and folding curves

For $j=0,1$, let $f_{j}$ be a diffeomorphism and $q_{j}$ a quadratic tangency associated with a hyperbolic fixed point $p_{j}$ satisfying the conditions of Theorem 3.1. We will define in this section front curves in $W^{u}\left(p_{j}\right)$ and folding curves in $W_{\text {loc }}^{u}\left(p_{j}\right)$ and show in the next section that these curves converge to straight segments which are preserved by any conjugacy homeomorphism between $f_{0}$ and $f_{1}$.

We set $f_{0}=f, p_{0}=p, q_{0}=q, r_{0}=r, \theta_{0}=\theta$ and $\lambda_{0}=\lambda$ for short. Similarly, let $f_{1}=f^{\prime}, p_{1}=p^{\prime}, q_{1}=q^{\prime}, r_{1}=r^{\prime}, \theta_{1}=\theta^{\prime}$ and $\lambda_{1}=\lambda^{\prime}$. Suppose that $(z, t)=(x, y, t)$ with $z=x+\sqrt{-1} y$ is a coordinate around $p$ with respect to which $f$ is linear. For a small $a>0$, let $D_{a}(p)$ be the disk $\{z \in \mathbb{C} ;|z| \leq a\}$. We may assume that $q$ is contained in the interior of $D_{a}(p) \times\{0\} \subset W_{\text {loc }}^{u}(p)$ and $\widehat{q}=f^{N}(q)$ is in the interior of the upper half $W_{\text {loc }}^{s+}(p)=\{0\} \times[0, a]$ of $W_{\text {loc }}^{s}(p)$ for some $N \in \mathbb{N}$. See Figure 3.1. Let $U_{a}(p)$ be


Figure 3.1: A saddle-focus $p$ and a homoclinic quadratic tangency $q$ in $D_{a}(p)$.
the circular column in the coordinate neighborhood defined by $U_{a}(p)=D_{a}(p) \times[0, a]$ and $V_{\widehat{q}}$ a small neighborhood of $\widehat{q}$ in $U_{a}(p)$. Suppose that $U_{a}(p)$ has the Euclidean metric induced from the linearizing coordinate on $U_{a}(p)$. By choosing the coordinate suitably and replacing $\theta$ by $-\theta$ if necessary, we may assume that the restriction $\left.f\right|_{D_{a}(p)}$ is represented as $r e^{\sqrt{-1} \theta} z$ for $z \in \mathbb{C}$ with $|z|<a$. Similarly, one can suppose that $\left.f^{\prime}\right|_{D_{a^{\prime}}\left(p^{\prime}\right)}$ is represented as $r^{\prime} e^{\sqrt{-1} \theta^{\prime}} z$ for some $a^{\prime}>0$. The orthogonal projection pr : $U_{a}(p) \rightarrow D_{a}(p)$ is defined by $\operatorname{pr}(x, y, t)=(x, y)$.

In this section, we construct an unstable bent disk $\widetilde{H}_{0}$ in $W^{u}(p) \cap U_{a}(p)$, the front curve $\widetilde{\gamma}_{0}$ in $\widetilde{H}_{0}$ and the folding curves $\gamma_{0}$ in $U_{a}(p)$. We also define the sequence of unstable bent disks $\widetilde{H}_{m}$ in $W^{u}(p) \cap U_{a}(p)$ converging to $\widetilde{H}_{0}$, which will be used in the next section to construct the sequence of front curves converging to $\widetilde{\gamma}_{0}$.

### 3.2.1 Construction of unstable bent disks, front curves and folding curves

We set $\widehat{q}=\left(0, t_{0}\right)$. Let $\widetilde{H}$ be the component of $W^{u}(p) \cap V_{\widehat{q}}$ containing $\widehat{q}$. One can retake the linearizing coordinate on $\mathbb{C}$ if necessary so that the line in $V_{\widehat{q}}$ passing through $\widehat{q}$ and parallel to the $x$-axis in $U_{a}(p)$ meets $\widetilde{H}$ transversely. Then $\widetilde{H}$ is represented as the graph of a $C^{r}$ function $x=\varphi(y, t)$ with

$$
\begin{equation*}
\varphi\left(0, t_{0}\right)=0, \quad \frac{\partial \varphi}{\partial t}\left(0, t_{0}\right)=0 \quad \text { and } \quad \frac{\partial^{2} \varphi}{\partial t^{2}}\left(0, t_{0}\right) \neq 0 . \tag{3.1}
\end{equation*}
$$

By the implicit function theorem, there exists a $C^{r-1}$ function $t=\eta(y)$ defined in a small neighborhood $V$ of 0 in the $y$-axis and satisfying $\eta(0)=t_{0}$ and $\partial \varphi(y, \eta(y)) / \partial t=0$. Then the curve $\widetilde{\gamma}$ in $V_{\widetilde{q}}$ parametrized by $(\varphi(y, \eta(y)), y, \eta(y))$ divides $\widetilde{H}$ into two components and $\gamma=\operatorname{pr}(\widetilde{\gamma})$ is a $C^{r-1}$ curve embedded in $D_{a}(p)$. Let $\widetilde{H}^{+}$(resp. $\left.\widetilde{H}^{-}\right)$be the closure of the upper (resp. lower) component of $\widetilde{H} \backslash \widetilde{\gamma}$. For a sufficiently large $n_{0} \in \mathbb{N}$, the component $\widetilde{H}_{0}$ of $f^{n_{0}}(\widetilde{H}) \cap U_{a}(p)$ containing $q_{0}=f^{n_{0}}(\widetilde{q})$ is an unstable bent disk in $U_{a}(p)$ such that $\partial \widetilde{H}_{0}$ is a simple closed $C^{r}$ curve in $\partial_{\text {side }} U_{a}(p)$, where

$$
\partial_{\text {side }} U_{a}(p)=\{(x, t) \in \mathbb{C} \times \mathbb{R} ;|z|=a, 0 \leq t<a\} \subset \partial U_{a}(p) .
$$

See Figure 3.2. We set $\widetilde{\gamma}_{0}=f^{n_{0}}(\widetilde{\gamma}) \cap \widetilde{H}_{0}, \widetilde{H}_{0}^{+}=f^{n_{0}}\left(\widetilde{H}^{+}\right) \cap \widetilde{H}_{0}, \widetilde{H}_{0}^{-}=f^{n_{0}}\left(\widetilde{H}^{-}\right) \cap \widetilde{H}_{0}$, $H_{0}=\operatorname{pr}\left(\widetilde{H}_{0}^{+}\right)=\operatorname{pr}\left(\widetilde{H}_{0}^{-}\right)$and $\gamma_{0}=\operatorname{pr}\left(\widetilde{\gamma}_{0}\right)$. Then $\widetilde{\gamma}_{0}$ is called the front curve of $\widetilde{H}_{0}$ and $\gamma_{0}$ is the folding curve of $H_{0}$.


Figure 3.2: The front curve $\widetilde{\gamma}_{0}$ divides $\widetilde{H}_{0}$ into the two sheets $\widetilde{H}_{0}^{+}$and $\widetilde{H}_{0}^{-}$. The folding curve $\gamma_{0}$ of $H_{0}$ is the orthogonal image of $\widetilde{\gamma}_{0}$.

We note that Nishizawa [Ni] has studied unstable bent disks similar to $\widetilde{H}_{0}$ as above in a different situation. In fact, he considered a 3-dimensional diffeomorphism $g$ which
has a saddle fixed point $s$ such that all the eigenvalues of $D g(s)$ are real and has a homoclinic quadratic tangency associated with $s$. Here we consider the component $\widetilde{H}_{0 ; u}^{-}$of $f^{u}\left(\widetilde{H}_{0}^{-}\right) \cap U_{a}(p)$ containing $f^{u}\left(q_{0}\right)$ for $u \in \mathbb{N}$. Since the homoclinic tangency $q$ is positively associated with $p$, one can show that there exists $\widetilde{H}_{0 ; u}^{-}$which meets $W^{s}(p)$ transversely at a point $\widehat{z}$ near $q$ by using an argument similar to that in [Ni, Lemma 4.4]. See Figure 3.3. To show the claim, the assumption of $\theta_{0} \neq 0 \bmod \pi$ in Theorem 3.1 is crucial. In fact,


Figure 3.3: The half disk $\widetilde{H}_{0 ; u}^{-}$meets $W^{s}(p)$ transversely at two points near $q$, one of which is $\widehat{z}$.
the condition implies that the following property:
(P) There exists an arbitrarily large $u$ such that the interior of $H_{0 ; u}=\operatorname{pr}\left(\widetilde{H}_{0 ; u}^{-}\right)$in $D_{a}(p)$ contains $q$.

