# A PRIORI BOUNDS FOR SOME INFINITELY RENORMALIZABLE QUADRATICS: I. BOUNDED PRIMITIVE COMBINATORICS 

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#### Abstract

We prove the a priori bounds for infinitely renormalizable quadratic polynomials for which we can find an infinite sequence of primitive renormalizations for which the ratios of the periods of successive renormalizations is bounded.


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Date: January 24, 2023.
This work was supported in part by Sloan Research Fellowship and NSF grants DMS-8920768 and DMS-9022140.
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## 1. Introduction

1.1. Statement of the main result. In this paper, we will prove a priori bounds for infinitely renormalizable quadratic polynomials of bounded primitive type. Let us begin with recalling some basic definitions of the holomorphic renormalization theory. Most of them can be found in any source on the subject, see [L3, L2, McM1].

A quadratic polynomial $f: z \mapsto z^{2}+c$ is called primitively renormalizable with period $p$ if there exist topological disks $V \supseteq U \ni 0$ such that $f^{p}: U \rightarrow V$ is a quadratic-like map with connected Julia set, and the domains $f^{n} U, n=0,1, \ldots, p-1$ are pairwise disjoint. This quadraticlike map is called a renormalization of $f$. If there is an infinite sequence of periods $p_{0}<p_{1}<\ldots$ such that $f$ is primitively $p_{k}$-renormalizable then $f$ is called infinitely primitively renormalizable. If additionally, there exists a $B$ such that $p_{k+1} / p_{k} \leqslant B, k=0,1, \ldots$, then $f$ is called infinitely primitively renormalizable of bounded type. Such a map has $a$ priori bounds if there exists an $\varepsilon>0$ and a sequence of quadratic-like renormalizations $f^{p_{k}}: U_{k} \rightarrow V_{k}$ such that $\bmod \left(V_{k} \backslash U_{k}\right) \geqslant \varepsilon$.

Main Theorem. Let $f$ be an infinitely primitively renormalizable quadratic polynomial of bounded type. Then $f$ has a priori bounds.

In the forthcoming notes [KL2], we will prove a priori bounds for a class of infinitely renormalizable maps of unbounded type.

For real quadratics of bounded type a priori bounds were proved by Sullivan [S], see also [LS, LY, MS]. They were also proved for a class of complex combinatorics of "high bounded type" [L1].
1.2. Consequences. The a priori bounds have numerous consequences. Let us list some of them (below $f_{c}: c \mapsto z^{2}+c$ stands for an infinitely primitively renormalizable quadratic polynomial of bounded type):

- The Mandelbrot set is locally connected at c, or equivalently, the polynomial $f_{c}$ is combinatorially rigid (see [L1]). The conjecture of local connectivity of the Mandelbrot set (the MLC Conjecture) formulated about 20 years ago by Douady and Hubbard (see [DH1]) is a central open problem in holomorphic dynamics. Previously, it was established for all quadratic maps which are not infinitely renormalizable (Yoccoz, see $[\mathrm{H}]$ ) and for the class of infinitely renormalizable maps of high type mentioned above (see [L1]).
- The Julia set $J\left(f_{c}\right)$ is locally connected (see [HJ, J]).
- The Feigenbaum-Coullet-Tresser Renormalization Conjecture is valid for primitive combinatorics. This conjecture was established in the work of Sullivan [S], McMullen [McM2] and Lyubich [L2] assuming a priori bounds (and thus, unconditionally, for real maps). Now, these results become unconditional for arbitrary primitive combinatorics.
- Universality and Hairiness of the Mandelbrot set at c. These properties were conjectured by Milnor [M] and proved in [L2] for maps with a priori bounds.
- The Basic Trichotomy for the measure and Hausdorff dimension of the Julia set $J\left(f_{c}\right)$ which was established in [AL] for maps with a priori bounds.
1.3. Outline of the proof. We will now give a brief top-down outline of the proof of the Main Theorem.

General strategy: improving of the lengths of the hyperbolic geodesics. Let $K_{n}$ be the filled Julia set of the renormalization $f^{p_{n}}: U_{n} \rightarrow V_{n}$,

$$
\mathcal{K}_{n}=\bigcup_{i=0}^{p_{n}-1} f^{i}\left(K_{n}\right),
$$

and let $\gamma_{n}$ be the peripheral hyperbolic geodesic in $\mathbb{C} \backslash \mathcal{K}_{n}$ going around $K_{n}$. The a priori bounds are equivalent to the assertion that the hyperbolic lengths of the $\gamma_{n}$ are bounded. Our strategy towards this end is to show that if the length of some $\gamma_{n}$ gets long then it was even longer before: There exist $M>0$ and $l>0$ such that

$$
\begin{equation*}
\left|\gamma_{n}\right|>M \Rightarrow\left|\gamma_{n-l}\right|>2 M \tag{1.1}
\end{equation*}
$$

Localization "in principle". The immediate and obvious difficulty in proving (1.1) is that the number of components of $\mathcal{K}_{n}$ grows exponentially with $n$, and of course is unbounded as $n$ goes to infinity. While there has been some success in handling an unbounded number of components, there is certainly a vast advantage to working with a bounded number, and this paper is able to make a major step forward with the help of such a localization.

The central idea is that, while $n$ is unbounded, $l$ is bounded, and, because we assume $p_{k+1} / p_{k} \leqslant B$ (for all $k$ ), the number of components of $\mathcal{K}_{n}$ that are surrounded by $\gamma_{n-l}$ is at most $B^{l}$. Moreover, if we let $f_{n-l} \equiv f^{p_{n-l}}: U_{n-l} \rightarrow V_{n-l}$ be a renormalization of period $p_{n-l}$ of $f$, then this $f_{n-l}$ is renormalizable of period $p_{n} / p_{n-l}$, and the small Julia sets for this renormalization are exactly the components of $\mathcal{K}_{n}$ surrounded by $\gamma_{n-l}$.

With this in mind, we will first state and prove a theorem about a primitively renormalizable quadratic-like map that will be closely related to (1.1). Before we do so, we will introduce some notation. Suppose that $f: U \rightarrow V$ is a quadratic-like map that is primitively $p$ renormalizable. We let $K \equiv K(f)$ be the filled-Julia set for $f$. We let $\mathcal{K}_{p} \equiv \mathcal{K}_{p}(f)$ be the union of small filled-Julia sets for this renormalization, and $K_{p}$ be the central small filled-Julia set. We let $\gamma$ be the closed geodesic in $V \backslash K$ that goes around $K$, and let $\gamma_{p}$ be the hyperbolic geodesic in $V \backslash \mathcal{K}_{p}$ that goes around $K_{p}$.
Theorem 1.2. There exists $\epsilon>0$, such that for all $p$ there exists $M$ : Suppose that $f$ is a quadratic-like map that is primitively p-renormalizable. Then, with the notation described just above,

$$
\begin{equation*}
l\left(\gamma_{p}\right)>M \Longrightarrow l(\gamma) \geqslant \epsilon p l\left(\gamma_{p}\right) . \tag{1.3}
\end{equation*}
$$

Now suppose $f$ in Theorem 1.2 is the renormalization $f_{n-l}$ described above, and $V_{n-l}$ is taken to be "pretty large", and let $q=p_{n} / p_{n-k}$. Then the $\gamma_{q}$ for the $f$ in Theorem 1.2 is a "good approximation" of the $\gamma_{n}$ for the original $f$, and the $\gamma$ for the $f$ in Theorem 1.2 is a "good approximation" of the $\gamma_{n-l}$ for the original $f$. We can choose $q$ large but bounded, and (1.1) would follow from (1.3) if these "good approximations" were actually equalities.

