JP Journal of Geometry and Topology
© 2014 Pushpa Publishing House, Allahabad, India
Published Online: February 2015
Available online at http://www.pphmj.com/journals/jpgt.htm Volume 16, Number 2, 2014, Pages 97-125

# POINT-PUSHING PSEUDO-ANOSOV MAPPING CLASSES AND THEIR ACTIONS ON THE CURVE COMPLEX 

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#### Abstract

Let $S$ be an analytically finite Riemann surface of type ( $p, n$ ) with $3 p+n>4$, which is equipped with a hyperbolic metric and contains at least one puncture $x$. Let $\mathcal{C}(S)$ be the curve complex endowed with a path metric $d_{\mathcal{C}}$. It is known that a point-pushing pseudo-Anosov mapping class $f$ on $S$ determines a filling closed geodesic $c$ on $S \cup\{x\}$, and that every non-preperipheral vertex $u$ in $\mathcal{C}(S)$ determines a simple curve $\widetilde{u}$ on $S \cup\{x\}$. In this paper, we consider the action of $f$ on $\mathcal{C}(S)$. We describe all geodesic segments in $\mathcal{C}(S)$ that connect $u$ and $f(u)$ when $d_{\mathcal{C}}(u, f(u))=2$. We also give sufficient conditions with respect to the intersection points between $c$ and $\tilde{u}$ for the path distance $d_{\mathcal{C}}(u, f(u))$ to be larger.


Received: August 8, 2014; Revised: September 28, 2014; Accepted: October 25, 2014
2010 Mathematics Subject Classification: Primary 32G15; Secondary 30F60.
Keywords and phrases: Riemann surfaces, pseudo-Anosov, Dehn twists, curve complex, filling curves.

Communicated by Yasuo Matsushita

## 1. Introduction

Let $S$ be an analytically finite Riemann surface of type $(p, n)$ with $3 p+n>4$, where $p$ is the genus and $n$ is the number of punctures on $S$. Assume that $S$ is equipped with a hyperbolic metric and contains at least one puncture $x$. In [5], Harvey introduced a curve complex $\mathcal{C}(S)$ on $S$, which is a simplicial complex where vertices are simple closed geodesics and a $k$ th dimensional simplex, denoted by $\mathcal{C}_{k}(S)$, is a collection of $k+1$ disjoint simple closed geodesics on $S$.

Let $\mathcal{C}_{0}^{\prime}(S)$ be the subset of $\mathcal{C}_{0}(S)$ consisting of boundary geodesics of twice punctured disks enclosing $x\left(\mathcal{C}_{0}^{\prime}(S)\right.$ is not empty if and only if $\left.n \geq 2\right)$. Let $\hat{\mathcal{C}}(S)$ denote the subcomplex of $\mathcal{C}(S)$ which consists of non-preperipheral simplexes, where a simplex $\left\{u_{0}, \ldots, u_{k}\right\}$ is called non-preperipheral if none of $u_{i}$ belongs to $\mathcal{C}_{0}^{\prime}(S)$. It is easily seen that $\mathcal{C}_{0}(S)=\hat{\mathcal{C}}_{0}(S)$ if $S$ is of type ( $p, n$ ) with $p \geq 2$ and $n=0$, 1 . Otherwise, we have $\mathcal{C}_{0}(S) \backslash \hat{\mathcal{C}}_{0}(S)=\mathcal{C}_{0}^{\prime}(S)$. Note that any two vertices in $\mathcal{C}_{0}^{\prime}(S)$ intersect, which means that $\mathcal{C}_{0}^{\prime}(S)$ is totally disconnected.

Write $\widetilde{S}=S \cup\{x\}$, and let $\widetilde{S}$ be equipped with a hyperbolic metric. The curve complex $\mathcal{C}(\widetilde{S})$ can be similarly defined so that there is a natural projection:

$$
\begin{equation*}
\varepsilon: \hat{\mathcal{C}}(S) \rightarrow \mathcal{C}(\widetilde{S}) . \tag{1.1}
\end{equation*}
$$

According to Birman and Series [3], we may choose a point $x$ that misses every simple closed geodesic on $\widetilde{S}$, which means that a vertex in $\mathcal{C}(\widetilde{S})$ can also be regarded as a vertex on $\mathcal{C}(S)$ (by simply removing the point $x$ ). Hence (1.1) admits a global section.

For any $u, v \in \mathcal{C}_{0}(S)$, the distance $d_{\mathcal{C}}(u, v)$ between $u$ and $v$ is defined as the minimum number of edges in $\mathcal{C}_{1}(S)$ joining $u$ and $v$, thereby $\mathcal{C}(S)$ is equipped with the path metric $d_{\mathcal{C}}$. Clearly, $d_{\mathcal{C}}(u, v)=1$ if and only if $u$,
$v$ are disjoint, and $d_{\mathcal{C}}(u, v) \geq 3$ if and only if $(u, v)$ fills $S$ in the sense that every closed geodesic intersects $u$ or $v$. It is well-known that $\mathcal{C}(S)$ is connected and is $\delta$-hyperbolic in the sense of Gromov [4]. See Masur and Minsky [8] for a detailed explanation.

Let $\mathscr{F}$ be the set of pseudo-Anosov mapping classes of $S$ onto itself that fix the puncture $x$ and are isotopic to the identity on $\widetilde{S}$ (see Thurston [10] for the definition and basic properties of pseudo-Anosov maps). Let $I: \widetilde{S} \times[0,1] \rightarrow \widetilde{S}$ denote the associated isotopy between an element $f \in \mathscr{F}$ and the identity, i.e., $I(\cdot, 0)=f$ and $I(\cdot, 1)=$ id. Then $I(x, t), t \in[0,1]$, defines an oriented closed curve $c$ on $\widetilde{S}$ passing through $x$. It was shown in Kra [7] that $c$ is freely homotopic to an oriented filling closed geodesic (call it $c$ also) and every such a geodesic determines a conjugacy class $K(c)$ of $c$ in $\mathscr{F}$. Let $\mathscr{S}$ denote the set of primitive oriented filling closed geodesics on $\widetilde{S}$. Obviously, $\mathscr{F}$ can be partitioned into conjugacy classes and there is a bijection between the set of these conjugacy classes and the set $\mathscr{S}$.

By Proposition 3.6 of [8], there exists a constant $a>0$, which depends only on ( $p, n$ ) such that for all pseudo-Anosov maps $f: S \rightarrow S$, all positive integers $m$, and all $u \in \mathcal{C}_{0}(S)$, it holds that $d_{\mathcal{C}}\left(u, f^{m}(u)\right) \geq a m$. By contrast, it was shown in $[14,16]$ that if $f \in \mathscr{F}$, then $d_{\mathcal{C}}\left(u, f^{m}(u)\right) \geq m$ for $1 \leq m \leq 4$. In [15], we showed that $d_{\mathcal{C}}(u, f(u)) \geq 3$ provided that $u \in \mathcal{C}_{0}^{\prime}(S)$.

The purpose of this paper is to study some questions on the distance $d_{\mathcal{C}}(u, f(u))$ for $u \in \hat{\mathcal{C}}_{0}(S)$, which is related to the set $\#\{\widetilde{u}, c\}$ of intersection points between $\widetilde{u} \in \mathcal{C}_{0}(\widetilde{S})$ and $c \in \mathscr{S}$, where and throughout the rest of the paper, we use the symbol $\widetilde{u}$ to denote $\varepsilon(u)$, which is the geodesic homotopic to $u$ on $\widetilde{S}$ if $u$ is also viewed as a curve on $\widetilde{S}$. The main results will be stated in the next section.

This paper is organized as follows: In Section 2, we state our main results. In Section 3, we collect some background information on the fibration $\varepsilon: \mathcal{C}(S) \rightarrow \mathcal{C}_{0}(\widetilde{S})$ between curve complexes; we then investigate the fibers $F_{\widetilde{u}}$ of the fibration and discuss some properties by means of a tessellation of $\mathbf{H}$. We show that there is an intimate relationship between $F_{\widetilde{u}}$ and a tessellation of $\mathbf{H}$. In Section 4, we study the distance between vertices in each fiber $F_{\widetilde{u}}$ in terms of the intersection numbers between the filling geodesics and vertices, and prove Theorem 2.1 and Theorem 2.2. The proof of Theorem 2.3 is given in Section 5.

## 2. Main Theorems

Fix $\widetilde{u} \in \mathcal{C}_{0}(\widetilde{S})$ and $c \in \mathscr{S}$. The geometric intersection number $i(c, \widetilde{u})$ between $c$ and $\widetilde{u}$ is defined as the number of points in $\#\{\widetilde{u}, c\}$, which is also given by

$$
i(c, \widetilde{u})=\min \left|c^{\prime} \cap \widetilde{u}^{\prime}\right|,
$$

where $c^{\prime}$ and $\widetilde{u}^{\prime}$ are in the homotopy classes of $c$ and $\widetilde{u}$, respectively.
It is known that $i(c, \tilde{u})=1$ if and only if $d_{\mathcal{C}}(u, f(u))=1$ for some $f \in K(c)$ and $u \in\left\{\varepsilon^{-1}(\widetilde{u})\right\}$ (Lemma 4.3).

Consider the case where $d_{\mathcal{C}}(u, f(u))=2$ for $u \in \hat{\mathcal{C}}_{0}(S)$. Denote by

$$
V(u, f)=\left\{v \in \mathcal{C}_{0}(S): d_{\mathcal{C}}(u, v)=1 \text { and } d_{\mathcal{C}}(f(u), v)=1\right\}
$$

For $\widetilde{u}=\varepsilon(u)$, we define a subset $W(\widetilde{u}, c)$ of $\mathcal{C}_{0}(\widetilde{S})$ as follows. Let $k=i(c, \widetilde{u})$ and let $\left\{Q_{1}, \ldots, Q_{k}\right\}$ be the intersection points between $\widetilde{u}$ and $c$. Let $\tilde{v} \in \mathcal{C}_{0}(\widetilde{S})$ be such that $\widetilde{v} \neq \tilde{u}$. Let $\left\{P_{1}, \ldots, P_{r}\right\}$ denote the intersection points between $c$ and $\widetilde{v}$. With the aid of a parametrization $c=c(\tau)$ for $\tau \in[0,1]$, we can write $P_{i}=c\left(\tau_{i}^{1}\right)$ and $Q_{j}=c\left(\tau_{j}^{2}\right)$ for some $\tau_{i}^{1}, \tau_{j}^{2} \in$ $(0,1)$. We call $\widetilde{u}$ and $\widetilde{v}$ are separated by $c$ if there is a parametrization $c=c(\tau)$ such that $\tau_{i}^{1}<\tau_{j}^{2}$ for all $1 \leq i \leq r$ and $1 \leq j \leq k$. Let $W^{*}(\widetilde{u}, c)$
denote the set of vertices $\tilde{v} \in \mathcal{C}_{0}(\widetilde{S})$ such that $\tilde{v}$ and $\tilde{u}$ are separated by $c$. It is clear that $W^{*}(\widetilde{u}, c)=\mathcal{C}_{0}(\widetilde{S})$ if $i(\widetilde{u}, c)=1$. Also, if $i(\widetilde{v}, c)=1$, then $\widetilde{v} \in W^{*}(\widetilde{u}, c)$.

Let $W(\widetilde{u}, c)$ be the subset of $W^{*}(\widetilde{u}, c)$ consisting of vertices $\widetilde{v}$ disjoint from $\widetilde{u}$. We first prove the following result.