Remark 3.4. (1) We here suppose $\theta=0 \bmod \pi$. Even in this case, if $f$ has the property (P), then the component of $W^{s}(p)$ containing $q$ and $W^{u}(p)$ have a homoclinic transverse intersection point. Then Theorems 3.1 and 3.3 will be proved quite similarly. Since $\theta=0$ $\bmod \pi$, all $f^{u}\left(\gamma_{0}\right)$ are tangent to a unique straight segment $\gamma_{\infty}$ in $D_{a}(p)$ at $p$. Thus the property ( P ) is satisfied if $\gamma_{\infty}$ does not pass through $q$.
(2) Even in the case of $-1<\lambda<0$, one can show that $f$ has the property ( P ) similarly by using $f^{2}$ instead of $f$ if $2 \theta \neq 0 \bmod \pi$. Moreover, since either $q$ or $f(q)$ is a homoclinic tangency positively associated with $p$, Theorems 3.1 and 3.3 hold without the assumption that $q$ is positively associated with $p$.

### 3.2.2 Construction of convergent sequence of unstable bent disks

Take $v \in \mathbb{N}$ such that $\widehat{z}_{0}=f^{v}(\widehat{z})$ is a point $(0, \widehat{t})$ contained in $U_{a}(p)$, where $\widehat{z}$ is the transverse intersection point of $\widetilde{H}_{0 ; u}^{-}$and $W^{s}(p)$ given in the previous subsection. Let $D$ be a small disk in $W^{u}(p) \cap U_{a}(p)$ whose interior contains $\widehat{z}_{0}$. The absolute slope $\sigma(\boldsymbol{v})$ of a
vector $\boldsymbol{v}=\left(v_{1}, v_{2}, v_{3}\right)$ in $U_{a}(p)$ with $\left(v_{1}, v_{2}\right) \neq(0,0)$ is given as

$$
\sigma(\boldsymbol{v})=\frac{\left|v_{3}\right|}{\sqrt{v_{1}^{2}+v_{2}^{2}}}
$$

The maximum absolute slope $\sigma(D)$ of $D$ is defined by

$$
\sigma(D)=\max \left\{\sigma(\boldsymbol{v}) ; \text { unit vectors } \boldsymbol{v} \text { in } U_{a}(p) \text { tangent to } D\right\}
$$

Fix $m_{0} \in \mathbb{N}$ such that, for any $m \in \mathbb{N} \cup\{0\}$, the component $D_{m}$ of $f^{m_{0}+m}(D) \cap U(p)$ containing $f^{m_{0}+m}\left(\widehat{z}_{0}\right)$ is a properly embedded disk in $U_{a}(p)$ with $\partial D_{m} \subset \partial_{\text {side }} U_{a}(p)$. Note that $D_{m}$ intersects $W_{\text {loc }}^{s}(p)$ transversely at $\left(0, \lambda^{m} t_{0}\right)$, where $t_{0}=\lambda^{m_{0}} \widehat{t}$. See Figure 3.4.


Figure 3.4: Trip from $\widetilde{H}_{0}^{-}$to $\widetilde{H}_{m}: f^{u+v}\left(\widetilde{H}_{0}^{-}\right) \supset D, f^{m_{0}}(D) \supset D_{0}, f^{m}\left(D_{0}\right) \supset D_{m}$ and $f^{N+n_{0}}\left(D_{m}\right) \supset \widetilde{H}_{m}$, where $N, n_{0}$ are the positive integers with $f^{N}(q)=\widetilde{q}$ and $f^{n_{0}}(\widetilde{q})=q_{0}$. The dotted line passing through $q$ represents a straight segment tangent to $\widetilde{\rho}$ at $q$.

The maximum absolute slope of $D_{m}$ satisfies

$$
\begin{equation*}
\sigma\left(D_{m}\right) \leq \sigma_{0} \lambda^{m} r^{-m} \tag{3.2}
\end{equation*}
$$

where $\sigma_{0}=\sigma(D) \lambda^{m_{0}} r^{-m_{0}}$. Consider a short straight segment $\rho$ in $U_{a}(p)$ meeting $\widetilde{H}_{0}$ orthogonally at $q_{0}$. Then $\widetilde{\rho}=f^{-\left(N+n_{0}\right)}(\rho)$ is a $C^{r}$ curve meeting $D_{a}(p)$ transversely at $q$, where $N, n_{0}$ are the positive integers given as above. One can choose $m_{0} \in \mathbb{N}$ so that, for any $m \in \mathbb{N} \cup\{0\}, \widetilde{\rho}$ meets $D_{m}$ transversely at a single point $\boldsymbol{w}_{m}=\left(z_{m}, s_{m}\right)$. Then (3.2) implies that $\left|t_{0} \lambda^{m}-s_{m}\right| \leq \widetilde{a} \sigma_{0} \lambda^{m} r^{-m}$, where $\widetilde{a}=\sup _{m \geq 0}\left\{\left|z_{m}\right|\right\}<\infty$. It follows that $s_{m}=t_{0} \lambda^{m}+O\left(\lambda^{m} r^{-m}\right)$. Since $\widetilde{\rho}$ has a tangency of order at least two with a straight segment at $q$,

$$
\begin{equation*}
\operatorname{dist}\left(\boldsymbol{w}_{m}, q\right)=\widetilde{t}_{0} \lambda^{m}+O\left(\lambda^{m} r^{-m}\right)+O\left(\lambda^{2 m}\right)=\widetilde{t}_{0} \lambda^{m}+o\left(\lambda^{m}\right) \tag{3.3}
\end{equation*}
$$

for some constant $\widetilde{t}_{0}>0$. By the inclination lemma, $D_{m}$ uniformly $C^{r}$ converges to $D_{a}(p)$. A short curve in $W^{s}(p)$ containing $q$ as an interior point meets $D_{m}$ transversely in
two points for all sufficiently large $m$. Let $\widetilde{H}_{m}$ be the component of $f^{N+n_{0}}\left(D_{m}\right) \cap U_{a}(p)$ containing $f^{N+n_{0}}\left(\boldsymbol{w}_{m}\right)$. Then $\widetilde{H}_{m} C^{r}$ converges to $\widetilde{H}_{0}$ as $m \rightarrow \infty$. By (3.1), there exist $C^{r}$ functions $\varphi_{m}(y, t) C^{r}$ converging to $\varphi$ and representing $\widetilde{H}_{m}$ as the graph of $x=\varphi_{m}(y, t)$. Then the front curve $\widetilde{\gamma}_{m}$ in $\widetilde{H}_{m}$ is defined as the front curve $\widetilde{\gamma}_{0}$ in $\widetilde{H}_{0}$. Since $\partial \varphi_{m}(y, t) / \partial t C^{r-1}$ converges to $\partial \varphi(y, t) / \partial t, \widetilde{\gamma}_{m}$ also $C^{r-1}$ converges to $\widetilde{\gamma}_{0}$. Note that $\widetilde{\gamma}_{m}$ divides $\widetilde{H}_{m}$ into the upper surface $\widetilde{H}_{m}^{+}$and the lower surface $\widetilde{H}_{m}^{-}$with $\widetilde{\gamma}_{m}=\widetilde{H}_{m}^{+} \cap \widetilde{H}_{m}^{-}$and $H_{m}=\operatorname{pr}\left(\widetilde{H}_{m}\right)=\operatorname{pr}\left(\widetilde{H}_{m}^{+}\right)=\operatorname{pr}\left(\widetilde{H}_{m}^{-}\right)$. The image $\gamma_{m}=\operatorname{pr}\left(\widetilde{\gamma}_{m}\right)$ is called the folding curve of $H_{m}$.

### 3.3 Limit straight segments

A curve $\gamma$ in $D_{a}(p)$ is called a straight segment if $\gamma$ is a segment with respect to the Euclidean metric on $D_{a}(p)$. In this section, we will construct a proper straight segment $\gamma_{0}^{\natural}$ in $D_{a}(p)$ with $p \notin \gamma_{0}^{\natural}$ which is mapped to a straight segment in $U_{a^{\prime}}\left(p^{\prime}\right)$ by $h$.