To make them equalities we introduce (in Section 8) the pseudo-quadratic-like map and the canonical renormalization. A pseudo-quadraticlike map $(i, f): \mathbf{U} \rightarrow \mathbf{V}$ is a generalization of a quadratic-like map where the inclusion from $\mathbf{U}$ to $\mathbf{V}$ is replaced by an immersion $i$ that has certain properties that imply the (multivalued) $f \circ i^{-1}$ has a (singlevalued) quadratic-like restriction. We can then, given our original quadratic polynomial $f$, define, for each $n$, the canonical renormalization $f_{n}: \mathbf{U}_{n} \rightarrow \mathbf{V}_{n}$ (really $\left(i_{n}, f_{n}\right)$ but we will suppress the $i_{n}$ here), which will be a pseudo-quadratic-like map. Then, given $n$ and $l<n$, and letting $q=p_{n} / p_{n-l}$, we'll see that $f_{n-l}$ has two crucial properties:
(1) The geodesic $\gamma$ around $K\left(f_{n-l}\right)$ in $\mathbf{V}_{n-l} \backslash K\left(f_{n-l}\right)$ has length equal to $\gamma_{n-l}$.
(2) $f_{n-l}$ is renormalizable of period $q$, and the geodesic $\gamma_{q}$ (for $f_{n-l}$ ) around $K_{q}\left(f_{n-l}\right)$ in $\mathbf{V}_{n-l} \backslash \mathcal{K}_{q}\left(f_{n-l}\right)$ has length greater than or equal to the length of $\gamma_{n}(f)$.
In Section 9, we then extend Theorem 1.2 to pseudo-quadratic-like maps, and (1.1) and our Main Theorem then follow.

Canonical weighted arc diagrams. Let us now briefly discuss the proof of Theorem 1.2. In Section 2 we develop the theory of the canonical weighted arc-diagram, which describes the "thin rectangles" of large
modulus which we will then use to estimate $\left|\gamma_{p}\right|$ as the sum of extremal widths of thin rectangles in $V \backslash \mathcal{K}_{p}$ crossed by $\gamma_{p}$.

Roughly speaking, the canonical weighted-arc diagram is

$$
W_{\mathrm{can}}=\sum W(\alpha) \alpha
$$

where the $\alpha$ are $\operatorname{arcs}^{1}$ in $V \backslash \mathcal{K}_{p}$ represented by paths crossing the thin rectangles and the $W(\alpha)$ are their widths. Holomorphic maps between Riemann surfaces, such as covering and inclusions, then induce relations between their canonical WAD's; these relations are described in Section 2. The most subtle of these relations is domination, which describes the relation between the canonical WAD's between Riemann surfaces $U \subset V$, and encodes the parallel and series laws for extremal length.

In Section 3 we reduce Theorem 1.2 to a similar statement expressed in terms of the canonical WAD of $V \backslash \mathcal{K}_{p}$, and then provide an outline of how to use the functorial properties of the canonical WAD (and the Covering Lemma[KL1], which is applied in Section 7) to reduce that statement, in turn, to a purely combinatorial one, which is proven in Sections 5 and 6, using the background in Section 4.
1.4. Terminology and Notation. We let:
$\mathbb{N}=\{1,2, \ldots\}$ be the set of natural numbers; $\mathbb{Z}_{\geqslant 0}=\mathbb{N} \cup\{0\}$;
$\mathbb{D}=\{z:|z|<1\}$ be the unit disk, $\mathbb{T}_{r}=\{z:|z|=r\}, \mathbb{T}=\mathbb{T}_{1} ;$
$\mathbb{A}(r, R)=\{z: r<|z|<R\}$.
A topological disk means a simply connected domain in some Riemann surface $S$. We will say a subset $K$ of $\mathbb{C}$ is an FJ-set ("filled Julia set") if $K$ is compact, connected, and full.
1.5. Acknowledgment. This paper was written in collaboration with Mikhail Lyubich following the mathematical work of the author. The primary result for this paper will be combined with the main result of a subsequent paper of Mikhail Lyubich and the author to form one result that is the goal of the series of papers of which this paper is the first.

This work has been partially supported by the Clay Mathematics Institute and the Simons Foundation. Part of it was done during the author's visit to the IMS at Stony Brook and the Fields Institute in Toronto. The author is thankful to these Institutions and Foundations.

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## 2. CANONICAL WEIGHTED ARC DIAGRAM

2.1. Arc diagrams. Let $S$ be a hyperbolic open Riemann surface of finite topology without cusps. It is conformally equivalent to the interior of a compact Riemann surface $\mathbf{S}$ with non-empty boundary. ${ }^{2}$ The boundary of $\mathbf{S}$ is called the ideal boundary of $S$. It is canonically attached to $S$ in the sense that any conformal isomorphism $S \rightarrow S^{\prime}$ extends to a homeomorphism $\mathbf{S} \rightarrow \mathbf{S}^{\prime}$.

A path in some topological space $Z$ is an embedded interval $\gamma: I \rightarrow Z$. It is called open, closed, or semiclosed depending on the nature of $I$. An open path $\gamma:(0,1) \rightarrow S$ is called proper if it extends to a closed path $\gamma:[0,1] \rightarrow \mathbf{S}$ such that $\gamma\{0,1\} \subset \partial \mathbf{S}$. Two proper paths $\gamma_{0}$ and $\gamma_{1}$ in $S$ are called properly homotopic if there is a homotopy $\gamma_{t}, t \in[0,1]$, connecting $\gamma_{0}$ to $\gamma_{1}$ through a family of proper paths. ${ }^{3}$

An arc on $S$ is a class of properly homotopic paths, $\alpha=[\gamma]$. An arc is called trivial if it has representing paths $\gamma: I \rightarrow \mathrm{~S}$ in arbitrary small neighborhoods of $\partial \mathbf{S}$. Let $\mathcal{A}(S)$ stand for the set of all non-trivial arcs on $S$.

Two different arcs are said to be non-crossing if they can be represented by non-crossing paths. An arc diagram is a family of pairwise non-crossing $\operatorname{arcs} \alpha_{i}$. Note that any arc diagram consists of at most $3|\chi(S)|$ arcs, where $\chi(S)$ is the Euler characteristic of $S$.

A weighted arc diagram (WAD) on $S$ is an arc diagram $\mathbf{a}=\left\{\alpha_{i}\right\}$ endowed with weights $w_{i} \in \mathbb{R}_{+}$. In this case, the arc diagram a is called the support of $W .{ }^{4}$ Let $\mathcal{W}(S)$ stand for the set of WAD's on $S$.

The set $\mathcal{W}(S)$ is partially ordered: $X \leqslant Y$ is $X(\alpha) \leqslant Y(\alpha)$ for any $\alpha \in \operatorname{supp} X$. We will also write $X \leqslant Y+c$ if $X(\alpha) \leqslant Y(\alpha)+c$ for any $\alpha \in \operatorname{supp} X$.

The sum of two WAD's, $X+Y$, is well defined whenever any two $\operatorname{arcs} \alpha \in \operatorname{supp} X$ and $\beta \in \operatorname{supp} Y$ are either the same or non-crossing. The difference $X-Y$ is always well defined if we let $(X-Y)(\alpha)=0$ whenever $X(\alpha) \leqslant Y(\alpha)$. Similarly, $X-c$ is well defined for any constant $c \geqslant 0$.

We will make use of two norms on the space of WAD's:

$$
\|W\|_{\infty}=\sup _{\alpha \in \mathcal{A}} W(\alpha) ; \quad\|W\|_{1}=\sum_{\alpha \in \mathcal{A}} W(\alpha) .
$$

[^1]If $f: U \rightarrow V$ is a holomorphic covering between two Riemann surfaces then there is a natural pull-back operation $f^{*}: \mathcal{W}(V) \rightarrow \mathcal{W}(U)$ acting on the WAD's.