Theorem 2.1. Let $f \in \mathscr{F}$ and let $c \in \mathscr{S}$ be determined by $f$. Let $u \in \hat{\mathcal{C}}_{0}(S)$. Assume that $d_{\mathcal{C}}(u, f(u))=2$. Then $V(u, f) \cap \mathcal{C}_{0}^{\prime}(S)=\varnothing$ and the puncture-forgetting projection (1.1) restricts to a bijection $V(u, f) \simeq$ $W(\widetilde{u}, c)$.

We now study the problem of when $d_{\mathcal{C}}(u, f(u)) \geq 3$. It seems likely that the distance between $u$ and $f(u)$ is proportional to the number of intersections between $\widetilde{u}$ and $c$. Unfortunately, this is not true. In fact, from the discussion in Section 3, there exist geodesics $c \in \mathscr{S}$ that intersect some $\widetilde{u}$ as many times as we expect, yet the distance $d_{\mathcal{C}}(u, f(u))$ remains small.

Nevertheless, with respect to a vertex $\widetilde{u} \in \mathcal{C}_{0}(\widetilde{S})$, any $c \in \mathscr{S}$ with $i(c, \widetilde{u})>1$ can be written as a curve concatenation $c_{1} \cdot c_{2}$ (from left to right), where $c_{1}$ and $c_{2}$ are constructed as follows. Given an orientation for $\tilde{u}$ and recall that $\#\{c \cap \tilde{u}\}=\left\{Q_{1}, \ldots, Q_{k}\right\}$. Let $c=c(\tau), \tau \in[0,1]$, be a parametrization and let $Q_{i}=c\left(\tau_{i}^{2}\right)$. Assume that $\tau_{1}^{2}<\tau_{2}^{2}<\cdots<\tau_{k}^{2}$. Let $\alpha$ be the path in $c$ connecting $Q_{1}$ and $Q_{k}$, and let $\beta$ be the path in $\widetilde{u}$ joining $Q_{k}$ and $Q_{1}$. Then the complement $\gamma=c \backslash \alpha$ is also a path in $c$ that joins $Q_{k}$ and $Q_{1}$. Let $c_{1}$ be the path $\alpha$ followed by $\beta$, and $c_{2}$ be the path $\beta^{-1}$ followed by $\gamma$. Obviously, we have $c=\alpha \cdot \gamma=(\alpha \cdot \beta) \cdot\left(\beta^{-1} \cdot \gamma\right)=c_{1} \cdot c_{2}$.

Let $\delta_{1}, \delta_{2}$ be the geodesics freely homotopic to $c_{1}$ and $c_{2}$, respectively, and let $\mathscr{S}^{*}$ be the subset of $\mathscr{S}$ consisting of geodesics that intersect each
$\widetilde{u} \in \mathcal{C}_{0}(\widetilde{S})$ at least twice. It is easy to see that $\delta_{1}, \delta_{2} \in \mathscr{S}$ implies $c \in \mathscr{S}^{*}$, this particularly implies that $\mathscr{S}^{*}$ is not empty. Let $F_{\widetilde{u}}$ denote the set $\left\{\varepsilon^{-1}(\widetilde{u})\right\}$ which consists of vertices $u$ in $\mathcal{C}_{0}(S)$ with $\varepsilon(u)=\widetilde{u}$. Theorem 2.1 leads to the following result.

Theorem 2.2. Let $c \in \mathscr{S}$ be determined by an element $f \in \mathscr{F}$. Let $\widetilde{u} \in \mathcal{C}_{0}(\widetilde{S})$. Suppose that $i(c, \widetilde{u}) \geq 3$ and with respect to $\widetilde{u}$, the filling geodesic $c$ can be expressed as $c_{1} \cdot c_{2}$, where $c_{1}$ and $c_{2}$ are freely homotopic to nontrivial geodesics $\delta_{1}$ and $\delta_{2}$ which satisfies the condition that $\delta_{1} \in \mathscr{S}$. Then $d_{\mathcal{C}}(u, f(u)) \geq 3$ for all $u \in F_{\widetilde{u}}$.

Remark. In the case where $i(\widetilde{u}, c)=2$, Lemma 3.3 together with Lemma 4.1 asserts that $d_{\mathcal{C}}(u, f(u)) \geq 3$ for most $u \in F_{\widetilde{u}}$, but for the remaining ones in $F_{\widetilde{u}}$, we have $d_{\mathcal{C}}(u, f(u))=2$.

Let $f, c$ and $\widetilde{u}$ be as in Theorem 2.2. Then $c=c_{1} \cdot c_{2}$, where $c_{1}$ can further be written as the concatenation $c_{1}=c_{11} \cdots c_{1, k-1}$, where $c_{11}$ is obtained from the path in $c$ joining $Q_{1}$ and $Q_{2}$, followed by the path $\beta_{1}$ in $\widetilde{u}$ joining $Q_{2}$ and $Q_{1}$. Inductively, for each $1 \leq i \leq k-1$, let $c_{1 i}$ be the closed curve obtained from the path $\beta_{i}^{-1}$ in $\tilde{u}$ joining $Q_{1}$ and $Q_{i}$, followed by the path in $c$ joining $Q_{i}$ and $Q_{i+1}$, then followed by the path $\beta_{i+1}$ in $\widetilde{u}$ joining $Q_{i+1}$ and $Q_{1}$. Let $\delta_{1 i}$ be the geodesic representative in the homotopy class of $c_{1 i}$.

Theorem 2.3. Let $f, c$ and $\widetilde{u}$ be as above. Assume that $k=i(c, \widetilde{u}) \geq 4$. If $c$ can be expressed as $\left(c_{11} \cdots c_{1, k-1}\right) \cdot c_{2}$, where at least two curves $c_{1 i}$ and $c_{1 j}, 1 \leq i, j \leq k-1$ and $|i-j| \geq 2$, are homotopic to filling closed geodesics, then there are $u \in F_{\widetilde{u}}$ such that $d_{\mathcal{C}}(u, f(u)) \geq 4$.

There arises the following question:

Question. Is there a primitive $f \in \mathscr{F}$ such that $d_{\mathcal{C}}(u, f(u))$ is arbitrarily large for $u \in \mathcal{C}(S)$ ?

## 3. Background and Tree Structures of Fibers in the Curve Complex

Let $\mathbf{H}$ denote the upper half plane that is endowed with the hyperbolic metric $\rho(z)|d z|=|d z| / y$. Let $\varrho: \mathbf{H} \rightarrow \widetilde{S}$ be the universal covering map with the covering group $G$. Then $G$ is a finitely generated Fuchsian group of the first kind and acts on $\mathbf{H}$ as a group of isometries with respect to $\rho(z)$. The group $G$ is torsion free and contains infinitely many hyperbolic Möbius transformations. $G$ also contains parabolic Möbius transformations if and only if $\widetilde{S}$ contains punctures $(n>1)$.

Choose $\hat{x} \in \mathbf{H}$ with $\varrho(\hat{x})=x$. Let $\mathscr{A}=\{h(\hat{x}): h \in G\} \subset \mathbf{H}$ and let $\dot{G}$ be the covering group of a universal covering map $\varrho^{\prime}: \mathbf{H} \rightarrow S$. Then $\mathbf{H} / \dot{G}$ $\cong S \cong(\mathbf{H} / G) \backslash\{x\}$ and there exists an exact sequence

$$
1 \rightarrow \Gamma \rightarrow \dot{G} \rightarrow G \rightarrow 1
$$

where $\Gamma$ is the covering group of a universal covering map $v: \mathbf{H} \rightarrow \mathbf{H} \backslash \mathscr{A}$.
Let $Q(G)$ (resp. $Q(\dot{G})$ ) be the group of quasiconformal automorphisms $w$ of $\mathbf{H}$ with $w G w^{-1}=G$ (resp. $w \dot{G} w^{-1}=\dot{G}$ ). Two elements $w, w_{1} \in Q(G)$ are said to be equivalent if $\left.w\right|_{\mathbf{R}}=\left.w_{1}\right|_{\mathbf{R}}$. Let $[w]$ be the equivalence class of $w$. The $x$-pointed mapping class group, denoted by $\operatorname{Mod}_{S}^{x}$, is the subgroup of the ordinary mapping class group $\operatorname{Mod}(S)$ that consists of mapping classes fixing $x$. In [1], Bers explicitly constructed an isomorphism $\varphi^{*}: Q(G) / \sim \rightarrow$ $\operatorname{Mod}_{S}^{x}$ (Theorem 10 of [1]), which is outlined below.

Let $Q_{0}(\dot{G}) \subset Q(\dot{G})$ be the subgroup of $Q(\dot{G})$ consisting of maps projecting (under $\varrho^{\prime}$ ) to maps on $S$ leaving the puncture $x$ fixed. For any $[w] \in Q(G) / \sim$, there is a map $w_{0} \in Q(G)$, which is obtained by performing
a local quasiconformal deformation within each fundamental region $D$ leaving the boundary $\partial D$ fixed, such that $w \sim w_{0}$ and $w_{0}(\hat{x}) \in \mathscr{A}$. Clearly, as a map of $\mathbf{H} \backslash \mathscr{A}, w_{0}$ can be lifted (through $\varrho^{\prime}$ ) to a map $\omega_{0} \in Q_{0}(\dot{G})$. Hence $\varphi^{*}([w])$ can be defined as the mapping class on $S$ represented by the projection of $\omega_{0}$ under $\varrho^{\prime}$.

Alternatively, let $T(\widetilde{S})$ denote the Teichmüller space of $\widetilde{S}$ and let $F(\widetilde{S})$ be the fiber space over $T(\widetilde{S})$ (see [1] for the definitions of $T(\widetilde{S})$ and $F(\widetilde{S})$ ). By Theorem 9 of [1], there is a biholomorphic map $\varphi: F(\widetilde{S}) \rightarrow T(S)$ that respects the forgetting map of $T(S)$ onto $T(\widetilde{S})$. By Theorem 1 of [1], $Q(G) / \sim$ is considered a group of holomorphic automorphisms of $F(\widetilde{S})$. Via $\varphi$ it thus defines a group of holomorphic automorphisms of $T(S)$, which is, by Royden's theorem [9], identified with the group $\operatorname{Mod}_{S}^{x}$.

Note that $G$ is centerless, for any $g, g^{\prime} \in G, g=g^{\prime}$ if and only if $g h g^{-1}=g^{\prime} h\left(g^{\prime}\right)^{-1}$ for all $h \in G$. We see that $G$ can be regarded as a normal subgroup of $Q(G) / \sim$. Thus, $\varphi^{*}$ restricts to an isomorphism of $G$ onto $\varphi^{*}(G)$ $\subset \operatorname{Mod}_{S}^{x}$. Write $[w]^{*}=\varphi^{*}([w])$ for $[w] \in Q(G) / \sim$ and $g^{*}=\varphi^{*}(g)$ for $g \in G$.

We now construct some special elements of $Q(G) \backslash G$. Let $\mathscr{N}$ be the disjoint union of small crescent neighborhoods of all geodesics in $\left\{\varrho^{-1}(\widetilde{u})\right\}$. Let $\widetilde{\mathscr{N}}=\varrho(\mathscr{N})$ and $t_{\widetilde{u}}$ the positive Dehn twist along $\widetilde{u}$ which is supported in $\widetilde{\mathscr{N}}$. Let $\tau: \mathbf{H} \rightarrow \mathbf{H}$ be a lift of $t_{\widetilde{u}}$. That is, $\tau$ satisfies the two conditions (i) $\tau G \tau^{-1}=G$ and (ii) $\varrho \tau=t_{\widetilde{u}} \varrho$. One can easily verify that $\tau \in Q(G)$ and thus that $[\tau]^{*} \in \operatorname{Mod}_{S}^{x}$.