### 3.3.1 Sequences of folding curves converging to straight segments

Let $\alpha$ be an oriented $C^{r-1}$ curve in $D_{a}(p)$ of bounded length. Since $r-1 \geq 2$, there exists the maximum absolute curvature $\kappa(\alpha)$ of $\alpha$. If $\alpha$ passes near the center 0 of $D_{a}(p)$ and satisfies $\kappa(\alpha)<1 / a$, then $\alpha$ has a unique point $z(\alpha)$ with $\operatorname{dist}(0, z(\alpha))=\operatorname{dist}(0, \alpha)$. In fact, if $\alpha$ had two points $z_{i}(i=1,2)$ with $\operatorname{dist}\left(0, z_{i}\right)=\operatorname{dist}(0, \alpha)$, then for a point $z_{3}$ in $\alpha$ with the maximum $\operatorname{dist}\left(0, z_{3}\right)$ between $z_{1}$ and $z_{2}$, the curvature of $\alpha$ at $z_{3}$ is not less than $1 / \operatorname{dist}\left(0, z_{3}\right) \geq 1 / a$, a contradiction. We denote by $\vartheta(\alpha) \bmod 2 \pi$ the angle between $\widehat{\alpha}$ and the positive direction of the $x$-axis at 0 , where $\widehat{\alpha}$ is the oriented curve in $D_{a}(p)$ obtained from $\alpha$ by the parallel translation taking $z(\alpha)$ to 0 .

By (3.3), there exists a constant $\widetilde{d}_{0}>0$ such that

$$
\begin{equation*}
\operatorname{dist}\left(\widetilde{\gamma}_{m}, \text { the } t \text {-axis }\right)=\widetilde{d}_{0}\left(\widetilde{t}_{0} \lambda^{m}+o\left(\lambda^{m}\right)\right)+o\left(\lambda^{m}\right)=\widetilde{d}_{0} \widetilde{t}_{0} \lambda^{m}+o\left(\lambda^{m}\right) \tag{3.4}
\end{equation*}
$$

Since $\gamma_{m} C^{r-1}$ converges to $\gamma_{0}, \kappa\left(\gamma_{m}\right)$ also converges to $\kappa\left(\gamma_{0}\right)$ as $m \rightarrow \infty$. This shows that

$$
\begin{equation*}
\sup _{m}\left\{\kappa\left(\gamma_{m}\right)\right\}=\kappa_{0}<\infty . \tag{3.5}
\end{equation*}
$$

It follows that, for all sufficiently large $m$, there exists a unique point $c_{m}$ of $\gamma_{m}$ with

$$
\operatorname{dist}\left(c_{m}, 0\right)=\operatorname{dist}\left(\gamma_{m}, 0\right)=\operatorname{dist}\left(\widetilde{c}_{m}, \text { the } t \text {-axis }\right)=\operatorname{dist}\left(\widetilde{\gamma}_{m}, \text { the } t \text {-axis }\right),
$$

where $\widetilde{c}_{m}$ is the point of $\widetilde{\gamma}_{m}$ with $\operatorname{pr}\left(\widetilde{c}_{m}\right)=c_{m}$.
Fix $w$ with $0<w<a / 2$ arbitrarily. For any $n \in \mathbb{N}$, let $m(n)$ be the minimum positive integer such that $f^{n}\left(c_{m}\right)$ is contained in $D_{w}(p)$ for any $m \geqq m(n)$. Then $\lim _{n \rightarrow \infty} m(n)=$ $\infty$ holds. For any $m \geq m(n)$, the component $\widetilde{H}_{m, n}$ of $f^{n}\left(\widetilde{\widetilde{H}}_{m}\right) \cap U_{a}(p)$ containing $\widetilde{c}_{m, n}=$ $f^{n}\left(\widetilde{c}_{m}\right)$ is a proper disk in $U_{a}(p)$ with $\partial \widetilde{H}_{m, n} \subset \partial_{\text {side }} U_{a}(p)$. Then $\widetilde{\gamma}_{m, n}=f^{n}\left(\widetilde{\gamma}_{m}\right) \cap \widetilde{H}_{m, n}$ is the front curve of $\widetilde{H}_{m, n}$ and $\gamma_{m, n}=\operatorname{pr}\left(\widetilde{\gamma}_{m, n}\right)$ is the folding curve of $H_{m, n}=\operatorname{pr}\left(\widetilde{H}_{m, n}\right)$. Then $c_{m, n}=\operatorname{pr}\left(\widetilde{c}_{m, n}\right)$ is a unique point of $\gamma_{m, n}$ closest to 0 . Here we orient $\widetilde{\gamma}_{m}=\widetilde{\gamma}_{m, 0}$ so
that $\widetilde{\gamma}_{m, 0} C^{r-1}$ converges as oriented curves to $\widetilde{\gamma}_{0}$ as $m \rightarrow \infty$. Suppose that $\gamma_{m, n}$ has the orientation induced from that on $\widetilde{\gamma}_{m, 0}$ via pr $\circ f^{n}$. In particular, it follows that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \vartheta\left(\gamma_{m, 0}\right)=\vartheta\left(\gamma_{0}\right) . \tag{3.6}
\end{equation*}
$$

We set $d_{m, n}=\operatorname{dist}\left(c_{m, n}, 0\right)$. By (3.4),

$$
\begin{equation*}
d_{m, n}=r^{n}\left(\widetilde{d}_{0} \widetilde{t}_{0} \lambda^{m}+o\left(\lambda^{m}\right)\right) . \tag{3.7}
\end{equation*}
$$

There exist subsequences $\left\{m_{j}\right\},\left\{n_{j}\right\}$ of $\mathbb{N}$ and $w \lambda / 2 \leq w_{0} \leq w$ such that

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \widetilde{d}_{0} \widetilde{t}_{0} \lambda^{m_{j}} r^{n_{j}}=w_{0} \tag{3.8}
\end{equation*}
$$

If necessary taking subsequences of $\left\{m_{j}\right\}$ and $\left\{n_{j}\right\}$ simultaneously, we may also assume that $\vartheta\left(\gamma_{m_{j}, n_{j}}\right)$ has a limit $\theta^{\natural}$. Since $f(z)=r e^{\sqrt{-1} \theta} z$ on $D_{a}(p)$, by (3.5) we have

$$
\kappa\left(\gamma_{m_{j}, n_{j}}\right) \leq r^{-n_{j}} \kappa\left(\gamma_{m_{j}, 0}\right) \leq r^{-n_{j}} \kappa_{0} \rightarrow 0 \quad \text { as } \quad j \rightarrow \infty .
$$

Thus the following lemma is obtained immediately.
Lemma 3.5. The sequence $\gamma_{m_{j}, n_{j}}$ uniformly converges as oriented curves to an oriented straight segment $\gamma_{0}^{\natural}$ in $D_{a}(p)$ with $\vartheta\left(\gamma_{0}^{\natural}\right)=\theta^{\natural}$ and $\operatorname{dist}\left(\gamma_{0}^{\natural}, 0\right)=w_{0}$.

We say that $\gamma_{0}^{\natural}$ is the limit straight segment of $\gamma_{m_{j}, n_{j}}$.

### 3.3.2 Limit straight segments preserved by the conjugacy

Let $U_{a^{\prime}}\left(p^{\prime}\right), U_{b^{\prime}}\left(p^{\prime}\right)$ be the circular columns defined as $U_{a}(p)$ for some $0<a^{\prime}<b^{\prime}$ which are contained in a coordinate neighborhood around $p^{\prime}$ with respect to which $f^{\prime}$ is linear. One can retake $a>0$ and choose such $a^{\prime}, b^{\prime}$ so that $U_{a^{\prime}}\left(p^{\prime}\right) \subset \underset{\tilde{H}}{\underset{H}{r}} \underset{a}{ }\left(U_{a}(p)\right) \subset U_{b^{\prime}}\left(p^{\prime}\right)$.