A proper lamination $\mathcal{F}$ on $S$ is a Borel set $Z \subset S$ explicitly realized as a union of disjoint proper paths called the leaves of $\mathcal{F}$. Any proper lamination ${ }^{5}$ can be written $\mathcal{F}=\bigcup_{\alpha} \mathcal{F}(\alpha)$, where $\mathcal{F}(\alpha)$ comprises the leaves of $\mathcal{F}$ that represent $\alpha$. The $\operatorname{arcs} \alpha \in \mathcal{A}$ for which $\mathcal{F}(\alpha)$ is nonempty assemble an arc diagram. Let us weight each of these arcs with the weight $W_{\mathcal{F}}(\alpha)$ equal to the extremal width $\mathcal{W}(\mathcal{F}(\alpha))$ of the sublamination $\mathcal{F}(\alpha)$ (viewed as a path family). In this way we obtain the WAD $W_{\mathcal{F}}=\sum_{\alpha} W_{\mathcal{F}}(\alpha) \cdot \alpha$ corresponding to $\mathcal{F}$.

Note that if $f: U \rightarrow V$ is a holomorphic covering between two Riemann surfaces and $\mathcal{F}$ is a proper lamination on $V$ then $f^{*}(\mathcal{F})$ is a proper lamination on $U$ and $W_{f *(\mathcal{F})}=f^{*}\left(W_{\mathcal{F}}\right)$.

Weighted arc diagrams that are $W_{\mathcal{F}}$ for some proper lamination $\mathcal{F}$ are called valid.
2.2. Canonical WAD. Let us consider the universal covering $\pi: \mathbb{D} \rightarrow$ int $\mathbf{S}$. Let $\Gamma$ be the Fuchsian group of deck transformations of $\pi$, and let $\Lambda \subset \mathbb{T}$ be its limit set. Since $\mathbf{S}$ has non-empty boundary, $\Lambda$ is a Cantor set. Moreover, $\pi$ extends continuously onto $\hat{S}=\overline{\mathbb{D}} \backslash \Lambda$, and the restriction of $\pi$ to any component $I$ of $\partial \hat{S}$ is a universal covering onto some component $I$ of $\partial \mathbf{S}$.

Let us pick two components, $I \neq J$, of $\partial \hat{S}$. The disk $\mathbb{D}$ with these two intervals as horizontal sides determines a quadrilateral $Q(I, J)$. This quadrilateral can be conformally uniformized, $\phi: Q(I, J) \rightarrow \mathbf{Q}(a)$, by a standard quadrilateral $\mathbf{Q}(a)=[0, a] \times[0,1]$ in such a way that $I$ and $J$ correspond to the horizontal sides of $\mathbf{Q}(a)$. The vertical foliation $\mathcal{F}(I, J)$ on $Q(I, J)$ is the $\phi$-pullback of the standard vertical foliation on $\mathbf{Q}(a)$.

Assume now that $a>2$, and let us cut off from $\mathbf{Q}(a)$ two side squares, $[0,1] \times[0,1]$ and $[a-1, a] \times[0,1]$. We call the left-over rectangle $\mathbf{Q}_{\text {can }}(a)$, and we let $Q_{\text {can }}(I, J)=\phi^{-1}\left(\mathbf{Q}_{\text {can }}(a)\right)$. The side quadrilaterals that we have cut off from $Q(I, J)$ are called its buffers.

Let $\mathcal{F}_{\text {can }}(I, J)$ be the restriction of $\mathcal{F}(I, J)$ to $Q_{\text {can }}(I, J)$. Obviously, for any deck transformation $\gamma \in \Gamma$, we have:

$$
\begin{equation*}
\mathcal{F}_{\text {can }}(\gamma(I), \gamma(J))=\gamma\left(\mathcal{F}_{\text {can }}(I, J)\right) \tag{2.1}
\end{equation*}
$$

Lemma 2.2. The rectangles $Q_{\text {can }}(I, J)$ are pairwise disjoint.

[^2]

Figure 2.1. The canonical foliation
Proof. Let us consider two rectangles, $Q \equiv Q(I, J)$ and $Q^{\prime} \equiv Q\left(I^{\prime}, J^{\prime}\right)$. Then we can find one interval from each pair, say $J$ and $J^{\prime}$, such that $J \neq J^{\prime}$ (so $J \cap J^{\prime}=\varnothing$ ), and such that there is a component $T$ of $\mathbb{T}-\left(J \cup J^{\prime}\right)$ such that $T \cap\left(I \cup I^{\prime}\right)=\varnothing$.

We let $B$ be the buffer of $Q(I, J)$ that has a horizontal side (which is a subset of $J$ ) that shares an endpoint with $T$; we define $B^{\prime}$ likewise. Then if any vertical leaf $\gamma$ of $\mathcal{F}_{\text {can }}(I, J)$ crossed any vertical leaf $\gamma^{\prime}$ of $\mathcal{F}_{\text {can }}\left(I^{\prime}, J^{\prime}\right)$, then every vertical leaf of $B$ would cross every vertical leaf of $B^{\prime}$. This would contradict Lemma 10.6.

This lemma allows us to define the canonical lamination $\mathcal{F}_{\text {can }}(\hat{S})$ as the union of the laminations $\mathcal{F}_{\text {can }}(I, J)$ for all pairs of different components $I$ and $J$ of $\partial \hat{S}$. By (2.1), this lamination is $\Gamma$-invariant, and hence it can be pushed forward to $\mathbf{S} \supset S$. In this way we obtain the canonical lamination on $S$ :

$$
\mathcal{F}_{\text {can }}(S)=\mathcal{F}_{\text {can }}(\mathbf{S})=\pi_{*}\left(\mathcal{F}_{\text {can }}(\hat{S})\right)
$$

The corresponding weighted arc diagram $\alpha \mapsto W_{\text {can }}(S, \alpha), \alpha \in \mathcal{A}(S)$, is called the canonical $W A D$ on $S .{ }^{6}$ By definition, it is valid.

We will now list several basic properties of the canonical WAD. The rest of the theory will be based on these properties in an essentially axiomatic way.
2.3. Property A: Maximality. Let $W_{\max }(S): \mathcal{A} \rightarrow \mathbb{R}_{+}$stand for the functional assigning to an arc $\alpha \in \mathcal{A}$ the extremal width of the family of all proper paths $\gamma$ in $S$ representing $\alpha$. (Note that $W_{\max }(S)$ is not a WAD as it is not supported on an arc diagram.)

[^3]Lemma 2.3. For any valid arc diagram $W$ on $S$, we have:

$$
W \leqslant W_{\max }(S) \leqslant W_{\operatorname{can}}(S)+2
$$

Proof. The first inequality is obvious, so let us focus on the second one.
It is trivial for any arc $\alpha \in \mathcal{A}$ with $W_{\max }(\alpha) \leqslant 2$. Let us consider some $\operatorname{arc} \alpha \in \mathcal{A}$ with $W_{\max }(\alpha)>2$. This arc connects two boundary components, $\sigma$ and $\omega$, of $S$.

The path family $\mathcal{G}(\alpha)$ representing $\alpha$ lifts to a path family $\hat{\mathcal{G}}(\alpha)$ consisting of all the paths in $\mathbb{D}$ that connect two appropriate arcs on $\mathbb{T}, I$ and $J$, covering $\sigma$ and $\omega$ respectively. Viewing $I$ and $J$ as the horizontal sides of the rectangle $Q(I, J)$ based on $\overline{\mathbb{D}}$, we obtain the desired estimate:

$$
\mathcal{W}(\mathcal{G}(\alpha)) \leqslant \mathcal{W}(\hat{\mathcal{G}}(\alpha)) \leqslant \mathcal{W}(Q(I, J))=W_{\text {can }}(S, \alpha)+2
$$

(where we have made use of Lemma 10.1 for the first estimate).