Fix $\tilde{u} \in \mathcal{C}_{0}(\widetilde{S})$. Let $\left\{\varrho^{-1}(\widetilde{u})\right\}$ denote the collection of all (disjoint) geodesics $\hat{u} \subset \mathbf{H}$ such that $\varrho(\hat{u})=\widetilde{u}$. Denote by $\mathscr{R}_{\widetilde{u}}$ the collection of all
components of $\mathbf{H} \backslash\left\{\varrho^{-1}(\widetilde{u})\right\}$ and by $F_{\widetilde{u}}$ the subset of $\mathcal{C}_{0}(S)$ that projects to $\widetilde{u}$ under the projection (1.1). For any $\Omega \in \mathscr{B}_{\tilde{u}}$, there exists a unique hyperbolic element $g \in G$ such that $\left.g^{-1} \tau\right|_{\Omega \backslash \mathscr{N}}=$ id (see [13] for more explanations). Let $\tau_{\Omega}=g^{-1} \tau$. Then $\tau_{\Omega} \in Q(G)$ and satisfies conditions (i) and (ii) above, which tells us that $\left[\tau_{\Omega}\right]^{*} \in \operatorname{Mod}_{S}^{x}$.

The complement of $\Omega$ in $\mathbf{H}$ consists of infinitely many half-planes, which are invariant half-planes under the action of $\tau$ and are called maximal (or first order) elements determined by $\tau$. Each maximal element contains infinitely many second order half-planes which are not invariant by the action of $\tau$, and so on (the actions of $\tau$ on higher order half-planes were investigated in [11]). Let $\mathscr{U}$ be the collection of all these half-planes of different orders. Then $\mathscr{U}$ is partially ordered defined by inclusion. That is, for any $n \geq 1$ and any element $\Delta_{n+1}$ of $\mathscr{U}$ with $(n+1)$ th order, there is a unique element $\Delta_{n} \in \mathscr{U}$ with $n$th order, such that $\Delta_{n+1} \subset \Delta_{n}$.

By Lemma 3.2 of [13], $\left[\tau_{\Omega}\right]^{*}=t_{u}$ for a $u \in F_{\widetilde{u}}$, where $t_{u}$ denotes the positive Dehn twist along $u$ on $S$. Conversely, for every $u \in F_{\widetilde{u}}$, there is a component $\Omega \in \mathscr{R}_{\widehat{u}}$ such that $\left[\tau_{\Omega}\right]^{*}=t_{u}$. In what follows, we write $\tau=\tau_{\Omega}$ and call the triple $(\tau, \Omega, \mathscr{U})$ the configuration corresponding to the vertex $u \in F_{\widetilde{u}}$. It is clear that for each $h \in G, \tau_{h(\Omega)}=h\left(g^{-1} \tau\right) h^{-1}$. Hence $\left[\tau_{h(\Omega)}\right]^{*}$

$$
\begin{gather*}
=\left[h\left(g^{-1} \tau\right) h^{-1}\right]^{*}=h^{*}\left[\tau_{\Omega}\right]^{*}\left(h^{*}\right)^{-1}=t_{h^{*}(u)} . \text { We thus obtain a map } \\
\chi: \mathscr{R}_{\widetilde{u}} \rightarrow F_{\widetilde{u}} \tag{3.1}
\end{gather*}
$$

defined by sending $\Omega$ to $u$. It is not difficult to prove that $\chi$ is surjective. To see that $\chi$ is also injective, we note that for any two regions $\Omega, \Omega^{\prime} \in \mathscr{R}_{\widetilde{u}}$ with $\Omega \neq \Omega^{\prime}$, they are either adjacent, by which we mean that $\Omega$ and $\Omega^{\prime}$ share a common geodesic boundary $e$ for $e \in\left\{\varrho^{-1}(\widetilde{u})\right\}$, or $\Omega$ and $\Omega^{\prime}$ are disjoint. If $\Omega$ and $\Omega^{\prime}$ are adjacent, then by Lemma 2.1 of [17],
$d_{\mathcal{C}}\left(\chi\left(\Omega_{1}\right), \chi\left(\Omega_{2}\right)\right)=1$, which occurs if and only if $\left\{\chi\left(\Omega_{1}\right), \chi\left(\Omega_{2}\right)\right\}$ are boundary components of an $x$-punctured cylinder on $S$. If $\Omega$ and $\Omega^{\prime}$ are disjoint, then there are maximal elements $\Delta \in \mathscr{U}$ and $\Delta^{\prime} \in \mathscr{U}^{\prime}$ such that $\Delta U \Delta^{\prime}=\mathbf{H}$ and $\partial \Delta \cap \partial \Delta^{\prime}=\varnothing$. By examining the action of $\tau$ and $\tau^{\prime}$ on $\mathbf{S}^{1}=\partial \mathbf{H}$, we conclude that $\tau^{n} \tau^{\prime m} \neq \tau^{\prime m} \tau^{n}$ for all large integers $m, n$. See Lemma 4 of [11] for more details. This implies that $t_{\chi(\Omega)}^{n} t^{m}{ }_{\chi\left(\Omega^{\prime}\right)} \neq$ $t_{\chi\left(\Omega^{\prime}\right)}^{m} \eta_{\chi(\Omega)}^{n}$. In particular, $\chi\left(\Omega^{\prime}\right) \neq \chi(\Omega)$.

We have thus proved the following lemma:
Lemma 3.1. The map $\chi$ defined as (3.1) is a bijection which satisfies the equivariant condition $\chi(g(\Omega))=g^{*}(\chi(\Omega))$ for each $g \in G$ and each $\Omega \in \mathscr{R}_{\widetilde{u}}$.

Remark. The bijection (3.1) can also be obtained from topological terms by Theorem 7.1 of Kent et al. [6]. In fact, Theorem 4.1 and Theorem 4.2 of Birman [2] state that there exists an exact sequence

$$
\begin{equation*}
1 \rightarrow \pi_{1}(\widetilde{S}, x) \rightarrow \operatorname{Mod}_{S}^{x} \rightarrow \operatorname{Mod}(\widetilde{S}) \rightarrow 1 \tag{3.2}
\end{equation*}
$$

where $\operatorname{Mod}(\widetilde{S})$ is the mapping class group and $\pi_{1}(\widetilde{S}, x)$ is the fundamental group of $\widetilde{S}$. Thus, there defines an injective map of $\pi_{1}(\widetilde{S}, x)$ into $\operatorname{Mod}_{S}^{x}$. Note that an isomorphism between $G$ and $\pi_{1}(\widetilde{S}, x)$ is obtained by choosing a lift $\hat{x}$ of $x$ in $\mathbf{H}$ to serve as the base point. As such, we obtain an injective map $\psi: G \rightarrow \operatorname{Mod}_{S}^{x}$ so that the image $\psi(G)$ consists of mapping classes projecting to the trivial mapping class on $\widetilde{S}$ as $x$ is filled in. Since $G$ keeps $\varrho^{-1}\{\widetilde{u}\}$ invariant, $G$ naturally acts on $\mathscr{R}_{\widetilde{u}}$. Hence $\psi(G)$ acts on $F_{\widetilde{u}}$. As such, the bijection $\chi$ in Lemma 3.1 can be given by identifying a region $\Omega$ that is the stabilizer of an $h \in G$ with a vertex $u \in F_{\widetilde{u}}$ that is the stabilizer of $\psi(h)$.

By a similar discussion to the above, we know that $\mathscr{R}_{\widetilde{u}}$ naturally inherits a structure of a tree. From the bijection (3.1), we know that $F_{\widetilde{u}}$ is also a tree. Hence $F_{\widetilde{u}}$ is path connected. We formulate the result as the following lemma:

Lemma 3.2. For each $\tilde{u} \in \mathcal{C}_{0}(\widetilde{S})$, the fiber $F_{\widetilde{u}}$ is path connected in $F_{\widetilde{u}}$. Moreover, for any $u, v \in F_{\widetilde{u}}$, there is one and only one path in $F_{\widetilde{u}}$ connecting $u$ and $v$.

For $\Omega, \Omega^{\prime} \in \mathscr{R}_{\widetilde{u}}$, we can introduce the distance $D\left(\Omega, \Omega^{\prime}\right)$ between $\Omega$ and $\Omega^{\prime}$ as follows. If there are elements $\Omega_{1}, \ldots, \Omega_{n-1} \in \mathscr{R}_{\widetilde{u}}$ which are obtained from Lemma 3.2, we then declare $D\left(\Omega, \Omega^{\prime}\right)=n$.

Lemma 3.3. We have $d_{\mathcal{C}}\left(\chi(\Omega), \chi\left(\Omega^{\prime}\right)\right) \leq D\left(\Omega, \Omega^{\prime}\right)$, and if $D\left(\Omega, \Omega^{\prime}\right)$ $=2$, then $d_{\mathcal{C}}\left(\chi(\Omega), \chi\left(\Omega^{\prime}\right)\right)=2$.

Proof. From Lemma 3.2, there is a sequence $\Omega_{0}=\Omega, \Omega_{1}, \ldots, \Omega_{n}=\Omega^{\prime}$ in $\mathscr{R}_{\widetilde{u}}$ such that for $0 \leq i \leq n-1, \Omega_{i}$ is adjacent to $\Omega_{i+1}$. By Lemma 2.1 of [17], we get $d_{\mathcal{C}}\left(\chi\left(\Omega_{i}\right), \chi\left(\Omega_{i+1}\right)\right)=1$. It follows from the triangle inequality that $d_{\mathcal{C}}\left(\chi(\Omega), \chi\left(\Omega^{\prime}\right)\right) \leq \sum_{i=0}^{n-1} d_{\mathcal{C}}\left(\chi\left(\Omega_{i}\right), \chi\left(\Omega_{i+1}\right)\right)=n$. Hence $d_{\mathcal{C}}\left(\chi(\Omega), \chi\left(\Omega^{\prime}\right)\right)$ $\leq D\left(\Omega, \Omega^{\prime}\right)$.

If $D\left(\Omega, \Omega^{\prime}\right)=2$, then $d_{\mathcal{C}}\left(\chi(\Omega), \chi\left(\Omega^{\prime}\right)\right) \leq 2$. Assume that $d_{\mathcal{C}}(\chi(\Omega)$, $\left.\chi\left(\Omega^{\prime}\right)\right)=1$, which says that $\chi(\Omega)$ and $\chi\left(\Omega^{\prime}\right)$ are disjoint. Thus, $\left\{\chi(\Omega), \chi\left(\Omega^{\prime}\right)\right\}$ are boundary components of an $x$-punctured cylinder on $S$. By Lemma 2.1 of [17], $\Omega$ and $\Omega^{\prime}$ are adjacent. But this means that $D\left(\Omega, \Omega^{\prime}\right)=1$. This is a contradiction.

## 4. Distances between Vertices in Fibers in the Curve Complex

For each $\tilde{u} \in \mathcal{C}_{0}(\widetilde{S})$, the group $G$ naturally acts on $\mathscr{R}_{\widetilde{u}}$. Let $g \in G$ be
an essential hyperbolic element, i.e., $g^{*} \in \mathscr{F}$. Let $\operatorname{axis}(g)$ be the axis of $g$; that is, $\operatorname{axis}(g) \subset \mathbf{H}$ is the unique geodesic so that $g(\operatorname{axis}(g))=\operatorname{axis}(g)$. Let $\Omega \in \mathscr{R}_{\widetilde{u}}$ and denote by $u=\chi(\Omega)$. We first investigate the situation where $\operatorname{axis}(g) \cap \Omega=\varnothing$. We have

Lemma 4.1. Suppose that $\Omega \cap \operatorname{axis}(g)=\varnothing$. Then $d_{\mathcal{C}}\left(u, g^{*}(u)\right) \geq 3$.
Proof. The proof can be found in Section 4 of [14].
Observe that for any $\Omega \in \mathscr{R}_{\widehat{u}}$, there always exists an essential hyperbolic element $g$ (and hence there are infinitely many such elements) in the conjugacy class of $g$ for which axis $(g) \cap \Omega \neq \varnothing$, and conversely, for any essential hyperbolic element $g \in G$, there are infinitely many $\Omega \in \mathscr{R}_{\widetilde{u}}$ such that $\operatorname{axis}(g) \cap \Omega \neq \varnothing$. We now proceed to explore the relationship between $D(\Omega, g(\Omega))$ and the intersection number between $\operatorname{axis}(g)$ and $\tilde{u}$.