Let $\widetilde{H}_{m, n}^{\prime}$ be the component of $h\left(\widetilde{H}_{m, n}\right) \cap U_{a^{\prime}}\left(p^{\prime}\right)$ defined as $\widetilde{H}_{m, n}$ and $\operatorname{pr}\left(\widetilde{H}_{m, n}^{\prime}\right)=H_{m, n}^{\prime}$. One can define the front and folding curves $\widetilde{\gamma}_{m, n}^{\prime}, \gamma_{m, n}^{\prime}$ in $\widetilde{H}_{m, n}^{\prime}$ and $H_{m, n}^{\prime}$ as $\widetilde{\gamma}_{m, n}, \gamma_{m, n}$ in $\widetilde{H}_{m, n}$ and $H_{m, n}$ respectively. See Figure 3.5.

Since $h$ is only supposed to be a homeomorphism, $h\left(\widetilde{\gamma}_{m, n}\right) \cap U_{a^{\prime}}\left(p^{\prime}\right)$ would not be equal to $\widetilde{\gamma}_{m, n}^{\prime}$. We will show that this equality holds in the limit. For the sequences $\left\{m_{j}\right\},\left\{n_{j}\right\}$ given in Section 3.3, we set $\widetilde{H}_{m_{j}, n_{j}}=\widetilde{H}_{(j)}, H_{m_{j}, n_{j}}=H_{(j)}, \widetilde{H}_{m_{j}, n_{j}}^{\prime}=\widetilde{H}_{(j)}^{\prime}$ and $H_{m_{j}, n_{j}}^{\prime}=$ $H_{(j)}^{\prime}$ for simplicity. Similarly, suppose that $\widehat{H}_{(j)}^{\prime}$ is the component of $W^{u}\left(p^{\prime}\right) \cap U_{b^{\prime}}\left(p^{\prime}\right)$ containing $\widetilde{H}_{(j)}^{\prime}$ and $\widehat{\gamma}_{m_{j}, n_{1}}^{\prime}$ is the front curve of $\widehat{H}_{(j)}^{\prime}$. The distance between $\boldsymbol{x}, \boldsymbol{y}$ in $U_{a}(p)$ is denoted by $d(\boldsymbol{x}, \boldsymbol{y})$ and that between $\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}$ in $U_{a^{\prime}}\left(p^{\prime}\right)$ by $d^{\prime}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)$.

The path metric on $\widetilde{H}_{(j)}$ is denoted by $d_{\widetilde{H}_{(j)}}$. That is, for any $\boldsymbol{x}, \boldsymbol{y} \in \widetilde{H}_{(j)}, d_{\widetilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})$ is the length of a shortest piecewise smooth curve in $\widetilde{H}_{(j)}$ connecting $\boldsymbol{x}$ with $\boldsymbol{y}$. The path metrics $d_{\tilde{H}_{(j)}^{\prime}}$ on $\widetilde{H}_{(j)}^{\prime}$ and $d_{\widehat{H}_{(j)}^{\prime}}$ on $\widehat{H}_{(j)}^{\prime}$ are defined similarly.
Lemma 3.6. (i) For any $\varepsilon>0$, there exists a constant $\eta(\varepsilon)>0$ independent of $j \in \mathbb{N}$ and satisfying the following conditions.


Figure 3.5: The image $h\left(\widetilde{H}_{(j)}\right)$ is contained in $\widehat{H}_{(j)}^{\prime}$, but $h\left(\widetilde{H}_{(j)}^{ \pm}\right)$is not necessarily contained in $\widehat{H}_{(j)}^{\prime \pm}$.

- $\lim _{\varepsilon \rightarrow 0} \eta(\varepsilon)=0$.
- Let $\boldsymbol{x}, \boldsymbol{y}$ be any points of $\widetilde{H}_{(j)}$ both of which are contained in one of $\widetilde{H}_{(j)}^{+}$and $\widetilde{H}_{(j)}^{-}$. If $d(\boldsymbol{x}, \boldsymbol{y})<\eta(\varepsilon)$, then $d_{\tilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})<\varepsilon$.
(ii) For any $\varepsilon>0$, there exists a constant $\delta(\varepsilon)>0$ independent of $j \in \mathbb{N}$ and satisfying the following conditions.
- $\lim _{\varepsilon \rightarrow 0} \delta(\varepsilon)=0$.
- Let $\boldsymbol{x}, \boldsymbol{y}$ be any points of $\widetilde{H}_{(j)}$ both of which are contained in one of $\widetilde{H}_{(j)}^{+}$and $\widetilde{H}_{(j)}^{-}$. If $d_{\tilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})<\delta(\varepsilon)$ and $\boldsymbol{x}^{\prime}=h(\boldsymbol{x})$ and $\boldsymbol{y}^{\prime}=h(\boldsymbol{y})$ are contained in $\widetilde{H}_{(j)}^{\prime}$, then $d_{\tilde{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)<\varepsilon$.

One can take these constants $\eta(\varepsilon), \delta(\varepsilon)$ so that they work also for $d_{\widetilde{H}_{(j)}^{\prime}}$ and $d_{\widehat{H}_{(j)}^{\prime}}$.
Proof. (i) The assertion is proved immediately from the fact that $\widetilde{H}_{(j)}^{ \pm}$uniformly converges to a disk $H^{\natural}$ in $D_{a}(p)$ such that $d(\boldsymbol{x}, \boldsymbol{y})=d_{H^{\natural}}(\boldsymbol{x}, \boldsymbol{y})$ for any $\boldsymbol{x}, \boldsymbol{y} \in H^{\natural}$.
(ii) Suppose that $\boldsymbol{x}, \boldsymbol{y} \in \widetilde{H}_{(j)}^{+}$. First we consider the case that both $\boldsymbol{x}^{\prime}$ and $\boldsymbol{y}^{\prime}$ are contained in one of $\widetilde{H}_{(j)}^{\prime+}$ and $\widetilde{H}_{(j)}^{\prime-}$, say $\widetilde{H}_{(j)}^{\prime+}$. If $d_{\widetilde{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right) \geq \varepsilon$, then it follows from the assertion (i) that $d^{\prime}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right) \geq \eta(\varepsilon)$. Since $h$ is uniformly continuous on $U_{a}(p)$, there exists a constant $\delta_{1}(\varepsilon)>0$ with $\lim _{\varepsilon \rightarrow 0} \delta_{1}(\varepsilon)=0$ and $d(\boldsymbol{x}, \boldsymbol{y}) \geq \delta_{1}(\varepsilon)$. Hence, in particular, $d_{\tilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y}) \geq$ $\delta_{1}(\varepsilon)$. Thus $d_{\widetilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})<\delta_{1}(\varepsilon)$ implies $d_{\widetilde{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)<\varepsilon$.

Next we suppose that $\boldsymbol{x}^{\prime} \in \widetilde{H}_{(j)}^{\prime+}$ and $\boldsymbol{y}^{\prime} \in \widetilde{H}_{(j)}^{\prime-}$. Consider a shortest curve $\alpha$ in $\widetilde{H}_{(j)}$ connecting $\boldsymbol{x}$ and $\boldsymbol{y}$. Since $\alpha^{\prime}=h(\alpha)$ is contained in $\widehat{H}_{(j)}^{\prime}, \alpha^{\prime}$ intersects $\widehat{\gamma}_{m_{j}, n_{j}}^{\prime}$ non-trivially.


Figure 3.6: The case of $\boldsymbol{x}, \boldsymbol{y} \in \widetilde{H}_{(j)}^{+}, \boldsymbol{x}^{\prime} \in \widetilde{H}_{(j)}^{\prime+}$ and $\boldsymbol{y}^{\prime} \in \widetilde{H}_{(j)}^{\prime-}$.