### 2.4. Property B: Natural behavior under coverings.

Lemma 2.4. If $f: U \rightarrow V$ is a finite-degree covering then $W_{\text {can }}(U)=$ $f^{*} W_{\text {can }}(V)$.
Proof. Let $\pi_{U}: \hat{U} \equiv \overline{\mathbb{D}} \backslash \Lambda_{U} \rightarrow U$ and $\pi_{V}: \hat{V} \equiv \overline{\mathbb{D}} \backslash \Lambda_{V} \rightarrow V$ be the universal coverings of $U$ and $V$, with deck transformations groups $\Gamma_{U}$ and $\Gamma_{V}$ respectively. Since $f$ has a finite degree, the group $\Gamma_{U}$ has a finite index in $\Gamma_{V}$. It follows that $\Lambda_{U}=\Lambda_{V}$, so that $\hat{U}=\hat{V}$. Hence $\mathcal{F}_{\text {can }}(\hat{U})=\mathcal{F}_{\text {can }}(\hat{V}) \equiv \mathcal{F}$. Then

$$
\mathcal{F}_{\text {can }}(U)=\mathcal{F} / \Gamma_{U}=f^{*}\left(\mathcal{F} / \Gamma_{V}\right)=f^{*}\left(\mathcal{F}_{\text {can }}(V)\right),
$$

and the conclusion follows.
2.5. Property C: Behavior under partially proper maps. Let $\mathcal{E}(S)$ stand for the set of ends of the Riemann surface $S$ (that, under our standing assumption, can be identified with the set of boundary components of $\mathbf{S})$. We say that $S$ is partially marked if we have chosen a subset $\mathcal{E}_{p}(S) \subset \mathcal{E}(S)$ of ends that we call "proper".

A map $e: U \rightarrow V$ between partially marked Riemann surfaces is called partially proper if it is proper on proper ends. An $\operatorname{arc} \alpha \in \mathcal{A}(S)$ is called horizontal if it connects proper ends of $S$. Let $\mathcal{A}^{\mathrm{h}}(S)$ stand for the set of horizontal arcs on $S$, and let $\mathcal{W}^{\mathrm{h}}(S)$ be the set of horizontal WAD's on $S$. The horizontal canonical $W A D W_{\text {can }}^{\mathrm{h}}(S)$ is the restriction of the canonical WAD to the set of horizontal arcs.

Notice that a partially proper map $e: U \rightarrow V$ induces a push-forward map on the horizontal arcs, $e_{*}: \mathcal{A}^{\mathrm{h}}(U) \rightarrow \mathcal{A}^{\mathrm{h}}(V)$, and hence, a pullback map on the horizontal arc diagrams, $e^{*}: \mathcal{W}^{\mathrm{h}}(V) \rightarrow \mathcal{W}^{\mathrm{h}}(U)$, defined by $e^{*}(Y)(\alpha)=Y\left(e_{*} \alpha\right)$.

Lemma 2.5. Let $U$ and $V$ be partially marked Riemann surfaces, and let $e: U \rightarrow V$ be a partially proper holomorphic map. Then

$$
W_{\mathrm{can}}^{\mathrm{h}}(U) \leqslant e^{*} W_{\mathrm{can}}^{\mathrm{h}}(V)
$$

Proof. Let us consider a horizontal arc $\alpha \in \operatorname{supp} W_{\text {can }}^{\mathrm{h}}(U)$. It connects two proper ends, $\sigma$ and $\omega$, and it lifts to an arc in $\hat{U}$ connecting some components $I$ and $J$ of $\partial \hat{U}$. Let $Q$ be the quadrilateral based on $\hat{U}$ with the horizontal sides $I$ and $J$. Let $I^{\prime}=\hat{e}(I), J^{\prime}=\hat{e}(J)$, and let $Q^{\prime}$ be the corresponding quadrilateral based on $\hat{V}$. Then $\hat{e}: Q \rightarrow Q^{\prime}$ maps the horizontal sides of $Q$ to the corresponding horizontal sides of $Q^{\prime}$. By Corollary $10.3, \mathcal{W}(Q) \leqslant \mathcal{W}\left(Q^{\prime}\right)$. But $\mathcal{W}(Q)=W_{\text {can }}(U, \alpha)+2$, $\mathcal{W}\left(Q^{\prime}\right)=W_{\text {can }}\left(V, e_{*}(\alpha)\right)+2$, and the desired conclusion follows.

Similarly, an arc $\alpha \in \mathcal{A}(S)$ is called vertical if it connects a proper end of $S$ to an improper one. The vertical canonical $W A D W_{\text {can }}^{\mathrm{v}}(S)$ is the restriction of the canonical WAD to the set of vertical arcs.
2.6. Property D: Domination. Let us now introduce an important relation between WAD's.

An integer WAD a (IWAD) is a WAD with integer coefficients, that is, a formal linear combination $\sum n_{i} \alpha_{i}$, where $\alpha_{i}$ is a arc diagram, and $n_{i} \in \mathbb{N}$. It will be also convenient to write a as a formal sum of arcs, $\mathbf{a}=\sum \alpha_{j}$, where any two different $\operatorname{arcs} \alpha_{j}$ are non-crossing. The order on WAD's induces a natural order on IWAD's.

Let us consider two Riemann surfaces, $U \subset V$. Given a path $\gamma$ on $V$, the restriction $\gamma \cap U$ has only finitely many non-trivial components, $\left(\gamma_{i}\right)_{i=1}^{n}$. They represent a sequence of $\operatorname{arcs}$ on $U, I(\gamma) \equiv\left(\alpha_{i} \equiv\left[\gamma_{i}\right]\right)_{i=1}^{n}$, called the itinerary of $\gamma$.

We say that a sequence of $\operatorname{arcs}\left(\alpha_{i}\right)$ on $U$ arrows an $\operatorname{arc} \beta$ on $V$ if there exists a path $\gamma$ representing $\beta$ such that $I(\gamma)=\left(\alpha_{i}\right)$. We will use notation $\left(\alpha_{i}\right) \longrightarrow \beta$ for the arrow relation.

Remark. Note that the endpoint of $\gamma_{i}$ is connected to the beginning of $\gamma_{i+1}$ by a path that goes through some component $K$ of $V \backslash U$. In this case, the end of $U$ corresponding to this component is not properly embedded into $V$. This remark is useful as it reduces a number of possibilities of how the arc $\beta$ can be composed by arcs $\alpha_{j}$.

We say that a $I W A D \mathbf{a}=\sum \alpha_{i}$ arrows $\beta$ if the $\operatorname{arcs} \alpha_{i}$ can be ordered so that the string of $\operatorname{arcs}\left(\alpha_{i}\right)$ arrows $\beta$. (In other words, we "abelianize" the arrow relation.) We use the same notation, $\mathbf{a} \longrightarrow \beta$, for this arrow relation.

Let us now consider two WAD's, $X \in \mathcal{W}(U)$ and $Y \in \mathcal{W}(V)$. We say that $X$ dominates $Y$, written

$$
X \multimap Y,
$$

if we can write

$$
X \geqslant \sum_{i} \sum_{j} w_{i j} \alpha_{i j}, \quad Y=\sum_{i} v_{i} \beta_{i}
$$

where, for each $i$,

$$
\left(\alpha_{i j}\right) j \longrightarrow \beta_{i}
$$

and

$$
\bigoplus_{j} w_{i j} \geqslant v_{i} .
$$

The basic example comes from laminations on $V$ :
Lemma 2.6. Given a valid $W A D Y$ on $V$, there exists a valid $W A D$ $X$ on $U$ such that $X \multimap Y$.

Proof. Since $Y$ is valid, $Y=W_{\mathcal{F}}$ for some lamination $\mathcal{F}$ on $V$. For $\beta \in \operatorname{supp} W_{\mathcal{F}}$, let $\mathcal{F}(\beta)$ be the sublamination of $\mathcal{F}$ assembled by the leaves $\gamma$ representing the arc $\beta$.