Lemma 4.2. Let $g \in G$ be an essential hyperbolic element, and let $\Omega \subset \mathscr{R}_{\widetilde{u}}$ be a region so that $\operatorname{axis}(g) \cap \Omega \neq \varnothing$. Then $i(\widetilde{u}, \varrho(\operatorname{axis}(g)))=$ $D(\Omega, g(\Omega))$.

Proof. Write $n=i(\widetilde{u}, \varrho(\operatorname{axis}(g)))$. Let $\Delta \in \mathscr{U}_{u}$ be the maximal element that covers the repelling fixed point of $g$. By Lemma 2.1 of [14], there is a maximal element $\Delta_{1} \in \mathscr{U}_{u}$ that contains $g(\mathbf{H} \backslash \Delta) . \Delta_{1}$ is disjoint from $\Delta$ and covers the attracting fixed point of $g$. Thus, $\Omega \subset \mathbf{H} \backslash\left(\Delta \cup \Delta_{1}\right)$. Parametrize $c(t)=\operatorname{axis}(g),-\infty<t<+\infty$, so that $c(t) \cap \partial \Delta=c\left(t_{0}\right)$. Also, write $z=$ $c\left(t_{0}\right)$. Let $c\left(t_{1}\right)=c(t) \cap \partial \Delta_{1}$. Then $c(t)$ continues to travel and enters a second-order element $\Delta_{2} \in \mathscr{U}_{u}$. Set $c\left(t_{2}\right)=c(t) \cap \partial \Delta_{2}$, and so on.

When $c(t)$ reaches the point $g(z)$ at time $t_{n}, c\left(t_{n}\right)=g(z)$ lies in $\partial \Delta_{n}$ for an $n$th order element $\Delta_{n} \in \mathscr{U}_{u}$. We thus obtain a finite sequence $-\infty<$ $t_{0}<t_{1}<t_{2}<\cdots<t_{n}<+\infty$, and see that $c(t), t \in\left[t_{0}, t_{n}\right]$ meets $n$ regions
of $\mathscr{R}_{\widetilde{u}}$. It is easy to check that $g(\Omega)$ lies in $\Delta_{n} \in \mathscr{U}_{u}$. It follows that $n=$ $D(\Omega, g(\Omega))$, as claimed.

Consider the case in which $d_{\mathcal{C}}\left(u, g^{*}(u)\right)=1$.
Lemma 4.3. Let $g \in G$ be an essential hyperbolic element. Let $\tilde{u} \in$ $\mathcal{C}_{0}(\widetilde{S})$. Then $i(\widetilde{u}, \varrho(\operatorname{axis}(g)))=1$ if and only if there is $u \in F_{\widetilde{u}}$ such that $d_{\mathcal{C}}\left(u, g^{*}(u)\right)=1$, in which case $u=\chi(\Omega)$ for an $\Omega \in \mathscr{R}_{\widetilde{u}}$ with $\Omega \cap$ $\operatorname{axis}(g) \neq \varnothing$.

Proof. If $i(\widetilde{u}, \varrho(\operatorname{axis}(g)))=1$, then by Lemma 4.2, $D(\Omega, g(\Omega))=1$ for an $\Omega \subset \mathscr{R}_{\widetilde{u}}$ with $\Omega \bigcap \operatorname{axis}(g) \neq \varnothing$. This means that $\Omega$ and $g(\Omega)$ are adjacent. By Lemma 2.1 of [17], $d_{\mathcal{C}}(\chi(\Omega), \chi(g(\Omega)))=1$. Therefore, if we let $u=\chi(\Omega)$, then

$$
d_{\mathcal{C}}(u, f(u))=d_{\mathcal{C}}\left(\chi(\Omega), g^{*}(\chi(\Omega))\right)=d_{\mathcal{C}}(\chi(\Omega), \chi(g(\Omega)))=1
$$

Conversely, suppose $d_{\mathcal{C}}(u, f(u))=1$. Let $\Omega \subset \mathscr{R}_{\widetilde{u}}$ be such that $u=$ $\chi(\Omega)$. By Lemma 3.1, we have $d_{\mathcal{C}}(\chi(\Omega), \chi(g(\Omega)))=d_{\mathcal{C}}\left(\chi(\Omega), g^{*}(\chi(\Omega))\right)$ $=1$. By Lemma 2.1 of [17] again, $\Omega$ and $g(\Omega)$ are adjacent. If $\Omega \cap$ $\operatorname{axis}(g)=\varnothing$, then by Lemma 4.1, $d_{\mathcal{C}}\left(\chi(\Omega), g^{*}(\chi(\Omega))\right) \geq 3$, which leads to a contradiction. We thus conclude that $\Omega \bigcap \operatorname{axis}(g) \neq \varnothing$. By Lemma 4.2, $i(\widetilde{u}, \varrho(\operatorname{axis}(g)))=1$.

The case where $d_{\mathcal{C}}(u, f(u))=2$ is more involved. Let $c=\varrho(\operatorname{axis}(g))$. If $i(\widetilde{u}, c)=1$, then by Lemma 4.1 and Lemma 4.3, either $d_{\mathcal{C}}(u, f(u)) \geq 3$ or $d_{\mathcal{C}}(u, f(u))=1$. So it must be the case that $i(\tilde{u}, c) \geq 2$. It is worthwhile pointing out that the number $i(\tilde{u}, c)$ can be arbitrarily large. In fact, one can easily construct a geodesic $c \in \mathscr{S} \backslash \mathscr{S}^{*}$ on a genus $p$ surface for $p>1$ such
that a vertex $\widetilde{v} \in \mathcal{C}_{0}(\widetilde{S})$ can be found with the property that $i(c, \widetilde{v})=1$ and $\widetilde{v}$ is disjoint from $\widetilde{u}$. Let $f \in \mathscr{F}$ be defined by $c$ by a point-pushing deformation, let $g \in G$ be such that $g^{*}=f$. In this case, $d_{\mathcal{C}}(u, f(u))=2$ for some $f \in K(c)$ and some $u \in F_{\widetilde{u}}$ (the author thanks the referee for his comments).

Figure 1 depicts such a filling curve $c$ which intersects $\widetilde{u}$ and $\widetilde{v}$, and goes around a hole as many times as possible. This illustrates that $i(\widetilde{u}, c)=$ $i(\widetilde{u}, \varrho(\operatorname{axis}(g)))$ can be arbitrarily large, whereas $d_{\mathcal{C}}(u, f(u))=2$ for $u=$ $\chi(\Omega)$, where $\Omega \in \mathscr{R}_{\widetilde{u}}$ with $\Omega \cap \operatorname{axis}(g) \neq \varnothing$. In particular, $i(\widetilde{u}, c)-$ $d_{\mathcal{C}}(u, f(u))$ could be unbounded.

Note that it could be the case where both $i(\widetilde{u}, c)$ and $d_{\mathcal{C}}(u, f(u))$ are larger than two even if $\Omega \bigcap \operatorname{axis}(g) \neq \varnothing$. However, if $i(\widetilde{u}, c)=2$, then Lemma 4.1 tells us that $d_{\mathcal{C}}(\chi(\Omega), \chi(g(\Omega))) \geq 3$ for all $\Omega \subset \mathscr{R}_{\widetilde{u}}$ with $\Omega \cap$ $\operatorname{axis}(g)=\varnothing$. But if $\Omega$ is so chosen that $\Omega \bigcap \operatorname{axis}(g) \neq \varnothing$, then by Lemma 4.2 and Lemma 3.3, $d_{\mathcal{C}}(u, f(u))=2$, where $u=\chi(\Omega)$.


Figure 1
Recall that $(\tau, \Omega, \mathscr{U})$ is the configuration corresponding to $u$. For other vertices $v \in \hat{\mathcal{C}}_{0}(S)$, we let $\left(\tau_{v}, \Omega_{v}, \mathscr{U}_{v}\right)$ denote the configuration corresponding to $v$. It is easy to verify that $g(\Omega)=\Omega_{w}$ if ( $\left.\tau_{w}, \Omega_{w}, \mathscr{U}_{w}\right)$ is the configuration corresponding to $w=g^{*}(u)$.

Assume that $\Omega \cap \operatorname{axis}(g) \neq \varnothing$. By Lemma 2.1 of [14], there is a pair $\left(\Delta, \Delta^{\prime}\right)$ of maximal elements of $\mathscr{U}_{u}$ such that $\Delta$ covers the repelling fixed point and $\Delta^{\prime}$ covers the attracting fixed point of $g$. In addition, we have $\Omega \subset \mathbf{H} \backslash\left(\Delta \cup \Delta^{\prime}\right)$ and $\bar{\Omega} \cap g(\bar{\Omega})=\varnothing$. Without loss of generality, we may assume that $g(\Delta) \cap \Delta^{\prime} \neq \varnothing$.

The following lemma will be used frequently in the proof of Theorem 2.1 and Theorem 2.2.

Lemma 4.4. With the above notations and conditions, we have $d_{\mathcal{C}}(v, f(u)) \geq 2$ provided that there is a maximal element $\Delta_{v} \in \mathscr{U}_{v}$ such that $\mathbf{H} \backslash \Delta_{v} \subset g(\Delta)$. Similarly, $d_{\mathcal{C}}(v, u) \geq 2$ if there is a maximal element $\Delta_{v} \in \mathscr{U}_{v}$ such that $\mathbf{H} \backslash \Delta_{v} \subset \Delta^{\prime}$. In particular, if $v$ is disjoint from both $u$ and $f(u)$, then no maximal element of $\mathscr{U}_{v}$ exists whose boundary lies entirely in $g(\Delta) \cap \Delta^{\prime}$.

Proof. We outline the argument here for completeness. See [14] for more details. If $d_{\mathcal{C}}(v, w)=1$, where $w=f(u)$, then $t_{w}(v)=v$. So for any maximal element $\Delta_{v}^{\prime} \in \mathscr{U}_{v}, t_{w}^{s}\left(\Delta_{v}^{\prime}\right) \in \mathscr{U}_{v}$ is also a maximal element for any integer $s \neq 0$. On the other hand, by the assumption, there is a maximal element $\Delta_{v} \in \mathscr{U}_{v}$ such that $\mathbf{H} \backslash \Delta_{v} \subset g(\Delta)$. By examining the action of $\tau_{w}$ on $\mathbf{H}$, we see that $\tau_{w}^{s}\left(\mathbf{H} \backslash \Delta_{v}\right)$ is disjoint from $\mathbf{H} \backslash \Delta_{v}$ for any $s \neq 0$. It follows that $\tau_{w}^{s}\left(\Delta_{v}\right) \cap \Delta_{v} \neq \varnothing$ and $\Delta_{v} \neq \tau_{w}^{s}\left(\Delta_{v}\right)$. We thus conclude that $\tau_{w}^{s}\left(\Delta_{v}\right)$ is not a maximal element of $\mathscr{U}_{v}$. This is a contradiction. Similarly, we can prove the other two statements.