Let $\boldsymbol{z}$ be one of the intersection points of $\alpha$ with $h^{-1}\left(\widehat{\gamma}_{m_{j}, n_{j}}^{\prime}\right)$. See Figure 3.6. Suppose that $d_{\tilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})<\delta_{1}(\varepsilon / 2)$. Since $d_{\tilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})=d_{\widetilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{z})+d_{\tilde{H}_{(j)}}(\boldsymbol{z}, \boldsymbol{y})$,

$$
d_{\tilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{z})<\delta_{1}(\varepsilon / 2) \quad \text { and } \quad d_{\tilde{H}_{(j)}}(\boldsymbol{z}, \boldsymbol{y})<\delta_{1}(\varepsilon / 2) .
$$

Since $\boldsymbol{x}^{\prime}, \boldsymbol{z}^{\prime} \in \widehat{H}_{(j)}^{\prime+}$ and $\boldsymbol{z}^{\prime}, \boldsymbol{y}^{\prime} \in \widehat{H}_{(j)}^{\prime-}$, by the result in the previous case we have $d_{\widehat{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{z}^{\prime}\right)<$ $\varepsilon / 2$ and $d_{\widehat{H}_{(j)}^{\prime}}\left(\boldsymbol{z}^{\prime}, \boldsymbol{y}^{\prime}\right)<\varepsilon / 2$, and hence

$$
d_{\widetilde{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)=d_{\widehat{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)<\varepsilon .
$$

Thus $\delta(\varepsilon):=\delta_{1}(\varepsilon / 2)$ satisfies the conditions of (ii).
The following result is a key of this section.
Lemma 3.7. For any $\varepsilon>0$, there exists $j_{0} \in \mathbb{N}$ such that, for any $j \geq j_{0}$,

$$
h\left(\widetilde{\gamma}_{m_{j}, n_{j}}\right) \cap \widetilde{H}_{(j)}^{\prime} \subset \mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right),
$$

where $\mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)$ is the $\varepsilon$-neighborhood of $\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}$ in $\widetilde{H}_{(j)}^{\prime}$.
Figure 3.7 illustrates the situation of Lemma 3.7.
Proof. For $\sigma= \pm$, we will show that $h^{-1}\left(\widetilde{H}_{(j)}^{\prime \sigma} \backslash \mathcal{N}_{\varepsilon}\left(\gamma_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)\right) \subset \widetilde{H}_{(j)}^{\sigma}$ for all sufficiently large $j$. Since $\left.h^{-1}\right|_{U_{a^{\prime}}\left(p^{\prime}\right)}$ is uniformly continuous, there exists $\nu(\varepsilon)>0$ such that, for any $\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime} \in U_{a^{\prime}}\left(p^{\prime}\right)$ with $d^{\prime}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)<\nu(\varepsilon)$, the inequality $d(\boldsymbol{x}, \boldsymbol{y})<\eta(\delta(\varepsilon))$ holds, where $\boldsymbol{x}=h^{-1}\left(\boldsymbol{x}^{\prime}\right), \boldsymbol{y}=h^{-1}\left(\boldsymbol{y}^{\prime}\right)$. Since both $\widetilde{H}_{(j)}^{\prime+}$ and $\widetilde{H}_{(j)}^{\prime-}$ uniformly converge to the same half disk $H^{\prime \dagger}$ in $D_{a^{\prime}}\left(p^{\prime}\right)$, there exists $j_{0} \in \mathbb{N}$ such that, for any $j \geq j_{0}$ and any $\boldsymbol{x}^{\prime} \in \widetilde{H}_{(j)}^{\prime \sigma} \backslash$ $\mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{(j)}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right), d^{\prime}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)$ is less than $\nu(\varepsilon)$, where $\boldsymbol{y}^{\prime}$ is the element of $\widetilde{H}_{(j)}^{\prime-\sigma}$ with $\operatorname{pr}\left(\boldsymbol{x}^{\prime}\right)=$ $\operatorname{pr}\left(\boldsymbol{y}^{\prime}\right)$. Then we have $d(\boldsymbol{x}, \boldsymbol{y})<\eta(\delta(\varepsilon))$. If both $\boldsymbol{x}$ and $\boldsymbol{y}$ were contained in one of $\widetilde{H}_{(j)}^{\sigma}$


Figure 3.7: The shaded region represents $\mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)$.
and $\widetilde{H}_{(j)}^{-\sigma}$, then by Lemma 3.6 (i) $d_{\widetilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})<\delta(\varepsilon)$. Then, by Lemma 3.6 (ii), $d_{\widetilde{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)$ would be less than $\varepsilon$. This contradicts that $\boldsymbol{x}^{\prime} \in \widetilde{H}_{(j)}^{\prime \sigma} \backslash \mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)$ and $\boldsymbol{y}^{\prime} \in \widetilde{H}_{(j)}^{\prime-\sigma}$. See Figure 3.8. Thus, if $\boldsymbol{y}$ is contained in $\widetilde{H}_{(j)}^{\sigma}$, then $\boldsymbol{x}$ is not in $\widetilde{H}_{(j)}^{\sigma}$. In particular, $\boldsymbol{x}$ is


Figure 3.8: The situation which does not actually occur. $d_{1}:=\operatorname{dist}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)<\nu(\varepsilon), d_{2}:=$ $\operatorname{dist}_{\widetilde{H}_{(j)}}(\boldsymbol{x}, \boldsymbol{y})<\delta(\varepsilon)$ and $d_{3}:=\operatorname{dist}_{\widetilde{H}_{(j)}^{\prime}}\left(\boldsymbol{x}^{\prime}, \boldsymbol{y}^{\prime}\right)<\varepsilon$.
not contained in $\widetilde{\gamma}_{m_{j}, n_{j}}=\widetilde{H}_{(j)}^{+} \cap \widetilde{H}_{(j)}^{-}$, and so $\widetilde{\gamma}_{m_{j}, n_{j}} \cap h^{-1}\left(\widetilde{H}_{(j)}^{\prime \sigma} \backslash \mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m, n}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)\right)=\emptyset$. Since $h^{-1}\left(\widetilde{H}_{(j)}^{\prime \sigma} \backslash \mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m, n}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)\right)$ is connected, it follows that $h^{-1}\left(\widetilde{H}_{(j)}^{\sigma} \backslash \mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m, n}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)\right) \subset \widetilde{H}_{(j)}^{\sigma}$ for $\sigma= \pm$, and hence $h^{-1}\left(\mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)\right) \supset \widetilde{\gamma}_{m_{j}, n_{j}} \cap h^{-1}\left(\widetilde{H}_{(j)}^{\prime}\right)$. This completes the proof.

From the proof of Lemma 3.7, we know that there exists a simple curve in $h\left(\widetilde{\gamma}_{m_{j}, n_{j}}\right) \cap$ $\widetilde{H}_{(j)}^{\prime}$ connecting the two components of $\partial \widetilde{H}_{(j)}^{\prime} \cap \partial \mathcal{N}_{\varepsilon}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}, \widetilde{H}_{(j)}^{\prime}\right)$. The following corollary says that the images of certain straight segments in $D_{a}(p)$ by the homeomorphism $h$ are
naturally straight segments in $D_{a^{\prime}}\left(p^{\prime}\right)$.
Corollary 3.8. For the limit straight segment $\gamma_{0}^{\natural}$ of $\gamma_{m_{j}, n_{j}}, h\left(\gamma_{0}^{\natural}\right) \cap D_{a^{\prime}}\left(p^{\prime}\right)$ is the limit straight segment of $\gamma_{m_{j}, n_{j}}^{\prime}$, i.e. $h\left(\gamma_{0}^{\natural}\right) \cap D_{a^{\prime}}\left(p^{\prime}\right)=\gamma_{0}^{\prime \dagger}$.

Proof. Since $\gamma_{0}^{\natural}$ is the limit straight segment of $\widetilde{\gamma}_{m_{j}, n_{j}}$ and $h$ is uniformity continuous, $h\left(\gamma_{0}^{\natural}\right) \cap D_{a^{\prime}}\left(p^{\prime}\right)$ is the limit of $h\left(\widetilde{\gamma}_{m_{j}, n_{j}}\right) \cap \widetilde{H}_{(j)}^{\prime}$. It follows from Lemma 3.7 that $h\left(\gamma_{0}^{\natural}\right) \cap D_{a^{\prime}}\left(p^{\prime}\right)$ is also the limit of $\operatorname{pr}\left(\widetilde{\gamma}_{m_{j}, n_{j}}^{\prime}\right)=\gamma_{m_{j}, n_{j}}^{\prime}$, that is, $h\left(\gamma_{0}^{\natural}\right) \cap D_{a^{\prime}}\left(p^{\prime}\right)$ is equal to the limit straight segment of $\gamma_{m_{j}, n_{j}}^{\prime}$.

For any straight segment $l$ in $D_{a}(p)$ such that $h(l)$ is also a straight segment in $D_{b^{\prime}}\left(p^{\prime}\right)$, we denote $h(l) \cap D_{a^{\prime}}\left(p^{\prime}\right)$ simply by $h(l)$. In particular, Corollary 3.8 implies that $h\left(\gamma_{0}^{\natural}\right)=$ $\gamma_{0}^{\prime 4}$.