Let us consider the slice of $\mathcal{F}$ on $U$, that is, let $\mathcal{H}=\mathcal{F} \cap U$ and $X=$ $W_{\mathcal{H}}$. To any leaf $\gamma$ of $\mathcal{F}$, let us associate its itinerary $I(\gamma)=\left(\alpha_{j}(\gamma)\right)$ on $U$. Let $\mathcal{I}(\beta)$ stand for the set of all non-trivial itineraries $\mathbf{a}=I(\gamma)$ corresponding to all possible leaves $\gamma$ of $\mathcal{F}(\beta)$. By definition, $\mathbf{a} \longrightarrow \beta$ for any $\mathbf{a} \in \mathcal{I}(\beta)$. Let $\mathcal{F}(\beta, \mathbf{a})$ stand for the sublamination of $\mathcal{F}(\beta)$ assembled by the leaves $\gamma$ with itinerary a, i.e., $I(\gamma)=\mathbf{a}$.

For $\mathbf{a}=\left(\alpha_{j}\right) \in \mathcal{I}(\beta)$, let $v(\beta, \mathbf{a})=W(\mathcal{F}(\beta, \mathbf{a}))$, and let $w_{j}(\beta, \mathbf{a})$ be the width of the lamination assembled by the segments of $\mathcal{F}(\beta, \mathbf{a}) \cap U$ corresponding to $\alpha_{j}$. By the Series Law,

$$
\bigoplus_{j} w_{j}(\beta, \mathbf{a}) \geqslant v(\beta, \mathbf{a}) .
$$

Moreover,

$$
X=\sum_{\beta} \sum_{\mathbf{a} \in \mathcal{I}(\beta)} \sum_{j} w_{j}(\beta, \mathbf{a}) \alpha_{j}, \quad Y=\sum_{\beta} \sum_{\mathbf{a} \in \mathcal{I}(\beta)} v(\beta, \mathbf{a}) \beta
$$

This would mean that $X \multimap Y$ if we knew that $\mathcal{I}(\beta)$ were finite. The rest of the argument will show that the terms can be grouped so that the decompositions become finite.

Let $T(\beta)$ be the set of all IWAD's $\boldsymbol{\alpha}=\sum \alpha_{j}$ corresponding to all itineraries $\mathbf{a}=\left(\alpha_{j}\right) \in \mathcal{I}(\gamma)$, and let $\eta: \mathcal{I}(\beta) \rightarrow T(\beta)$ be the corresponding abelianization projection. Let $\left(e_{i}\right)_{i=1}^{n}$ be the maximal family of different non-trivial arcs represented by the leaves of $\mathcal{F} \cap U$. Then any

IWAD $\boldsymbol{\alpha} \in T(\beta)$ has a form $\sum n_{i} e_{i}, n_{i} \in \mathbb{Z}_{+}$, and hence represents an element of the free semi-group $\left(\mathbb{Z}_{+}\right)^{n}$. By Lemma 2.7 below, there is a finite set $B(\beta) \subset T(\beta)$ such that any $\boldsymbol{\alpha} \in T(\beta)$ is greater or equal than some $\mathbf{a}^{\prime} \in B(\beta)$. Let us show that the labeling set $\mathcal{I}(\beta)$ can be replaced with $B(\beta)$.

Let us select a projection $\pi: \mathcal{I}(\beta) \rightarrow B(\beta)$ factored through $\eta$ and such that $\mathbf{a}^{\prime} \equiv \pi(\mathbf{a}) \leqslant \eta(\mathbf{a})$. Then the components of $\mathbf{a}^{\prime}=\sum \alpha_{k}^{\prime}$ is a subset of components of $\mathbf{a}=\left(\alpha_{j}\right)$, so that, we can select an injective $\operatorname{map} j(k)=j_{\mathbf{a}}(k)$ such that $\alpha_{k}^{\prime}=\alpha_{j(k)}$. Let

$$
w_{k}^{\prime}\left(\beta, \mathbf{a}^{\prime}\right)=\sum_{\pi(\mathbf{a})=\mathbf{a}^{\prime}} w_{j_{\mathbf{a}}(k)}(\beta, \mathbf{a}), \quad v^{\prime}\left(\beta, \mathbf{a}^{\prime}\right)=\sum_{\pi(\mathbf{a})=\mathbf{a}^{\prime}} v(\beta, \mathbf{a}) .
$$

Then

$$
X \geqslant \sum_{\beta} \sum_{\mathbf{a}^{\prime} \in B(\beta)} \sum_{k} w_{k}^{\prime}\left(\beta, \mathbf{a}^{\prime}\right) \alpha_{k}^{\prime}, \quad Y=\sum_{\beta} \sum_{\mathbf{a}^{\prime} \in B(\beta)} v^{\prime}\left(\beta, \mathbf{a}^{\prime}\right) \beta
$$

Moreover, by Lemma 11.8

$$
\bigoplus_{k} w_{k}^{\prime}\left(\beta, \mathbf{a}^{\prime}\right) \geqslant v^{\prime}\left(\beta, \mathbf{a}^{\prime}\right) .
$$

Thus, $X \multimap Y$.
Let us consider a free Abelian semigroup $S=\mathbb{Z}_{+}^{n}$ with the standard basis $\left(e_{i}\right)_{i=1}^{n}$. It is ordered coordinatewise: $\sum x_{i} e_{i} \geqslant \sum y_{i} e_{i}$ if $x_{i} \geqslant y_{i}$ for all $i$. For any $x \in S$, let $S_{x}=\{y \in S: y \geqslant x\}$ stand for the cone with the vertex at $x$.

Lemma 2.7. Let $T$ be an arbitrary subset of the semi-group $S$. Then there exist a finite subset $B \subset T$ such that $T \subset \bigcup_{x \in B} S_{x}$.
Proof. Suppose first that $T$ is non-empty. Let $a=\sum a_{i} e_{i}$ be an element of $T$. Let $\pi_{i}: \mathbb{Z}_{+}^{n} \rightarrow \mathbb{Z}$ be the projection onto the $i^{\text {th }}$ coordinate. For $i \leqslant n$, and $k<a_{i}, k \in \mathbb{Z}_{+}$, let $T_{k}^{i}=\pi_{i}^{-1}(\{k\}) \cap T$. Then, by induction on $n$, there is a finite $B_{k}^{i} \subset T_{k}^{i}$ such that

$$
T_{k}^{i} \subset \bigcup_{x \in B_{k}^{i}} S_{x}
$$

Then let

$$
B=\{a\} \cup \bigcup_{1 \leqslant i \leqslant n} \bigcup_{k<a_{i}} B_{k}^{i} ;
$$

it has the desired property.
We can now prove Property D of canonical WAD's:

Lemma 2.8. Let $U \subset V$. Then there exists a WAD $B \in W(U)$ with $\|B\|_{\infty} \leqslant 2$ such that

$$
W_{\text {can }}(U)+B \multimap W_{\text {can }}(V)
$$

Proof. Since $W_{\text {can }}(V)$ is valid, Lemma 2.6 gives us a WAD $X$ on $U$ such that $X \multimap W_{\text {can }}(Y)$. By the Maximality Property, $X \leqslant W_{\text {can }}(U)+2$. Hence there exists a WAD $B$ with $\operatorname{supp} B \subset \operatorname{supp} X,\|B\|_{\infty} \leqslant 2$, and $X \leqslant W_{\text {can }}(U)+B$. The conclusion follows.