Alternatively, the condition that there is a maximal element $\Delta_{v} \in \mathscr{U}_{v}$ such that $\mathbf{H} \backslash \Delta_{v} \subset g(\Delta)$ implies that $\bar{\Omega}_{v} \cap \bar{\Omega}_{w}=\varnothing$, which says that $\Omega_{v}$ and $\Omega_{w}$ are disjoint but not adjacent. The result is stated as follows for future references.

Lemma 4.5. Let $u, v \in \hat{\mathcal{C}}_{0}(S)$ be such that $d_{\mathcal{C}}(u, v)=1$. Let $\left(\tau_{u}, \Omega_{u}, \mathscr{U}_{u}\right)$ and $\left(\tau_{v}, \Omega_{v}, \mathscr{U}_{v}\right)$ be the configurations corresponding to $u$ and $v$, respectively. Then $\bar{\Omega}_{u}$ and $\bar{\Omega}_{v}$ cannot be disjoint.

Let $u^{\prime} \in \mathcal{C}_{0}^{\prime}(S)$. Then $\widetilde{u}^{\prime}$ is trivial. By Lemma 5.1 and Lemma 5.2 of [15], $u^{\prime}$ corresponds to a parabolic fixed point $z^{\prime}$ of $G$; that is, there is a primitive parabolic element $T \in G$ such that $T\left(z^{\prime}\right)=z^{\prime}$. If $d_{\mathcal{C}}\left(u, u^{\prime}\right)=1$, then $t_{u} t_{u^{\prime}}=t_{u^{\prime}} t_{u}$. Thus, $\tau_{u} T=T \tau_{u}$, which in turn implies that $\tau_{u}$ fixes $z^{\prime}$, and thus that $z^{\prime} \in \Omega_{u} \cap \mathbf{S}^{1}$. The converse remains valid and the result can be stated in the following lemma (see [15] for more detailed argument):

Lemma 4.6. Let $u \in \hat{\mathcal{C}}_{0}(S)$ and $u^{\prime} \in \mathcal{C}_{0}^{\prime}(S)$. Let $z^{\prime} \in \mathbf{S}^{1}$ be the parabolic fixed point of $G$ corresponding to $u^{\prime}$. Then $d_{\mathcal{C}}\left(u, u^{\prime}\right)=1$ if and only if $z^{\prime} \in \Omega_{u} \cap \mathbf{S}^{1}$.

Return to the case in which $\Omega=\Omega_{u}$ and $d_{\mathcal{C}}(\chi(\Omega), \chi(g(\Omega)))=2$ but $D(\Omega, g(\Omega))$ is arbitrary (where we recall that $u=\chi(\Omega)$, and $g^{*}=f$ ).

Lemma 4.7. Let $\widetilde{u} \in \mathcal{C}_{0}(\widetilde{S})$ be a non-trivial vertex, and let $g \in G$ be an essential hyperbolic element so that $i(\widetilde{u}, \varrho(\operatorname{axis}(g))) \geq 2$. Suppose that $\Omega \in \mathscr{R}_{\widetilde{u}}$ is such that $d_{\mathcal{C}}(\chi(\Omega), \chi(g(\Omega)))=2$. Then $V\left(\chi(\Omega), g^{*}\right) \cap \mathcal{C}_{0}^{\prime}(S)$ $=\varnothing$. Furthermore, we have $\varepsilon\left(V\left(\chi(\Omega), g^{*}\right)\right) \subset W(\widetilde{u}, \varrho(\operatorname{axis}(g)))$, where $\varepsilon$ is the projection given in (1.1).

Proof. If $\Omega \bigcap \operatorname{axis}(g)=\varnothing$, then by Lemma 4.1, we see that $d_{\mathcal{C}}(\chi(\Omega)$, $\chi(g(\Omega))) \geq 3$. This contradicts the hypothesis. So it must be the case that $\Omega \cap \operatorname{axis}(g) \neq \varnothing$. Write $u=\chi(\Omega)$. By Lemma 3.1, $\chi(g(\Omega))=g^{*}(\chi(\Omega))=$ $g^{*}(u)=f(u)$. Since $d_{\mathcal{C}}(u, f(u))=2,(u, f(u))$ does not fill $S$. Choose $v \in$ $V(u, f)$, which means that $d_{\mathcal{C}}(v, u)=d_{\mathcal{C}}(v, f(u))=1$.

If $\widetilde{v}=\varepsilon(v)$ is trivial, i.e., $v \in \mathcal{C}_{0}^{\prime}(S)$, then by Lemma 4.6, $z \in$ $(\Omega \cap g(\Omega)) \cap \mathbf{S}^{1}$. This is impossible since $\bar{\Omega} \cap g(\bar{\Omega})=\varnothing$. This proves $V\left(\chi(\Omega), g^{*}\right) \cap \mathcal{C}_{0}^{\prime}(S)=\varnothing$.

We thus assume that $v \in \hat{\mathcal{C}}_{0}(S)$. Then $\widetilde{v} \in \mathcal{C}_{0}(\widetilde{S})$. Suppose that $\widetilde{v}$ is not in $W(\tilde{u}, \varrho(\operatorname{axis}(g)))$. By Lemma 4.4, we need to find a maximal element $\Delta_{v} \in \mathscr{U}_{v}$ such that $\mathbf{H} \backslash \Delta_{v} \subset \Delta^{\prime}$ or $\mathbf{H} \backslash \Delta_{v} \subset g(\Delta)$. To do so, we first notice that $g \in G$ is an essential hyperbolic element. As a result, axis $(g)$ cannot lie in $\Omega_{v}$. There remain two cases:

Case 1. There is a maximal element $\Delta_{v} \in \mathscr{U}_{v}$ such that $\operatorname{axis}(g) \subset \Delta_{v}$. Since $\tilde{u}$ is disjoint from $\tilde{v}$ and since $\partial \Delta_{v} \in\left\{\varrho^{-1}(\tilde{v})\right\}$, we see that $\partial \Delta_{v}$ is disjoint from all geodesics drawn in Figure 2(a).

Notice that $\partial \Delta^{\prime}, \partial \Delta, \partial g(\Delta) \in\left\{\varrho^{-1}(\widetilde{u})\right\}$, so they are all mutually disjoint. As it turns out, $\left(\mathbf{H} \backslash \Delta_{v}\right) \cap \mathbf{S}^{1}$ lies in one of the eight subarcs of $\mathbf{S}^{1} \backslash\{A, B, E$, $F, H, L, M, N\}$. A possible $\Delta_{v}$ is drawn in Figure 2(a). As we see, $\mathbf{H} \backslash \Delta_{v}$ $\subset g(\Delta)$ or $\mathbf{H} \backslash \Delta_{v} \subset \Delta_{v}^{\prime}$, and we are done.

Case 2. $\operatorname{axis}(g) \cap \Omega_{v} \neq \varnothing$. There is a pair $\left(\Delta_{v}, \Delta_{v}^{\prime}\right)$ of maximal elements of $\mathscr{U}_{v}$, where $\Delta_{v}$ covers the repelling fixed point of $g$ and $\Delta_{v}^{\prime}$ covers the attracting fixed point of $g$, so that $\Omega_{v} \subset \mathbf{H} \backslash\left(\Delta_{v} \cup \Delta_{v}^{\prime}\right)$.

By hypothesis, $\widetilde{v}$ is non-trivial and does not lie in $W(\widetilde{u}, \varrho(\operatorname{axis}(g)))$. Write $\#\{c(\tau), \widetilde{v}\}=\left\{P_{1}, \ldots, P_{r}\right\}$. There is a parametrization $c=c(\tau)$ and an intersection point $P_{q} \in\left\{P_{1}, \ldots, P_{r}\right\}, P_{q}=c\left(\tau_{q}^{1}\right)$, such that $\tau_{1}^{2}<\tau_{q}^{1}$, where $c\left(\tau_{1}^{2}\right)$ is the first intersection point in $\#\{c(\tau), \tilde{u}\} . P_{q}$ corresponds to the intersection point $\hat{P}_{q}$ shown in Figure 2(b).

(a)

(b)

Figure 2
Since $u, v$ are disjoint, $\widetilde{u}, \widetilde{v}$ are also disjoint, which implies that $\left\{\varrho^{-1}(\widetilde{u})\right\}$ and $\left\{\varrho^{-1}(\widetilde{v})\right\}$ are disjoint. But $\partial \Delta, \partial \Delta^{\prime}, \partial g(\Delta) \in\left\{\varrho^{-1}(\widetilde{u})\right\}$. We conclude that there is one geodesic $[X Y] \in\left\{\varrho^{-1}(\widetilde{v})\right\}$ that passes through $\hat{P}_{q}$ and lies entirely in the region $g(\Delta) \cap \Delta^{\prime}$.

If $\partial \Delta_{v} \subset \mathbf{H} \backslash g(\Delta)$, then $\mathbf{H} \backslash \Delta_{v} \subset \Delta^{\prime}$. If $\partial \Delta_{v} \subset g(\Delta) \cap \Delta^{\prime}$ (note that $\partial \Delta_{v}$ may or may not be equal to $[X Y]$ ), then also $\mathbf{H} \backslash \Delta_{v} \subset \Delta^{\prime}$. We are done. If $\partial \Delta_{v} \subset \mathbf{H} \backslash\left(\Delta \cup \Delta^{\prime}\right)$, then since $[X Y] \subset g(\Delta) \cap \Delta^{\prime}, \Delta_{v}^{\prime}$ must contain $[X Y]$, which tells us that $\mathbf{H} \backslash \Delta_{v}^{\prime} \subset g(\Delta)$. Finally if $\partial \Delta_{v} \subset \Delta$, then $\Delta_{v}^{\prime}$ contains [ $X Y$ ]. It follows that $\mathbf{H} \backslash \Delta_{v}^{\prime} \subset g(\Delta)$.

It remains to consider the case where $\partial \Delta_{v} \in\left\{\partial \Delta, \partial \Delta^{\prime}, \partial g(\Delta)\right\}$. If $\partial \Delta_{v}=$ $\partial \Delta$, then $t_{v}=t_{u}$. Thus, $u=v$. If $\partial \Delta_{v}=\partial g(\Delta)$, then $t_{w}=t_{v}$. So $w=v$. If $\partial \Delta_{v}=\partial \Delta^{\prime}$, then $\Delta_{v}^{\prime}$ contains [XY]. It follows that $\mathbf{H} \backslash \Delta_{v}^{\prime} \subset g(\Delta)$. Geometrically, in this case, $v=u_{0}$, where $u_{0}$ together with $u$ forms the boundary of an $x$-punctured cylinder.

Proof of Theorem 2.1. By Lemma 4.7, we see that the restriction $\left.\varepsilon\right|_{V(u, f)}$ sends $V(u, f)$ to $W(\widetilde{u}, c)$. We first claim that $\left.\varepsilon\right|_{V(u, f)}$ is onto.

Indeed, pick $\widetilde{v} \in W(\widetilde{u}, c)$, and let $c=c(\tau)$ be a parametrization for $c$ such that

$$
\begin{equation*}
\tau_{i}^{1}<\tau_{j}^{2} \text { for all } 1 \leq i \leq r \text { and } 1 \leq j \leq k \tag{4.1}
\end{equation*}
$$

where $\#\{\tilde{v}, c\}=\left\{c\left(\tau_{1}^{1}\right), \ldots, c\left(\tau_{r}^{1}\right)\right\}$ and $\#\{\tilde{u}, c\}=\left\{c\left(\tau_{1}^{2}\right), \ldots, c\left(\tau_{k}^{2}\right)\right\}$. Take $x=c(0)$ and move $x$ along $c(\tau)$. After meeting $\widetilde{v}$ at $P_{1}, \ldots, P_{r}$ in the order, $c$ starts meeting $\widetilde{u}$ at order $Q_{1}, \ldots, Q_{k}$ and never meets $\widetilde{v}$ again before arriving at $x$ because of (4.1). It turns out that $v$ is disjoint from $g^{*}(u)$ as well (where $v$ is obtained from $\widetilde{v}$ by deleting $x$ ). We see that $v$ is disjoint from $u$ and $f(u)$. By the assumption, $d_{\mathcal{C}}(u, f(u))=2$ and $d_{\mathcal{C}}(u, v)=$ $d_{\mathcal{C}}(f(u), v)=1$. By the triangle inequality, we obtain

$$
2=d_{\mathcal{C}}(u, f(u)) \leq d_{\mathcal{C}}(u, v)+d_{\mathcal{C}}(v, f(u))=2
$$

It follows that $d_{\mathcal{C}}(u, f(u))=d_{\mathcal{C}}(u, v)+d_{\mathcal{C}}(v, f(u))$. Hence $v \in V(u, f)$ and $\varepsilon(v)=\widetilde{v}$. This proves that $\left.\varepsilon\right|_{V(u, f)}$ is onto.