### 3.4 Proof of Theorem 3.1

Suppose that $\mathrm{St}_{a}(p)$ is the set of oriented proper straight segments in $D_{a}(p)$ passing through 0 , that is, each element of $\operatorname{St}_{a}(p)$ is an oriented diameter of the disk $D_{a}(p)$. For any $l \in \operatorname{St}_{a}(p)$ and $n \in \mathbb{N}$, the component of $f^{n}(l) \cap U_{a}(p)$ containing 0 is also an element of $\operatorname{St}_{a}(p)$. We denote the element simply by $f^{n}(l)$.

Since $\left.f^{n}\right|_{D_{a}(p)}$ preserves angles on $D_{a}(p)$, by (3.6), for any $k, n \in \mathbb{N}$,

$$
\vartheta\left(\gamma_{m, n}\right)-\vartheta\left(\gamma_{m+k, n}\right)=\vartheta\left(\gamma_{m, 0}\right)-\vartheta\left(\gamma_{m+k, 0}\right) \rightarrow \vartheta\left(\gamma_{0}\right)-\vartheta\left(\gamma_{0}\right)=0
$$

as $m \rightarrow \infty$. Moreover it follows from (3.7) that $\lim _{j \rightarrow \infty} d_{m_{j}+k, n_{j}}=w_{0} \lambda^{k}$. By these facts together with Lemma 3.5, one can show that $\gamma_{m_{j}+k, n_{j}}$ uniformly converges as $m \rightarrow \infty$ to a straight segment $\gamma_{k}^{\natural}$ in $U_{a}(p)$ with

$$
\begin{equation*}
\vartheta\left(\gamma_{k}^{\natural}\right)=\theta^{\natural} \quad \text { and } \quad d\left(0, \gamma_{k}^{\natural}\right)=w_{0} \lambda^{k} . \tag{3.9}
\end{equation*}
$$

Thus we have obtained the parallel family $\left\{\gamma_{k}^{\natural}\right\}$ of oriented straight segments in $D_{a}(p)$. See Figure 3.9. By Corollary 3.8, $\left\{\gamma_{k}^{\prime \dagger}\right\}$ with $\gamma_{k}^{\prime \natural}=h\left(\gamma_{k}^{\natural}\right)$ is also a parallel family of oriented straight segments in $D_{a^{\prime}}\left(p^{\prime}\right)$. Since $\gamma_{k}^{\prime \natural}$ is the limit of $\gamma_{m_{j}+k, n_{j}}^{\prime}$ as $j \rightarrow \infty$, we have the equations

$$
\begin{equation*}
\vartheta\left(\gamma_{k}^{\prime \boldsymbol{\natural}}\right)=\theta^{\prime \text { Ł }} \quad \text { and } \quad d\left(0, \gamma_{k}^{\prime \text { Ø }}\right)=w_{0}^{\prime} \lambda^{\prime k} . \tag{3.10}
\end{equation*}
$$

corresponding to (3.9) for some $\theta^{\prime \natural}$ and $w_{0}^{\prime}>0$. Let $\gamma_{\infty}^{\natural} \in \operatorname{St}_{a}(p)$ (resp. $\gamma_{\infty}^{\prime \natural} \in \operatorname{St}_{a^{\prime}}\left(p^{\prime}\right)$ ) be the limit of $\gamma_{k}^{\natural}$ (resp. $\gamma_{k}^{\prime \natural}$ ).

Proof of Theorem 3.1. By Lemma 3.5 and (3.7), $w_{0}=\lim _{j \rightarrow \infty} \widetilde{d}_{0} \widetilde{t}_{0} \lambda^{m_{j}} r^{n_{j}}$. This implies that

$$
\lim _{j \rightarrow \infty}\left(\frac{m_{j}}{n_{j}} \log \lambda+\log r\right)=\lim _{j \rightarrow \infty} \frac{1}{n_{j}} \log \frac{w_{0}}{\widetilde{d}_{0} \widetilde{t}_{0}}=0
$$



Figure 3.9: The images of the parallel straight segments $\gamma_{k}^{\natural}$ in $D_{a}(p)$ by $h$.
and hence $\lim _{j \rightarrow \infty} \frac{m_{j}}{n_{j}}=-\frac{\log r}{\log \lambda}$. Applying the same argument to $\gamma_{m_{j}, n_{j}}^{\dagger}$, we also have $\lim _{j \rightarrow \infty} \frac{m_{j}}{n_{j}}=-\frac{\log r^{\prime}}{\log \lambda^{\prime}}$. This shows the part (D1) of Theorem 3.1.

Now we will prove the part (D2). For any $n \in \mathbb{N} \cup\{0\}$, we set $f^{n}\left(\gamma_{\infty}^{\natural}\right)=\gamma_{\infty, n}^{\natural}$ and $f^{\prime n}\left(\gamma_{\infty}^{\prime 4}\right)=\gamma_{\infty, n}^{\prime \natural}$. By Corollary 3.8,

$$
\begin{equation*}
h\left(\gamma_{\infty, n}^{\natural}\right)=h\left(f^{n}\left(\gamma_{\infty}^{\natural}\right)\right)=f^{\prime n}\left(h\left(\gamma_{\infty}^{\natural}\right)\right)=f^{\prime n}\left(\gamma_{\infty}^{\prime \natural}\right)=\gamma_{\infty, n}^{\prime} . \tag{3.11}
\end{equation*}
$$

We identify $\operatorname{St}_{a}(p)$ with the unit circle $S^{1}=\{z \in \mathbb{C} ;|z|=1\}$ by corresponding $l \in \operatorname{St}_{a}(p)$ to $e^{\sqrt{-1} \vartheta(l)}$. Then the action of $f$ on $\operatorname{St}_{a}(p)$ is equal to the $\theta$-rotation $R_{\theta}$ on $S^{1}$ defined by $R_{\theta}(z)=e^{\sqrt{-1} \theta} z$.

If $\theta / 2 \pi=v / u$ for coprime positive integers $u, v$ with $0 \leq v<u$. Since $h\left(\gamma_{\infty}^{\natural}\right)=\gamma_{\infty}^{\prime}$, we have $f^{\prime k}\left(\gamma_{\infty}^{\prime 4}\right) \neq \gamma_{\infty}^{\prime}$, for $k=1, \ldots, u-1$ and $f^{\prime \prime}\left(\gamma_{\infty}^{\prime \prime}\right)=\gamma_{\infty}^{\prime \prime}$. This implies that $\theta^{\prime} / 2 \pi=v^{\prime} / u$ for some $v^{\prime} \in \mathbb{N}$ with $0 \leq v^{\prime}<u$. Since $\left.h\right|_{D_{a}(p)}: D_{a}(p) \rightarrow D_{a^{\prime}}\left(p^{\prime}\right)$ is a homeomorphism with the correspondence $h\left(R_{\theta}^{k}\left(\gamma_{\infty}^{\natural}\right)\right)=R_{\theta^{\prime}}^{k}\left(\gamma_{\infty}^{\prime}\right)(k=0,1, \ldots, u-1)$, there exists an orientationpreserving homeomorphism $\eta_{0}: S^{1} \rightarrow S^{1}$ with $\eta_{0}\left(e^{\sqrt{-1}\left(\theta^{\natural}+k \theta\right)}\right)=e^{\sqrt{-1}\left(\theta^{\text {h }}+k \theta^{\prime}\right)}$ for $k=$ $0,1, \ldots, u-1$. We set $\Gamma=\left\{e^{\sqrt{-1}\left(\theta^{\natural}+k \theta\right)} ; k=0,1, \ldots, u-1\right\}$ and $\Gamma^{\prime}=\left\{e^{\sqrt{-1}\left(\theta^{\natural}+k \theta^{\prime}\right)} ; k=\right.$ $0,1, \ldots, u-1\}$. Then $\left[e^{\sqrt{-1} \theta^{\natural}}, e^{\sqrt{-1}\left(\theta^{\natural}+\theta\right)}\right) \cap \Gamma$ consists of $v$ points, where $[a, b)$ denotes the positively oriented half-open interval in $S^{1}$ for $a, b \in S^{1}$ with $a \neq b$. Since moreover $\eta_{0}\left(\left[e^{\sqrt{-1} \theta^{\natural}}, e^{\sqrt{-1}\left(\theta^{\natural}+\theta\right)}\right) \cap \Gamma\right)=\left[e^{\sqrt{-1} \theta^{\natural}}, e^{\sqrt{-1}\left(\theta^{\natural}+\theta^{\prime}\right)}\right) \cap \Gamma^{\prime}$ consists of $v^{\prime}$ points, it follows that $v=v^{\prime}$, and hence $\theta=\theta^{\prime}$.