Lemma 2.9. If

$$
\bigoplus\left(x_{i}+b_{i}\right) \geqslant y
$$

then

$$
\bigoplus x_{i} \geqslant y-\sum b_{i} .
$$

Proof.

$$
\frac{\partial}{\partial x_{i}} \oplus x_{k}=\frac{\left(\oplus x_{k}\right)^{2}}{x_{i}^{2}} \leqslant 1
$$

Therefore

$$
\bigoplus\left(x_{i}+b_{i}\right) \leqslant \bigoplus x_{i}+\sum b_{i} .
$$

Lemma 2.10. If $X+B \multimap Y$, then $X \multimap Y-\left\|B_{1}\right\|_{1}$.
Proof. Now suppose $X+B \multimap Y$. Formally we can write $X+B \geqslant \sum_{i} T_{i}$, $Y=\sum Y_{i}$, where $T_{i}=\sum_{j} w_{i j} \alpha_{i j}, Y_{i}=v_{i} \beta_{i}$, and

$$
\begin{equation*}
\left(\alpha_{i j}\right)_{j} \rightarrow \beta_{i}, \tag{2.11}
\end{equation*}
$$

and

$$
\bigoplus_{j} w_{i j}>v_{i} .
$$

By the general theory of positive vectors in $\mathbb{R}^{n}$, we can write $T_{i}=$ $X_{i}+B_{i}$, where $X \geqslant \sum X_{i}$, and $B \geqslant \sum B_{i}$. So writing $X_{i}=\sum_{j} w_{i j}^{X} \alpha_{i j}$, and $B_{i}=\sum_{j} w_{i j}^{B} \alpha_{i j}$, we obtain

$$
\bigoplus_{j}\left(w_{i j}^{X}+w_{i j}^{B}\right) \geqslant v_{i},
$$

so, by Lemma 2.9,

$$
\bigoplus_{j} w_{i j}^{X} \geqslant v_{i}-\left\|B_{i}\right\|_{1},
$$

and therefore (using (2.11))

$$
X \geqslant \sum_{i} X_{i} \multimap \sum_{i}\left(v_{i}-\left\|B_{i}\right\|_{1}\right) \beta_{i} \geqslant Y-\|B\|_{1} .
$$

Together with Lemma 2.8, this implies:
Corollary 2.12. Let $U \subset V$. Then

$$
W_{\text {can }}(U) \multimap W_{\text {can }}(V)-6|\chi(U)| .
$$

2.7. Property E: Intersection number and hyperbolic length. Let $S$ be a compact Riemann surface with boundary. We have defined $W_{\text {can }}(S)$, and we have also defined the the weighted arc-diagram $M^{S}$ in Appendix C. We can relate $M^{S}$ to $W_{\text {can }}^{S}$ via the following:

Lemma 2.13. Let $Q$ be a quadrilateral with horizontal sides $I_{1}$ and $I_{2}$, endowed with the hyperbolic metric. Let $\Gamma_{1}$ and $\Gamma_{2}$ be the hyperbolic geodesics in $\Pi$ that share the endpoints with $I_{1}$ and $I_{2}$ respectively. Then

$$
\mathcal{W}(Q)=-\frac{2}{\pi} \log \operatorname{dist}\left(\Gamma_{1}, \Gamma_{2}\right)+O(1)
$$

Proof. We can map $Q$ conformally to the infinite strip

$$
\Pi \equiv\{z: 0<\Im z<\pi\}
$$

such that $\Gamma_{1}$ and $\Gamma_{2}$ map to vertical transverse segments with real part 0 and $d$ respectively, where $d \equiv \operatorname{dist}\left(\Gamma_{1}, \Gamma_{2}\right)$. Then the vertical sides of $Q$ map to $J_{1} \equiv[0, d] \times\{0\}$ and $J_{2} \equiv[0, d] \times\{\pi\}$. Then $\mathcal{W}(Q)$ is equal to $\mathcal{L}(\Theta)$, where $\Theta$ is the family of paths in $\Pi$ connecting $J_{1}$ and $J_{2}$. By two applications of the reflection principle (see [A2]), we find that $\mathcal{L}(\Theta)=4 \mathcal{L}\left(\Theta^{\prime}\right)$, where $\Theta^{\prime}$ is the family of paths connecting $J_{1}$ to the boundary of $\Pi^{\prime} \equiv\{z:|\Im z|<\pi / 2\}$. By the round annulus theorem[McM2]

$$
\mathcal{L}\left(\Theta^{\prime}\right)=\frac{1}{2 \pi} \log \frac{\pi / 2}{d}+O(1)=-\frac{1}{2 \pi} \log d+O(1)
$$

when $d$ is bounded above. The theorem follows.
We then have the following corollary:
Lemma 2.14. For any arc $\alpha \in \mathcal{A}(S)$,

$$
\left|W_{\mathrm{can}}^{S}(\alpha)-\frac{2}{\pi} M^{S}(\alpha)\right|<C_{0}
$$

Proof. Given $\alpha$, let us consider a lift $\tilde{\alpha}$ to the universal cover. It connects two arcs $I_{1}, I_{2}$ on the circle that cover the boundary curves of $S$ that the endpoints of $\alpha$ lie on. We can lift $h: \mathbf{S} \rightarrow S$ to $\tilde{h}: \boldsymbol{R} \rightarrow \mathbb{D}$, where $\boldsymbol{R}$ is the convex hull of the limit set of the deck transformation group. Then letting $\Gamma_{1}, \Gamma_{2}$ connect the endpoints of $I_{1}, I_{2}$ respectively, we find that the transverse geodesic arc $\boldsymbol{\alpha}$ lifts to the common perpendicular segment to $\Gamma_{1}$ and $\Gamma_{2}$. So $\operatorname{dist}\left(\Gamma_{1}, \Gamma_{2}\right)=L(\boldsymbol{\alpha})$, and the Lemma follows from Lemma 2.13.

From Theorem 12.6.1 and Lemma 2.14 we can immediately conclude:
Corollary 2.15. Let $S$ be a compact Riemann surface with boundary, and endow Int $S$ with its Poincaré metric. Let $\gamma$ be a peripheral closed geodesic in Int $S$. Then

$$
L(\gamma)=\pi\left\langle W_{\text {can }}(S), \gamma\right\rangle+O(1 ; \chi(S))
$$

With $S$ as in Corollary 2.15, we let $\Gamma_{S}$ be the formal sum of the peripheral closed geodesics in $\operatorname{Int} S$. Then

$$
\left\langle W_{\text {can }}(S), \Gamma_{S}\right\rangle=2\left\|W_{\text {can }}(S)\right\|_{1},
$$

because every arc in $S$ intersects exactly two peripheral geodesics exactly once (or intersects exactly one exactly twice). Therefore, by Corollary 2.15,

$$
\begin{equation*}
\left\|W_{\text {can }}(S)\right\|_{1}=\frac{\pi}{2}|\Gamma|+O(1 ; p) . \tag{2.16}
\end{equation*}
$$

We can use this equation to compare the total weight of $W_{\text {can }}(S)$ between two surfaces:

Theorem 2.17. Suppose $S$ and $S^{\prime}$ are the interiors of compact Riemman surfaces, and there is a holomorphic map $i: S^{\prime} \rightarrow S$ that is homotopic to a homeomorphism. Then

$$
\left\|W_{\text {can }}(S)\right\|_{1} \leqslant\left\|W_{\text {can }}\left(S^{\prime}\right)\right\|_{1}+O(1 ; \chi(S))
$$

Proof. Because $i$ is homotopic to a homeomorphism, it bijectively maps the free homotopy classes of the peripheral geodesics of $S^{\prime \prime}$ to the same for $S$. Then, by the Schwarz Lemma, because $i$ is holomorphic, we have $\left|\Gamma_{S^{\prime}}\right| \geqslant\left|\Gamma_{S}\right|$. The Theorem then follows from the application of (2.16) to $\Gamma_{S^{\prime}}$ and $\Gamma_{S} \mid$.

## 3. Theorem 1.2 in terms of $W_{\text {can }}$

In this brief section we will prove Theorem 1.2, given a simple theorem about $W_{\text {can }}\left(V \backslash \mathcal{K}_{p}\right)$. We then outline the proof of this latter theorem, which is then proven in Sections 4-7.
3.1. Quadratic-like maps and their renormalizations. For the basic theory of quadratic-like maps, see [DH2].