We next claim that $\left.\varepsilon\right|_{V(u, f)}$ is one-to-one. Suppose there are $v_{1}, v_{2} \in$ $V(u, f)$ so that $\varepsilon\left(v_{1}\right)=\varepsilon\left(v_{2}\right)$. Let $\tilde{v}=\varepsilon\left(v_{i}\right), i=1,2$. Let $\left(\tau_{i}, \Omega_{i}, \mathscr{U}_{i}\right)$ be the configurations corresponding to $v_{i}$.

If, say, for $v_{1}$ we have $\Omega_{1} \cap \operatorname{axis}(g)=\varnothing$, then by Lemma 4.1, $d_{\mathcal{C}}\left(v_{1}, f\left(v_{1}\right)\right) \geq 3$. Since $v_{1} \in V(u, f)$, we have $d_{\mathcal{C}}\left(u, v_{1}\right)=d_{\mathcal{C}}\left(f(u), v_{1}\right)$ $=1$. Hence the triangle inequality yields that

$$
\begin{aligned}
2 & =d_{\mathcal{C}}(u, f(u))=d_{\mathcal{C}}\left(u, v_{1}\right)+d_{\mathcal{C}}\left(v_{1}, f(u)\right) \\
& =d_{\mathcal{C}}\left(f(u), f\left(v_{1}\right)\right)+d_{\mathcal{C}}\left(v_{1}, f(u)\right) \geq d_{\mathcal{C}}\left(v_{1}, f\left(v_{1}\right)\right) \geq 3
\end{aligned}
$$

This is a contradiction. Similarly, we assert that $\Omega_{2} \cap \operatorname{axis}(g) \neq \varnothing$. So we only need to handle the case in which $\Omega_{i} \cap \operatorname{axis}(g) \neq \varnothing$. Since $\varepsilon\left(v_{1}\right)=$ $\varepsilon\left(v_{2}\right), v_{1}, v_{2} \in F_{\widetilde{v}}$. From the bijection (3.1), we see that $\Omega_{1}, \Omega_{2} \in \mathscr{R}_{\widetilde{v}}$. As
an immediate consequence, either $\Omega_{1}, \Omega_{2}$ are disjoint, or $\Omega_{1}, \Omega_{2}$ are adjacent, or $\Omega_{1}=\Omega_{2}$.

The condition $v_{1} \in V(u, f)$ says that $v_{1}$ is disjoint from $u$ and $w$, where $w=f(u)$. By the argument of Lemma 4.4, $\tau_{u}^{r} \tau_{w}^{-s}$ sends every maximal element of $\mathscr{U}_{1}$ to a maximal element. In particular, for the maximal element $\Delta_{v_{1}} \in \mathscr{U}_{1}$ that contains the repelling fixed point of $g$ and $\operatorname{axis}(g) \cap \partial \Delta_{v_{1}}$ $\neq \varnothing$, it cannot occur that $\mathbf{H} \backslash \Delta_{v_{1}} \subset g(\Delta)$ or $\mathbf{H} \backslash \Delta_{v_{1}} \subset \Delta^{\prime}$. This further implies that $\bar{\Omega}_{1}$ contains $\operatorname{axis}(g) \cap\left(g(\Delta) \cap \Delta^{\prime}\right)$ (see Figure 2(b)). Similarly, the condition $v_{2} \in V(u, f)$ leads to that $\bar{\Omega}_{2}$ also contains axis $(g) \cap$ $\left(g(\Delta) \cap \Delta^{\prime}\right)$. Hence $\bar{\Omega}_{1} \cap \bar{\Omega}_{2}$ contains axis $(g) \cap\left(g(\Delta) \cap \Delta^{\prime}\right)$. In particular, $\bar{\Omega}_{1} \cap \bar{\Omega}_{2} \neq \varnothing$. This tells us that $\bar{\Omega}_{1}$ and $\bar{\Omega}_{2}$ are not disjoint. If they are adjacent, $\bar{\Omega}_{1} \cap \bar{\Omega}_{2}$ is the axis of a simple hyperbolic element of $G$, which coincides with $\operatorname{axis}(g)$. This implies that $g$ is simple, which is clearly a contradiction. So we must have $\Omega_{1}=\Omega_{2}$, which means $\chi\left(\Omega_{1}\right)=\chi\left(\Omega_{2}\right)$. That is, $v_{1}=v_{2}$. This shows that $\left.\varepsilon\right|_{V(u, f)}$ is one-to-one, and hence the proof of Theorem 2.1 is complete.

Proof of Theorem 2.2. Suppose $d_{\mathcal{C}}(u, f(u)) \leq 2$ for some $u \in F_{\widetilde{u}}$. If $d_{\mathcal{C}}(u, f(u))=1$, then $\chi(\Omega)=u$ and $\Omega$ satisfies $\Omega \cap \operatorname{axis}(g) \neq \varnothing$. By Lemma 4.3, $i(\widetilde{u}, c)=1$. In particular, $c$ cannot be expressed as a concatenation of two filling curves $c_{1}$ and $c_{2}$. Assume that $d_{\mathcal{C}}(u, f(u))=2$. Then by Lemma 4.3 and Lemma 4.1, $i(\tilde{u}, c) \geq 2$ (as discussed in Section 3, the converse is not true, but $i(\widetilde{u}, c)=2$ does imply that $d_{\mathcal{C}}(u, f(u))=2$ for $u=\chi(\Omega) \in F_{\widetilde{u}}$, where $\left.\Omega \cap \operatorname{axis}(g) \neq \varnothing\right)$.

From Theorem 2.1, we know that $V(u, f) \cap \mathcal{C}_{0}^{\prime}(S)=\varnothing$ and $\varepsilon(V(u, f))$ $=W(\widetilde{u}, c)$. Choose $v \in V(u, f)$ and write $\widetilde{v}=\varepsilon(v)$. Since $u, v$ are disjoint, $\widetilde{u}$ and $\widetilde{v}$ are also disjoint and by Theorem 2.1 again, they are separated by $c$,
which means that with the aid of a parametrization $c=c(\tau), \tau \in[0,1]$, all points $Q_{j}$ in $\#\{c, \tilde{u}\}$ can be expressed as $Q_{j}=c\left(\tau_{j}^{2}\right)$ and all points $P_{i}$ in $\#\{c, \tilde{v}\}$ can be expressed as $P_{i}=c\left(\tau_{j}^{1}\right)$, where $\tau_{i}^{1}<\tau_{j}^{2}$ for all $1 \leq i \leq r$ and $1 \leq j \leq k$. Without loss of generality, we assume $\tau_{1}^{2}=\min \left\{\tau_{1}^{2}, \ldots, \tau_{k}^{2}\right\}$ and $\tau_{k}^{2}=\max \left\{\tau_{1}^{2}, \ldots, \tau_{k}^{2}\right\}$. Then $c(\tau), \tau_{1}^{2}<\tau<\tau_{k}^{2}$, is a path in $c$ connecting $Q_{1}=c\left(\tau_{1}^{2}\right)$ and $Q_{k}=c\left(\tau_{k}^{2}\right)$, which is denoted by $c_{0}$. We see that $c_{1}=c_{0} \cdot \beta$ is homotopic to $\delta_{1}$.

We claim that $\widetilde{v}$ is disjoint from $\delta_{1}$. Since $\widetilde{v} \in W(\widetilde{u}, c)$, it is clear that $c_{0}$ does not meet $\widetilde{v}$. But $\widetilde{v}$ is disjoint from $\widetilde{u}$ and $\beta$ is a part of $\widetilde{u}$. We see that $\widetilde{v}$ does not meet $\beta$ either, which says $\widetilde{v}$ does not meet $c_{1}$ (because $c_{1}$ is the union of $\beta$ and $c_{0}$ ). This further implies that $\widetilde{v}$ is disjoint from $\delta_{1}$, contradicting that $\delta_{1} \in \mathscr{S}$. This proves Theorem 2.2.

## 5. Proof of Theorem 2.3

As usual, we let $c \in \mathscr{S}$ be determined by an element $f \in \mathscr{F}$. Let $\widetilde{u} \in$ $\mathcal{C}_{0}(\widetilde{S})$ be given an orientation. With respect to $\tilde{u}$, the geodesic $c$ can be expressed as a concatenation $c_{1} \cdot c_{2}$, where $c_{1}$ is freely homotopic to a geodesic $\delta_{1}$. Moreover, $c_{1}$ can further be written as the concatenation $c_{11}, \ldots, c_{1, k-1}$ of curves, where each $c_{1 i}$ is freely homotopic to nontrivial geodesics $\delta_{1 i}$. Set $\sum=\left\{\delta_{11}, \ldots, \delta_{1, k-1}\right\}$.

We prove Theorem 2.3 by establishing the following result:
Theorem 5.1. If $\Sigma$ contains two filling close geodesics $\delta_{1 i}, \delta_{1 j}$ for $1 \leq i, j \leq k-1$ and $|i-j| \geq 2$, then $d_{\mathcal{C}}(u, f(u)) \geq 4$ for those $u=\chi(\Omega)$ $\in F_{\widetilde{u}}$ for which $\Omega \cap \operatorname{axis}(g) \neq \varnothing$.

Proof. Let $\left(\Delta, \Delta^{\prime}\right)$ be the pair of maximal elements of $\mathscr{U}_{u}$ such that $\Delta$
and $\Delta^{\prime}$ cover the repelling and attracting fixed points of $g$, respectively. We know that $\Omega \subset \mathbf{H} \backslash\left(\Delta \cup \Delta^{\prime}\right)$. If $d_{\mathcal{C}}(u, f(u))=1$, then by Lemma 4.3, $i(\widetilde{u}, c)=1$, contradicting that $c=\delta_{1} \in \mathscr{S}^{*}$. If $d_{\mathcal{C}}(u, f(u))=2$, then by Theorem 2.1, for any $v \in V(u, f)$, we have $v \notin \mathcal{C}_{0}^{\prime}(S)$ and $\varepsilon(v) \in W(\widetilde{u}, c)$. This implies that $\varepsilon(v)$ does not intersect $\delta_{1}$, which says that $\delta_{1} \notin \mathscr{S}$. So $\delta_{1} \notin \mathscr{S}^{*}$.

Assume that $d_{\mathcal{C}}(u, f(u))=3$, where we write $u=\chi(\Omega)$ for an $\Omega \in \mathscr{R}_{\widetilde{u}}$ satisfying $\Omega \cap \operatorname{axis}(g) \neq \varnothing . \Omega$ Let $\left[u, v_{1}, v_{2}, f(u)\right]$ be a geodesic path in $\mathcal{C}_{0}(S)$ connecting $u$ and $f(u)$. Note that $v_{1}, v_{2} \in \mathcal{C}_{0}^{\prime}(S) \cup \hat{\mathcal{C}}_{0}(S)$.