Next we suppose that $\theta / 2 \pi$ is irrational. Then, for any $l \in \operatorname{St}_{a}(p)$, there exists a subsequence $\left\{n_{k}\right\}$ of $\mathbb{N}$ such that the sequence $\gamma_{\infty, n_{k}}^{\natural}$ uniformly converges to $l$ as $k \rightarrow \infty$. By (3.11), $\gamma_{\infty, n_{k}}^{\prime \emptyset}$ uniformly converges to $l^{\prime}=h(l)$. Since $\gamma_{\infty}^{\prime \natural}, n_{k} \in \operatorname{St}_{a^{\prime}}\left(p^{\prime}\right), l^{\prime}$ is also an element of $\mathrm{St}_{a^{\prime}}\left(p^{\prime}\right)$. Thus we have a homeomorphism $\eta: S^{1} \rightarrow S^{1}$ with respect to which $R_{\theta}$ and $R_{\theta^{\prime}}$ are conjugate. Since the rotation number is invariant under topological conjugations, $\theta / 2 \pi=\theta^{\prime} / 2 \pi \bmod 1$ holds. This completes the proof of the part (D2).

### 3.5 Proof of Theorem 3.3

In this section, we will prove Theorem 3.3. Suppose that $f, f^{\prime}$ are elements of $\operatorname{Diff}^{r}(M)$ satisfying the conditions of Theorems 3.1 and $\theta / 2 \pi$ is irrational.

Since $\theta=\theta^{\prime} \bmod 2 \pi$, for any $k, j \in \mathbb{N}$,

$$
\begin{equation*}
\vartheta\left(\gamma_{\infty, k}^{\natural}\right)-\vartheta\left(\gamma_{\infty, j}^{\natural}\right)=\vartheta\left(\gamma_{\infty, k}^{\prime \natural}\right)-\vartheta\left(\gamma_{\infty, j}^{\prime \natural}\right)=(k-j) \theta \quad \bmod 2 \pi . \tag{3.12}
\end{equation*}
$$

Let $l_{j}(j=1,2)$ be any elements of $\mathrm{St}_{a}(p)$. As in the proof of Theorem 3.1, there exist subsequences $\left\{n_{k}\right\},\left\{n_{j}\right\}$ of $\mathbb{N}$ such that the sequencers $\left\{\gamma_{\infty, n_{k}}^{\natural}\right\},\left\{\gamma_{\infty, n_{j}}^{\natural}\right\}$ uniformly converge to $l_{1}$ and $l_{2}$ respectively. Then, $\left\{\gamma_{\infty, n_{k}}^{\prime \natural}\right\},\left\{\gamma_{\infty}^{\prime}, n_{j}\right\}$ also uniformly converge to the elements $l_{1}^{\prime}=h\left(l_{1}\right)$ and $l_{2}^{\prime}=h\left(l_{2}\right)$ of $\mathrm{St}_{a^{\prime}}\left(p^{\prime}\right)$ respectively. Then, by (3.12),

$$
\begin{equation*}
\vartheta\left(l_{2}\right)-\vartheta\left(l_{1}\right)=\vartheta\left(l_{2}^{\prime}\right)-\vartheta\left(l_{1}^{\prime}\right) \quad \bmod 2 \pi . \tag{3.13}
\end{equation*}
$$

For the proof of Theorem 3.3, we need another family of straight segments in $D_{a}(p)$. Fix an integer $a_{0}$ with

$$
a_{0}>\max \left\{\frac{\log (2 r)}{\log \left(\lambda^{-1}\right)}, \frac{\log \left(2 r^{\prime}\right)}{\log \left(\lambda^{\prime-1}\right)}\right\} .
$$

For any $k \geq 0$, we consider the straight segment $\xi_{k}^{\natural}=f^{k}\left(\gamma_{a_{0} k}^{\natural}\right) \cap D_{a}(p)$. By (3.9),

$$
\begin{equation*}
\vartheta\left(\xi_{k}^{\natural}\right)-\vartheta\left(\xi_{0}^{\natural}\right)=k \theta \quad \bmod 2 \pi \quad \text { and } \quad d\left(0, \xi_{k}^{\natural}\right)=w_{0} \lambda^{a_{0} k} r^{k}<2^{-k} w_{0} . \tag{3.14}
\end{equation*}
$$

Similarly, by (3.10), $\xi_{k}^{\prime \natural}=h\left(\xi_{k}^{\natural}\right)$ is a straight segment in $D_{a^{\prime}}\left(p^{\prime}\right)$ with

$$
\begin{equation*}
\vartheta\left(\xi_{k}^{\prime 4}\right)-\vartheta\left(\xi_{0}^{\prime \boldsymbol{4}}\right)=k \theta \quad \bmod 2 \pi \quad \text { and } \quad d\left(0, \xi_{k}^{\prime 4}\right)=w_{0}^{\prime} \lambda^{\prime a_{0} k} r^{\prime k}<2^{-k} w_{0}^{\prime} . \tag{3.15}
\end{equation*}
$$

Proof of Theorem 3.3. Let $\alpha$ be the element of $\operatorname{St}_{a}(p)$ with $\vartheta\left(\xi_{0}^{\natural}\right)-\vartheta(\alpha)=\pi / 2$ and $\alpha^{\prime}=$ $h(\alpha) \in \mathrm{St}_{a^{\prime}}\left(p^{\prime}\right)$. We will show that $\theta_{\alpha^{\prime}}:=\vartheta\left(\xi_{0}^{\prime \natural}\right)-\vartheta\left(\alpha^{\prime}\right)$ is also equal to $\pi / 2 \bmod 2 \pi$. See Figure 3.10. In fact, since $\theta / 2 \pi$ is irrational, by (3.14) there exists a subsequence


Figure 3.10: Correspondence of straight segments via $h$.
$\xi_{k_{j}}^{\natural}$ uniformly converges to $\alpha$. Since $\left.h\right|_{D_{a}(p)}$ is uniformly continuous, $\xi_{k_{j}}^{\prime \natural}$ also uniformly converges to $\alpha^{\prime}$. On the other hand, since $\vartheta\left(\xi_{k_{j}}^{\natural}\right)-\vartheta(\alpha)=k_{j} \theta+\pi / 2 \bmod 2 \pi$ and $\vartheta\left(\xi_{k_{j}}^{\prime}\right)-$ $\vartheta\left(\alpha^{\prime}\right)=k_{j} \theta+\theta_{\alpha^{\prime}} \bmod 2 \pi$,

$$
\theta_{\alpha^{\prime}}-\frac{\pi}{2}=\left(\vartheta\left(\xi_{k_{j}}^{\prime \natural}\right)-\vartheta\left(\alpha^{\prime}\right)\right)-\left(\vartheta\left(\xi_{k_{j}}^{\natural}\right)-\vartheta(\alpha)\right) \rightarrow 0 \bmod 2 \pi
$$

as $j \rightarrow \infty$. Thus we have $\theta_{\alpha^{\prime}}=\pi / 2 \bmod 2 \pi$.
We denote by $z(\boldsymbol{x}) \in \mathbb{C}$ the entry of $\boldsymbol{x} \in D_{a}(p)$ with respect to the linearizing coordinate on $D_{a}(p)$. Similarly, the entry of $\boldsymbol{x}^{\prime} \in D_{a^{\prime}}\left(p^{\prime}\right)$ is denoted by $z^{\prime}\left(\boldsymbol{x}^{\prime}\right)$. Let $\boldsymbol{x}_{0}$ be the intersection point of $\alpha$ and $\xi_{0}^{\natural}$, and let $\boldsymbol{x}_{0}^{\prime}=h\left(\boldsymbol{x}_{0}\right)$. One can set $z\left(\boldsymbol{x}_{0}\right)=\rho_{0} e^{\sqrt{-1} \omega_{0}}$ and $z^{\prime}\left(\boldsymbol{x}_{0}^{\prime}\right)=\rho_{0}^{\prime} e^{\sqrt{-1} \omega_{0}^{\prime}}$ for some $\rho_{0}>0, \rho_{0}^{\prime}>0$ and $\omega_{0}, \omega_{0}^{\prime} \in \mathbb{R}$. We define the new linearizing coordinate on $D_{a^{\prime}}\left(p^{\prime}\right)$ by using the linear conformal map such that, for any $\boldsymbol{x}^{\prime} \in D_{a^{\prime}}\left(p^{\prime}\right)$, $z^{\prime \text { new }}\left(\boldsymbol{x}^{\prime}\right)=\rho_{0} \rho_{0}^{\prime-1} e^{\sqrt{-1}\left(\omega_{0}-\omega_{0}^{\prime}\right)} z^{\prime}\left(\boldsymbol{x}^{\prime}\right)$. Then $z\left(\boldsymbol{x}_{0}\right)=z^{\prime \text { new }}\left(\boldsymbol{x}_{0}^{\prime}\right)$ holds.