A quadratic-like map is a double branched covering $f: U \rightarrow V$, where $U$ and $V$ are topological disc in $\mathbb{C}$, and $U \Subset V \neq \mathbb{C}$. Such a map has a unique critical point, which will be placed at the origin. The filled Julia set $K(f)$ is the set of non-escaping points:

$$
K(f)=\left\{z: f^{n} z \in U, n=0,1, \ldots\right\} .
$$

It is either connected or a Cantor set depending on whether $0 \in K(f)$ or not. (In what follows, we will often skip "filled", since we will never deal with the actual "Julia sets".)

Two quadratic-like maps are called hybrid equivalent if they are quasi-conformally conjugate by a map $h$ with $\bar{\partial} h=0$ a.e. on $K(f)$. The Douady and Hubbard Straightening Theorem asserts that any quadratic-like map with connected Julia set is hybrid equivalent to a quadratic-like restriction of a unique quadratic polynomial.

A quadratic-like map $f$ is called primitively renormalizable with period $p$ if there exists a topological disk $U^{\prime} \ni 0$ such that the domains $f^{k}\left(U^{\prime}\right), k=0,1, \ldots, p-1$, are contained in $U$ and pairwise disjoint, and $g=f^{p}: U^{\prime} \rightarrow f^{p}\left(U^{\prime}\right)$ is a quadratic-like map with connected Julia set. In this case, $K_{0}=K(g)$, as well as $K_{j}=f^{j}\left(K^{\prime}\right), j=0,1, \ldots, p-1$, are called the little Julia set, and the invariant set

$$
\mathcal{K}=\bigcup_{j=0}^{p-1} K_{j}
$$

is called the cycle of little Julia sets.
3.2. Horizontal and vertical arc diagrams. In the further applications, the Riemann surface $S$ will be $V \backslash \mathcal{K}$, where $V$ is a topological disk, and $\mathcal{K}$ is the union of $p$ disjoint full continua (in reality, the cycle of little Julia sets of some quadratic-like map).

Under these circumstances, a proper path (and the corresponding arc) in $V \backslash \mathcal{K}$ is called horizontal if it connects two little Julia sets, and is called vertical if it connects a little Julia set to $\partial V$. Given an arc diagram $W$ on $V \backslash \mathcal{K}$, let $W^{\mathrm{h}}$ and $W^{\mathrm{v}}$ stand for its horizontal and vertical parts respectively, and let $W^{\mathrm{v}+\mathrm{h}}$ stand for their sum. In particular, we can consider $W_{\text {can }}^{\mathrm{h}}, W_{\text {can }}^{\mathrm{v}}$ and $W_{\text {can }}^{\mathrm{v}+\mathrm{h}}$.

We can now state a theorem which will be proven in Section 7:
Theorem 3.1. Let $f: \mathbf{U} \rightarrow \mathbf{U}$ be primitively renormalizable quadratic-like-map with period $p$, and let $\mathcal{K}$ be its cycle of little Julia sets. Then there exists $M=M(p)$ such that: If $\left\|W_{\mathrm{can}}^{\mathrm{v}+\mathrm{h}}(\mathbf{U} \backslash \mathcal{K})\right\|_{1}>M$ then

$$
\left\|W_{\mathrm{can}}^{\mathrm{v}}(\mathbf{U} \backslash \mathcal{K})\right\|_{1} \geqslant C^{-1}\left\|W_{\mathrm{can}}^{\mathrm{h}}(\mathbf{U} \backslash \mathcal{K})\right\|_{1}
$$

with an absolute constant $C>0$.
Given Theorem 3.1, we can prove Theorem 1.2:
We are now ready to show that the modulus of a $\psi$-quadratic-like map is improving under the renormalization.

Let $f: \mathbf{U} \rightarrow \mathbf{U}$ be a quadratic-like map with filled Julia set $K$. The simple closed geodesic $\gamma$ in the annulus $\mathbf{U} \backslash K$ is called the geodesic associated with $f$. Let $|\gamma|$ stand for its hyperbolic length in $\mathbf{U} \backslash K$.

Proof of Theorem 1.2. Let $K_{p}(j)=f^{j}\left(K_{p}\right)$ be the $j^{\text {th }}$ small filled-Julia set, and let $\gamma_{p}(j)$ be the peripheral simple closed geodesic in $\mathbf{U} \backslash \mathcal{K}_{p}$ going around $K_{p}(j)$. Let $\Gamma$ be peripheral closed geodesic in $\mathbf{U} \backslash \mathcal{K}_{p}$ homotopic to $\partial \mathbf{U}$. We let $W_{\text {can }} \equiv W_{\text {can }}\left(\mathbf{U} \backslash \mathcal{K}^{\prime}\right)$.

The geodesic $\Gamma$ intersects each vertical arc once and does not intersect horizontal arcs. By Corollary 2.15,

$$
\begin{equation*}
|\Gamma| \geqslant c\left\|W_{\text {can }}^{\mathrm{v}}\right\|_{1}+O(1) \tag{3.2}
\end{equation*}
$$

Let $W_{\text {can }} \mid j$ be the part of $W_{\text {can }}$ supported on the arcs landing at $K_{p}(j)$. The geodesic $\gamma_{p}(j)$ does not intersect $W_{\text {can }}-W_{\text {can }} \mid j$ and intersects each arc of $\operatorname{supp} W_{\text {can }} \mid j$ once. By Corollary 2.15,

$$
\begin{equation*}
\left|\gamma_{p}(j)\right|=c\left\|W_{\text {can }} \mid j\right\|_{1}+O(1) . \tag{3.3}
\end{equation*}
$$

But by the Schwarz Lemma, the geodesics $\gamma_{p}(j)$ have comparable lengths (see [McM1, Theorem 9.3]):

$$
\frac{1}{2}\left|\gamma_{p}\right| \leqslant\left|\gamma_{p}(j)\right| \leqslant\left|\gamma_{p}\right|
$$

Putting this together with (3.3), we see that all the $\left\|W_{\text {can }} \mid j\right\|_{1}$ are also comparable, provided $\left|\gamma_{p}\right|$ is sufficiently big (bigger than some absolute $L_{0}$ ). Hence

$$
\left\|W_{\text {can }} \mid 0\right\|_{1}=\frac{1}{p}\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}\right\|_{1},
$$

and together with (3.3), we obtain:

$$
\left|\gamma_{p}\right|=\frac{1}{p}\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}\right\|_{1} .
$$

Now assume $\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}\right\|_{1}>M(\bar{p})$. Then Theorem 3.1 implies that $\left\|W_{\text {can }}\right\|_{1}$ is comparable with $\left\|W_{\text {can }}^{\mathrm{v}}\right\|_{1}$, and we can conclude

$$
\left|\gamma_{p}\right| \leqslant \frac{C}{p}\left\|W_{\text {can }}^{\mathrm{v}}\right\|_{1},
$$

which, combined with (3.2), implies

$$
\left|\gamma_{p}\right| \leqslant \frac{C}{p}|\Gamma| .
$$

We now need only observe that by the Schwarz Lemma, $|\Gamma| \leqslant|\gamma|$ (since $\left.K \supset \mathcal{K}_{p}\right)$.
3.3. Outline for the proof of Theorem 3.1. In this informal subsection we will outline the proof of Theorem 1.2 that will be presented in Sections 5, 6, and 7, using the background in Section 4.