Consider the case in which $v_{i} \in \hat{\mathcal{C}}_{0}(S), i=1,2$. Denote by $\left(\tau_{i}, \Omega_{i}, \mathscr{U}_{i}\right)$ the configuration corresponding to $v_{i}$. Let $\Delta_{i} \in \mathscr{U}_{i}$ be the maximal element that contains axis $(g)$ if $\Omega_{i} \cap \operatorname{axis}(g)=\varnothing$ (Figure 2(a) with $\left.\Delta_{v}=\Delta_{i}\right)$; and let $\Delta_{i}, \Delta_{i}^{\prime} \in \mathscr{U}_{i}$ be the pair of maximal elements so that $\Delta_{i}$ covers the repelling fixed point $B$ and $\Delta_{i}^{\prime}$ covers the attracting fixed point $A$ of $g$ if $\Omega_{i} \cap \operatorname{axis}(g) \neq \varnothing$ (Figure 2(b) with $\Delta=\Delta_{i}$ and $\left.\Delta^{\prime}=\Delta_{i}^{\prime}\right)$.

Since $\widetilde{v}_{i}$ and $\widetilde{u}$ are simple closed geodesics, all geodesics in $\left\{\varrho^{-1}\left(\widetilde{v}_{i}\right)\right\}$ are mutually disjoint, and so are the geodesics in $\left\{\varrho^{-1}(\widetilde{u})\right\}$. In particular, since $\left\{\partial \Delta, \partial \Delta^{\prime}, \partial g(\Delta)\right\} \subset\left\{\varrho^{-1}(\widetilde{u})\right\}, \quad \partial \Delta, \quad \partial \Delta^{\prime}$ and $\partial g(\Delta)$ are mutually disjoint. If, in addition, $v_{i}$ and $u$ are disjoint, then $\left\{\varrho^{-1}\left(\widetilde{v}_{i}\right)\right\}$ is also disjoint from $\left\{\varrho^{-1}(\widetilde{u})\right\}$ which particularly implies that $\left\{\varrho^{-1}\left(\widetilde{v}_{i}\right)\right\}$ is disjoint from $\left\{\partial \Delta, \partial \Delta^{\prime}, \partial g(\Delta)\right\}$. We will use these properties in each case below.

There are several possibilities: (1) $v_{1}, v_{2} \in \mathcal{C}_{0}^{\prime}(S)$, (2) $v_{1} \in \mathcal{C}_{0}^{\prime}(S)$ and $v_{2} \in \hat{\mathcal{C}}_{0}(S)$, (3) $v_{1} \in \hat{\mathcal{C}}_{0}(S)$ and $v_{2} \in \mathcal{C}_{0}^{\prime}(S)$, and (4) $v_{1}, v_{2} \in \hat{\mathcal{C}}_{0}(S)$.

Case 1. $v_{1}, v_{2} \in \mathcal{C}_{0}^{\prime}(S)$. We claim that this case cannot occur, as any two twice punctured disks enclosing $x$ have an overlap near $x$, as it turns out, $v_{1}, v_{2}$ must intersect, contradicting that $\left[v_{1}, v_{2}\right]$ is an edge in $\mathcal{C}_{1}(S)$.

Case 2. $v_{1} \in \mathcal{C}_{0}^{\prime}(S)$ and $v_{2} \in \hat{\mathcal{C}}_{0}(S)$. By Lemmas 5.1 and 5.2 of [12], $v_{1}$ corresponds to a parabolic fixed point $z_{1}$ of $G$, where $z_{1} \in \mathbf{S}^{1}$.

If $\Omega_{2} \cap \operatorname{axis}(g)=\varnothing$, then since $d_{\mathcal{C}}\left(v_{2}, f(u)\right)=1$, either $\Delta_{2} \subset g(\Delta)$ or $g(\Delta) \subset \Delta_{2}$. The former cannot occur, as $\Delta_{2}$ covers the point $A$ while $g(\Delta)$ does not. Hence we have $g(\Delta) \subset \Delta_{2}$. Notice that $v_{1}$ is disjoint from $v_{2}$. By Lemma 4.6, $z_{1} \in \Omega_{2} \cap \mathbf{S}^{1}$. This implies that $z_{1} \in\left(\mathbf{H} \backslash \Delta_{2}\right) \cap \mathbf{S}^{1}$. Hence $z_{1} \in(\mathbf{H} \backslash g(\Delta)) \cap \mathbf{S}^{1}$. But $\mathbf{H} \backslash g(\Delta) \subset \Delta^{\prime}$. Therefore, $z_{1}$ is covered by $\Delta^{\prime}$. From Lemma 4.6, $d_{\mathcal{C}}\left(v_{1}, u\right) \geq 2$, contradicting that $d_{\mathcal{C}}\left(v_{1}, u\right)=1$.

If $\Omega_{2} \cap \operatorname{axis}(g) \neq \varnothing$, then $\Omega_{2} \subset \mathbf{H} \backslash\left(\Delta_{2} \cup \Delta_{2}^{\prime}\right)$. By the assumption, $d_{\mathcal{C}}\left(v_{2}, f(u)\right)=1$. Hence either $g(\Delta) \subset \Delta_{2}$ or $\Delta_{2} \subset g(\Delta)$. Suppose that the former occurs. Then $d_{\mathcal{C}}\left(v_{1}, v_{2}\right)=1$ implies (from Lemma 4.6) that $z_{1} \in$ $\Omega_{2} \cap \mathbf{S}^{1}$. We deduce that $z_{1}$ is covered by $\Delta^{\prime}$. By Lemma 4.6 again, $d_{\mathcal{C}}\left(v_{1}, u\right) \geq 2$. This leads to a contradiction.

Suppose now $\Delta_{2} \subset g(\Delta)$. By the assumption, $\delta_{1} \in \mathscr{S}^{*}$. This means that $\left\{\varrho^{-1}\left(\tilde{v}_{2}\right)\right\}$ intersects $\hat{c}_{0}$ at least twice (where and below $\hat{c}_{0}$ denotes the geodesic segment $\operatorname{axis}(g) \cap g(\Delta) \cap \Delta^{\prime}$ ); that is to say, there are geodesics $\left[X_{2} Y_{2}\right],\left[X_{2}^{\prime} Y_{2}^{\prime}\right] \in\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ that are entirely contained in $g(\Delta) \cap \Delta^{\prime}$, where [ $\left.X_{2}^{\prime} Y_{2}^{\prime}\right]$ is the one closest to $[M N]=\partial g(\Delta)$. If $\partial \Delta_{2}=[M N]$ or $\left[X_{2}^{\prime}, Y_{2}^{\prime}\right]$, then since $\Omega_{2} \subset \mathbf{H} \backslash\left(\Delta_{2} \cup \Delta_{2}^{\prime}\right)$ and since $d_{\mathcal{C}}\left(v_{1}, v_{2}\right)=1$, by Lemma 4.6, $z_{1} \in \Omega_{2} \cap \mathbf{S}^{1}$, and hence $z_{1} \in \Delta^{\prime}$. By Lemma 4.6 again, $d_{\mathcal{C}}\left(u, v_{1}\right) \geq 2$, contradicting the hypothesis. If $\partial \Delta_{2}=\left[X_{2} Y_{2}\right]$ or any other geodesics of $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ intersecting $\hat{c}_{0}$, then by Lemma 4.6, $\Delta_{2}^{\prime}$ must contain [ $X_{2}^{\prime} Y_{2}^{\prime}$ ]
and so $\Delta_{2}^{\prime} \cap g(\Delta) \neq \varnothing$. It follows that $d_{\mathcal{C}}\left(f(u), v_{2}\right) \geq 2$. This is again a contradiction.

Case 3. $v_{1} \in \hat{\mathcal{C}}_{0}(S)$ and $v_{2} \in \mathcal{C}_{0}^{\prime}(S)$. This time, $v_{2}$ corresponds to a parabolic fixed point $z_{2}$ of $G$. By the assumption, $\widetilde{v}_{1}$ is nontrivial and $u, v_{1}$ are disjoint.

Suppose that $\Omega_{1} \cap \operatorname{axis}(g)=\varnothing$. By the assumption, $d_{\mathcal{C}}\left(u, v_{1}\right)=1$. Hence either $\Delta_{1} \subset \Delta$ or $\Delta \subset \Delta_{1}$. The former case does not occur, as $\Delta_{1}$ covers the point $A$ while $\Delta$ does not. So we must have $\Delta \subset \Delta_{1}$; or $\mathbf{H} \backslash \Delta_{1} \subset \mathbf{H} \backslash \Delta$. The assumption says $d_{\mathcal{C}}\left(v_{1}, v_{2}\right)=1$. By Lemma 4.6, $z_{2} \in$ $\Omega_{1} \cap \mathbf{S}^{1}$. But $\Omega_{1} \subset \mathbf{H} \backslash\left(\Delta_{1} \cup \Delta_{1}^{\prime}\right)$. So $z_{2}$ is covered by $g(\Delta)$. It follows that $d_{\mathcal{C}}\left(f(u), v_{2}\right) \geq 2$.

Suppose that $\Omega_{1} \cap \operatorname{axis}(g) \neq \varnothing$. In this case, $\Omega_{1} \subset \mathbf{H} \backslash\left(\Delta_{1} \cup \Delta_{1}^{\prime}\right)$. By the assumption, $d_{\mathcal{C}}\left(u, v_{1}\right)=1$. This implies that $\Delta \subset \Delta_{1}$ or $\Delta_{1} \subset \Delta$. Assume that $\Delta_{1} \subset \Delta$. Since $\left\{\varrho^{-1}\left(\widetilde{v}_{1}\right)\right\}$ intersects $\hat{c}_{0}$ at least twice, let $\left[X_{1} Y_{1}\right],\left[X_{1}^{\prime} Y_{1}^{\prime}\right]$ be such geodesics, where $\left[X_{1}^{\prime} Y_{1}^{\prime}\right]$ is the one in $\left\{\varrho^{-1}\left(\widetilde{v}_{1}\right)\right\}$ closest to $[M N]$ $=\partial g(\Delta)$. Note that $\left[X_{1} Y_{1}\right],\left[X_{1}^{\prime} Y_{1}^{\prime}\right]$ are contained in $\Delta^{\prime} \cap g(\Delta)$. It follows that $\Delta_{1}^{\prime}$ contains [ $X_{1}^{\prime} Y_{1}^{\prime}$ ]. But $z_{2} \in \Omega_{1} \cap \mathbf{S}^{1}$ and $\Omega_{1} \subset \mathbf{H} \backslash\left(\Delta_{1} \cup \Delta_{1}^{\prime}\right)$. We conclude that $z_{2}$ is covered by $g(\Delta)$. So $d_{\mathcal{C}}\left(v_{2}, f(u)\right) \geq 2$.

Assume now that $\Delta \subset \Delta_{1}$. The assumption leads to that $\left\{\varrho^{-1}\left(\widetilde{v}_{1}\right)\right\}$ intersects $\hat{c}_{0}$ at least twice. Let $\left[X_{1} Y_{1}\right],\left[X_{1}^{\prime} Y_{1}^{\prime}\right]$ be as above. Then since $d_{\mathcal{C}}\left(u, v_{1}\right)=1, \Delta_{1} \subset \mathbf{H} \backslash \Delta^{\prime}$. Hence by Lemma 4.6, $\Delta_{1}^{\prime}$ contains [ $X_{1} Y_{1}$ ]. It follows that $z_{2} \in \Omega_{1} \cap \mathbf{S}^{1}$ must be covered by $g(\Delta)$. This contradicts that $d_{\mathcal{C}}\left(v_{2}, f(u)\right)=1$.