For any $\boldsymbol{x} \in \xi_{0}^{\natural}$, there exists $l \in \operatorname{St}_{a}(p)$ with $\{\boldsymbol{x}\}=\xi_{0}^{\natural} \cap l$. Then $\boldsymbol{x}^{\prime}=h(\boldsymbol{x})$ is the intersection of $\xi_{0}^{\prime \dagger}$ and $l^{\prime}=h(l)$. By (3.13), $\vartheta(l)-\vartheta(\alpha)=\vartheta\left(l^{\prime}\right)-\vartheta\left(\alpha^{\prime}\right) \bmod 2 \pi$ and hence $z(\boldsymbol{x})=z^{\prime \text { new }}\left(\boldsymbol{x}^{\prime}\right)$. We say the property that $h$ is identical on $\xi_{0}^{\natural}$. Since $\theta / 2 \pi$ is irrational, there exists $k_{*} \in \mathbb{N}$ satisfying

$$
\frac{\pi}{3} \leq \vartheta\left(\xi_{k_{*}}^{\natural}\right)-\vartheta\left(\xi_{0}^{\natural}\right) \leq \frac{\pi}{2} \quad \bmod 2 \pi .
$$

Then $\xi_{k_{*}}^{\natural}$ meets $\xi_{0}^{\natural}$ at a single point $\boldsymbol{x}_{k_{*}}$ in $D_{a}(p)$. For $\alpha_{k_{*}}=f^{k_{*}}(\alpha)$ and $\alpha_{k_{*}}^{\prime}=h\left(\alpha_{k_{*}}\right)$, we have $\vartheta\left(\xi_{k_{*}}^{\natural}\right)-\vartheta\left(\alpha_{k_{*}}\right)=\vartheta\left(\xi_{k_{*}}^{\prime \prime}\right)-\vartheta\left(\alpha_{k_{*}}^{\prime}\right)=\pi / 2$. Since $h$ is identical at $\boldsymbol{x}_{k_{*}}, h$ is proved to be identical on $\xi_{k_{*}}^{\natural}$ by an argument as above. Then one can show inductively that, for any $n \in \mathbb{N}, h$ is identical on $\xi_{n k_{*}}^{\natural}$. See Figure 3.11. By (3.14), $\lim _{n \rightarrow \infty} d\left(0, \xi_{n k_{*}}^{\natural}\right)=0$. Since


Figure 3.11: Correspondence via $h$ with respect to the new coordinate on $D_{a^{\prime}}\left(p^{\prime}\right)$.
moreover $k_{*} \theta / 2 \pi$ is irrational, $\overline{\bigcup_{n=1}^{\infty} \xi_{n k_{*}}^{\natural}}$ is equal to $D_{a}(p)$. This shows that $h$ is identical on $D_{a}(p)$. In particular, this implies that $\left.h\right|_{D_{a}(p)}$ is a linear conformal map with respect to the original coordinates. We write $z(q)=\rho_{1} e^{\sqrt{-1} \omega_{1}}$ and $z^{\prime}\left(q^{\prime}\right)=\rho_{1}^{\prime} e^{\sqrt{-1} \omega_{1}^{\prime}}$. It follows
from the assumption of $h(q)=q^{\prime}$ in our theorems that $h(z)=\rho_{1}^{\prime} \rho_{1}^{-1} e^{\sqrt{-1}\left(\omega_{1}^{\prime}-\omega_{1}\right)} z$ for any $\underset{\sim}{z} \in \mathbb{C}$ with $|z| \leq a$. In particular, this implies that $\left.h\right|_{W_{\text {loc }}^{u}(p)}$ is a linear conformal map. Let $\widetilde{h}$ be any other conjugacy homeomorphism between $f$ and $f^{\prime}$ satisfying the conditions in Theorems 3.1 and 3.3. In particular, $\widetilde{h}(p)=p^{\prime}$ and $\widetilde{h}(q)=q^{\prime}$ hold. Since $z(q)=\rho_{1} e^{\sqrt{-1} \omega_{1}}$ and $z^{\prime}\left(q^{\prime}\right)=\rho_{1}^{\prime} e^{\sqrt{-1} \omega_{1}^{\prime}}$, one can show as above that $\widetilde{h}(z)=\rho_{1}^{\prime} \rho_{1}^{-1} e^{\sqrt{-1}\left(\omega_{1}^{\prime}-\omega_{1}\right)} z$ for any $z \in \mathbb{C}$ with $|z| \leq a$ and hence $\left.\widetilde{h}\right|_{D_{a}(p)}=\left.h\right|_{D_{a}(p)}$. This shows the assertion (E2) of Theorem 3.3 and $r=r^{\prime}$. Then, by the assertion (D1) of Theorem 3.1, we also have $\lambda=\lambda^{\prime}$. This completes the proof.

Let $\widehat{z}$ be the homoclinic transverse point of $W^{u}(p)$ and $W^{s}(p)$ given in Subsection 3.2.1. Fix a sufficiently large $n \in \mathbb{N}$ with $s=f^{-n}(\widehat{z}) \in D_{p}(a)$. Then $s^{\prime}=h(s)$ is contained in $D_{b^{\prime}}\left(p^{\prime}\right)$. The following corollary shows that $z(s) / z(q)$ is a modulus for $f$. Recall that $z(\boldsymbol{x}) \in \mathbb{C}$ is the entry of $\boldsymbol{x}$ with respect to the complex linearizing coordinate on $D_{a}(a)$. The complex number $z^{\prime}\left(\boldsymbol{x}^{\prime}\right)$ is defined similarly for $\boldsymbol{x}^{\prime} \in D_{a^{\prime}}\left(p^{\prime}\right)$.

Corollary 3.9. Let $f$, $f^{\prime}$ be elements of $\operatorname{Diff}^{r}(M)$ satisfying the conditions of Theorems 3.1 and 3.3, and let $h$ be a conjugacy homeomorphism between $f$ and $f^{\prime}$ with $h(p)=p^{\prime}$ and $h(q)=q^{\prime}$. If $\left.h\right|_{W_{\text {loc }}^{u}(p)}$ is orientation-preserving, then $z(s) / z(q)=z^{\prime}\left(s^{\prime}\right) / z^{\prime}\left(q^{\prime}\right)$. Otherwise, $z(s) / z(q)=\overline{z^{\prime}\left(s^{\prime}\right) / z^{\prime}\left(q^{\prime}\right)}$.

Proof. Here we only consider the case that $h$ is orientation-preserving. Since $\left.h\right|_{D_{a}(p)}$ is a linear conformal map, the triangle with vertices $0, z(q), z(s)$ is similar to that with vertices $0, z^{\prime}\left(q^{\prime}\right), z^{\prime}\left(s^{\prime}\right)$ with respect to Euclidean geometry. This shows $z(s) / z(q)=z^{\prime}\left(s^{\prime}\right) / z^{\prime}\left(q^{\prime}\right)$.

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[1] S. Hashimoto, Moduli of surface diffeomorphisms with cubic tangencies, to appear in Tokyo J. Math.
[2] S. Hashimoto, S. Kiriki, and T. Soma, Moduli of 3-dimensional diffeomorphisms with saddle foci, Discrete Cont. Dynam. Sys. 38 (2018), No. 10, 5021-5037.