Case 4. $v_{1}, v_{2} \in \hat{\mathcal{C}}_{0}(S)$. Then for $i=1,2, v_{i}$ corresponds to the configurations $\left(\tau_{i}, \Omega_{i}, \mathscr{U}_{i}\right)$. By the assumption, $v_{1}$ and $v_{2}$ are disjoint.

Suppose that $\Omega_{1} \cap \operatorname{axis}(g)=\varnothing$ and $\Omega_{2} \cap \operatorname{axis}(g)=\varnothing$. By the same argument as in Case $3, \Delta \subset \Delta_{1}$. By Lemma 4.5, $d_{\mathcal{C}}\left(u, v_{1}\right)=1$ implies that $\Omega_{1} \cap \Omega \neq \varnothing$. Hence $\mathbf{H} \backslash \Delta_{1} \subset \mathbf{H} \backslash\left(\Delta \cup \Delta^{\prime}\right)$. If $\Delta_{1} \subset \Delta_{2}$, i.e., $\mathbf{H} \backslash \Delta_{2} \subset \mathbf{H} \backslash \Delta_{1}$, then $\mathbf{H} \backslash \Delta_{2} \subset g(\Delta)$. By Lemma 4.5, $d_{\mathcal{C}}\left(f(u), v_{2}\right) \geq 2$, contradicting that $d_{\mathcal{C}}\left(f(u), v_{2}\right)=1$. Suppose $\Delta_{2} \subset \Delta_{1}$, i.e., $\mathbf{H} \backslash \Delta_{1} \subset \mathbf{H} \backslash \Delta_{2}$. See Figure 3(a). Note that $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ intersects $\hat{c}_{0}$ more than once, and all geodesics in $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ are disjoint, as a member of $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}, \partial \Delta_{2}$ must be disjoint from any geodesic in $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ intersecting $\hat{c}_{0}$. It follows that $\mathbf{H} \backslash \Delta_{2} \subset g(\Delta)$. By Lemma 4.5, $d_{\mathcal{C}}\left(f(u), v_{2}\right) \geq 2$, contradicting that $d_{\mathcal{C}}\left(f(u), v_{2}\right)=1$.

If $\Omega_{1} \cap \operatorname{axis}(g)=\varnothing$ and $\Omega_{2} \cap \operatorname{axis}(g) \neq \varnothing$, then by the discussion above, we know that $\Delta \subset \Delta_{1}$ and $\mathbf{H} \backslash \Delta_{1} \subset \mathbf{H} \backslash\left(\Delta \cup \Delta^{\prime}\right)$. By the assumption, there are at least two geodesics in $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ intersecting $\hat{c}_{0}$. Let $\left[X_{2} Y_{2}\right]$ be the one which is closest to $[E F]=\partial \Delta^{\prime}$. See Figure 3(b).


Figure 3
Recall that $\left(\Delta_{2}, \Delta_{2}^{\prime}\right)$ is the pair of the maximal elements of $\mathscr{U}_{2}$ that crosses axis $(g)$. We claim that $\Delta_{2}^{\prime}$ contains $\left[X_{2} Y_{2}\right]$. Otherwise, by Lemma 2.1 of [14], $\Delta_{2}$ contains [ $X_{2} Y_{2}$ ]. Since $\Omega_{2} \subset \mathbf{H} \backslash\left(\Delta_{2} \cup \Delta_{2}^{\prime}\right)$ and $\Omega_{1} \subset \mathbf{H} \backslash \Delta_{1}$,
we see that $\bar{\Omega}_{2} \cap \bar{\Omega}_{1}=\varnothing$. By Lemma 4.4, $d_{\mathcal{C}}\left(v_{1}, v_{2}\right) \geq 2$, contradicting that $d_{\mathcal{C}}\left(v_{1}, v_{2}\right)=1$. We conclude that $\Delta_{2}^{\prime}$ contains $\left[X_{2} Y_{2}\right]$. But then $\Delta_{2}^{\prime} \cap$ $g(\Delta) \neq \varnothing$. So again, by Lemma 4.5, $d_{\mathcal{C}}\left(v_{1}, v_{2}\right) \geq 2$, which also leads to a contradiction.

The case where $\Omega_{1} \cap \operatorname{axis}(g) \neq \varnothing$ and $\Omega_{1} \cap \operatorname{axis}(g)=\varnothing$ can be handled similarly.

If $\Omega_{1} \cap \operatorname{axis}(g) \neq \varnothing$ and $\Omega_{2} \cap \operatorname{axis}(g) \neq \varnothing$, we note that $\left\{\varrho^{-1}\left(\widetilde{v}_{1}\right)\right\}$, $\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}$ and $\left\{\varrho^{-1}(\widetilde{u})\right\}$ are mutually disjoint. From the assumption, $\left\{\varrho^{-1}(\widetilde{u})\right\}$ intersects $\hat{c}_{0}$ at points $q_{1}, \ldots, q_{k}$, where $q_{1}=[E F] \cap \hat{c}_{0}$ and $q_{k}=[M N] \cap \hat{c}_{0}$. See Figure 4.

Note that $c_{1 i}$ and $c_{1 j}$ correspond to two intervals $E_{i}$ and $E_{j}$ on $\hat{c}_{0}$. Since $|i-j| \geq 2$, there is a geodesic [CD] that separates $E_{i}$ and $E_{j}$. Denote $q=[C D] \cap \hat{c}_{0}$. Then the geodesic segment $\hat{c}_{0}$, which is a portion of $\operatorname{axis}(g)$, can be written as $\left[q_{1} q\right] \cup\left[q q_{k}\right]$. Thus, $E_{i} \subset\left[q_{1} q\right]$ and $E_{j} \subset\left[q q_{k}\right]$. By the assumption, $\delta_{1 i}$ and $\delta_{1 j} \in \mathscr{S}$ for some $i, j$ with $j-i \geq 2$. This means that $\widetilde{v}_{1}$ and $\widetilde{v}_{2}$ intersect $\delta_{1 i}$ and $\delta_{1 j}$. Therefore, the segments [ $q_{1} q$ ] and $\left[q q_{k}\right]$ both intersect geodesics in $\Sigma_{1}=\left\{\hat{v} \in\left\{\varrho^{-1}\left(\widetilde{v}_{1}\right)\right\}: \hat{v} \subset g(\Delta) \cap \Delta^{\prime}\right\}$. Similarly, $\left[q_{1} q\right]$ and $\left[q q_{k}\right]$ also intersect geodesics in $\Sigma_{2}=\left\{\hat{v} \in\left\{\varrho^{-1}\left(\widetilde{v}_{2}\right)\right\}\right.$ : $\left.\hat{v} \subset g(\Delta) \cap \Delta^{\prime}\right\}$.

For simplicity, let $\left[X_{1}^{\prime} Y_{1}^{\prime}\right] \in \Sigma_{1}$ denote the geodesic that is closest to the geodesic $[M N]$, and let $\left[X_{1} Y_{1}\right] \in \Sigma_{1}$ denote the geodesic that is closest to the geodesic $[E F]$. Then $\left[X_{1}^{\prime} Y_{1}^{\prime}\right]$ lies in the component of $\mathbf{H}$ bounded by [ $M N$ ] and [CD], and $\left[X_{1} Y_{1}\right]$ lies in the component of $\mathbf{H}$ bounded by $[E F]$ and $[C D]$. See Figure 4(a). Likewise, let $\left[X_{2}^{\prime} Y_{2}^{\prime}\right] \in \Sigma_{2}$ be the geodesic closest to $[M N]$, and $\left[X_{2} Y_{2}\right] \in \Sigma_{2}$ be the geodesic closest to $[E F]$. It is also
obvious that $\left[X_{2}^{\prime} Y_{2}^{\prime}\right]$ lies in the component of $\mathbf{H}$ bounded by [MN] and $[C D]$, and $\left[X_{2} Y_{2}\right]$ lies in the component of $\mathbf{H}$ bounded by $[E F]$ and $[C D]$. See Figure 4(b). Notice that geodesics in $\Sigma_{1}$ and $\Sigma_{2}$ are boundaries of elements of $\mathscr{U}_{1}$ and $\mathscr{U}_{2}$, respectively. Consider the pairs $\left(\Delta_{1}, \Delta_{1}^{\prime}\right)$ and $\left(\Delta_{2}, \Delta_{2}^{\prime}\right)$ of maximal elements of $\mathscr{U}_{1}$ and $\mathscr{U}_{2}$, respectively. Recall that $\Delta_{1}^{\prime}$ and $\Delta_{2}^{\prime}$ cover the attracting fixed point $A$ of $g$, while $\Delta_{1}$ and $\Delta_{2}$ cover the repelling fixed point $B$ of $g$.

We claim that $\partial \Delta_{1}^{\prime}=\left[X_{1} Y_{1}\right]$ or $\partial \Delta_{1}^{\prime}$ is disjoint from $\Delta^{\prime}$. If not, then by Lemma 2.1 of [14], either $\partial \Delta_{1}=\left[X_{1} Y_{1}\right]$ or $\left[X_{1} Y_{1}\right]$ is contained in $\Delta_{1}$. Since $\Omega_{1} \subset \mathbf{H} \backslash\left(\Delta_{1} \cup \Delta_{1}^{\prime}\right)$, we see that $\Omega_{1}$ is disjoint from $\Omega$. From Lemma 4.5, $d_{\mathcal{C}}\left(u, v_{1}\right) \geq 2$, contradicting that $d_{\mathcal{C}}\left(u, v_{1}\right)=1$. So $\Delta_{1}^{\prime}$ contains [CD].

We also claim that $\partial \Delta_{2}=\left[X_{2}^{\prime} Y_{2}^{\prime}\right]$ or $\partial \Delta_{2}$ is disjoint from $g(\Delta)$. Suppose not. By Lemma 2.1 of [14] again, either $\partial \Delta_{2}^{\prime}=\left[X_{2}^{\prime} Y_{2}^{\prime}\right]$ or $\left[X_{2}^{\prime} Y_{2}^{\prime}\right]$ is contained in $\Delta_{2}^{\prime}$. But $\Omega_{2} \subset \mathbf{H} \backslash\left(\Delta_{2} \cup \Delta_{2}^{\prime}\right)$. Hence $\Omega_{2}$ is disjoint from $g(\Omega)$. Clearly, $\chi(g(\Omega))=g^{*}(\chi(\Omega))=f(u)$. From Lemma 4.5, $d_{\mathcal{C}}\left(f(u), v_{2}\right)$ $\geq 2$, contradicting that $d_{\mathcal{C}}\left(f(u), v_{2}\right)=1$. So $\Delta_{2}$ contains [CD].

(a)

(b)

Figure 4

In particular, $\Delta_{1}^{\prime} \cap \Delta_{2} \neq \varnothing$. As members of $\varrho^{-1}\left(\widetilde{v}_{1}\right)$ and $\varrho^{-1}\left(\widetilde{v}_{1}\right), \partial \Delta_{1}^{\prime}$ is disjoint from $\partial \Delta_{2}$. Since $\Omega_{1} \subset \mathbf{H} \backslash\left(\Delta_{1} \cup \Delta_{1}^{\prime}\right)$ and $\Omega_{2} \subset \mathbf{H} \backslash\left(\Delta_{2} \cup \Delta_{2}^{\prime}\right)$, we conclude that $\bar{\Omega}_{1} \cap \bar{\Omega}_{2}=\varnothing$. By Lemma 4.5 , we see that $d_{\mathcal{C}}\left(v_{1}, v_{2}\right) \geq 2$, contradicting that $d_{\mathcal{C}}\left(v_{1}, v_{2}\right)=1$.

## Acknowledgment

The author is grateful to the referees for their insights, careful reading of this paper and for their helpful comments and suggestions.

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