Suppose that $f: U \rightarrow V$ is a primitively $p$-renormalizable quadraticlike map. Theorem 1.2 relates two hyperbolic lengths in $V \backslash \mathcal{K}_{p}$, which we will reinterpret in terms of $W_{\text {can }}\left(V \backslash \mathcal{K}_{p}\right)$ in light of eqrefeq:lengthwcan. Moreover, letting $K_{p}(k)=f^{k}\left(K_{p}\right)$ be the $k^{\text {th }}$ small filled-Julia set, and $\gamma_{1}(k)$ be the geodesic in $V \backslash \mathcal{K}_{p}$ going around it, we know that $\left|\gamma_{1}\right| \leqslant\left|\gamma_{1}(k)\right| \leqslant 2\left|\gamma_{1}\right|$. It is then not too hard to see that Theorem 1.2 follows from the following:

$$
\begin{equation*}
\left\|W^{\mathrm{v}+\mathrm{h}}\right\|_{1}>M \Longrightarrow\left\|W^{\mathrm{v}}\right\| \geqslant \epsilon\left\|W^{\mathrm{v}+\mathrm{h}}\right\| . \tag{3.4}
\end{equation*}
$$

We will prove (3.4) as follows:
Restrictions of the canonical WAD's and the entropy argument. We consider the nest of the domains $\mathbf{U}^{m}$, pullbacks of $V$ under the iterates of $f$, and restrict the canonical WAD to these domains. Let $X^{m}$ stand for the horizontal parts of these restrictions (with appropriate constants subtracted). We show that $f^{*} X^{m} \multimap X^{m+1}$. This allows us to conclude that eventually the diagrams $X^{n}$ are aligned with the Hubbard tree (§4, 5) and that

$$
\left\|X^{10 p}\right\|_{1} \leqslant \frac{1}{2}\left\|X^{0}\right\|_{1}
$$

which implies

$$
\left\|W_{\mathrm{can}}^{\mathrm{h}}\left(\mathbf{U}^{10 p} \backslash \mathcal{K}_{p}\right) \left\lvert\, \leqslant \frac{1}{2}\right.\right\| W_{\mathrm{can}}^{\mathrm{h}}\left(\mathbf{U}^{0} \backslash \mathcal{K}_{p}\right) \|+O(1 ; p)
$$

(§5). The latter property depends on the positivity of entropy of the Hubbard tree dynamics: this is the main place where we use that the renormalizations are of primitive type.

Push-forward via the Covering Lemma (§7). The horizontal width released under restrictions is turned into the vertical width, which implies

$$
\left\|W_{\text {can }}\left(\mathbf{U}^{10 p} \backslash \mathcal{K}_{p}\right)\right\|_{1} \geqslant \frac{1}{2}\left\|W_{\text {can }}^{\mathrm{h}}\left(\mathbf{U}^{10 p} \backslash \mathcal{K}_{p}\right)\right\|_{1}-O(1 ; p) .
$$

We can then apply the Covering Lemma[KL1] to push forward $W_{\text {can }}\left(\mathbf{U}^{10 p} \backslash\right.$ $\mathcal{K}_{p}$ ) by $f^{10 p}$ and obtain (3.4).

## 4. Life on Hubbard trees

4.1. Topological disked trees and aligned graphs. Let us consider a Riemann surface $U$ with finitely many open Jordan disks $D_{j} \Subset U$ with disjoint closures, and let $\mathcal{D}=\cup D_{j}$. An embedded 1-complex
$H_{k} \subset U \backslash \mathcal{D}$ is called proper if all its tips (i.e., valence 1 vertices) belong to $\partial \mathcal{D}$. Assume that we have finitely many disjoint proper graphs $H_{k}$ such that $H=\mathcal{D} \cup_{k} H_{k}$ is simply connected. In this case we say that $H$ is a (topological) disked tree.

We say that a path in $H \backslash \mathcal{D}$ is aligned with a disked tree $H$ if it connects the boundaries of two disks in $\mathcal{D}$. The $\operatorname{arcs} \alpha$ in $U \backslash \mathcal{D}$ represented by aligned paths are also called aligned with $H$. Let $\mathbf{H}$ stand for the family of the aligned paths/arcs. (Since there is a natural one-to-one correspondence between the aligned arcs and the aligned paths, we will not distinguish notationally these families).

Let $\mathbf{G}$ be the abstract graph whose vertices are disks $D_{i}$ and edges are the paths/arcs of $\mathbf{H}$. We call it the graph of aligned arcs.
Lemma 4.1. The graph $\mathbf{G}$ is a tree of complete graphs. ${ }^{7}$
Proof. Let $\mathcal{D}_{k}$ stand for the family of the disks $D_{i}$ that touch the graph $H_{k}$. Let $\mathbf{G}_{k} \subset G$ be the graph whose vertices are disks $D_{i} \subset \mathcal{D}_{k}$ and edges are the path of $\mathbf{H}$ contained in $\mathbf{G}_{k}$. As $H_{k}$ is path connected, any two disks $D_{i}$ and $D_{j}$ of $\mathcal{D}_{k}$ can be connected by an $\operatorname{arc}$ in $\mathbf{G}_{k}$, so that, the graph $\mathbf{G}_{k}$ is full.

Since $H$ is a tree, it is easy to check that the graphs $\mathbf{G}_{k}$ are organized in a tree as well.
4.2. Superattracting model. To a primitively $p$-renormalizable map $f$, one can associate a certain superattracting quadratic polynomial $F$ of period $p$ as follows (compare [McM1, $\S B]$ ). First, we can straighten $f$ to a quadratic polynomial, so we can assume that it is a quadratic polynomial in the first place. Then, collapsing the little Julia sets $K_{j}$ to points $c_{j}$, we obtain a topological sphere $S$. Moreover, the map $f$ descends to a degree two map $f_{0}: S \rightarrow S$ with a periodic critical point $c_{0}\left(f_{0}\right.$ collapses the sets $\left.-K_{j} \subset \mathbb{C} \backslash \mathcal{K}, j=1, \ldots, p-1\right)$, to points $c_{j+1}$. One can check that this map does not have Thurston obstructions (see [DH3]), so it can be realized as a superattracting quadratic polynomial. This polynomial is the desired $F$.

We call $F$ the superattracting model for $f$. We let $\mathcal{O}=\left\{F^{k}(0)\right\}_{k=0}^{p-1}$ stand for the superattracting cycle of $F$, and $\mathcal{D}$ stand for its immediate basin of attraction.

Lemma 4.2. Let $f:(U, \mathcal{K}) \rightarrow(V, \mathcal{K})$ be a primitively renormalizable quadratic-like map, and let $F:(\mathbb{C}, \mathcal{D}) \rightarrow(\mathbb{C}, \mathcal{D})$ be its superattracting model. There is a natural one-to-one correspondence between horizontal/vertical arcs in $U \backslash \mathcal{K}$ (resp., $V \backslash f^{-1}(\mathcal{K})$ ) and horizontal/vertical

[^4]arcs in $\mathbb{C} \backslash \mathcal{D}\left(\right.$ resp. $\left.\mathbb{C} \backslash F^{-1}(\mathcal{D})\right)$. The correspondence between the horizontal arcs is compatible with the arrow relation.

This obvious statement allows us to replace $f$ with its superattracting model $F$, as long as we are dealing with the combinatorics of arc diagrams. The advantage of the model is that it possesses a well-defined Hubbard tree.
4.3. Disked Hubbard tree. Let $D_{k}$ be the component of $\mathcal{D}$ containing $c_{k}$. It is known that the $\bar{D}_{k}$ 's are Jordan disks. Given a set $X \subset K(F)$, the path hull of $X$ is defined as the smallest path connected closed subset of $K(F)$ containing $X$ satisfying the property that it intersects any component of int $K(F)$ by the union of internal rays. The Hubbard tree $H \equiv H_{F}$ is defined as the path hull of the basin $\mathcal{D}$. It is a disked tree in the sense of $\S 4.1$. Moreover, it is invariant under $F$; in fact, $F(H)=H$.

Lemma 4.3. The valence of any disk $D_{k}$ of the Hubbard tree is at most 2.
4.4. No periodic horizontal arcs. Given a proper path $\gamma$ in $\mathbb{C} \backslash \mathcal{D}$ that begins at some $D_{i}, i \in \mathbb{Z} / p \mathbb{Z}$, let $F^{*} \gamma$ stand for the union of the proper lifts of $\gamma$ that begin at $D_{i-1}$. (Note that a horizontal path $\gamma$ has at most one proper lift.)

This lifting operation descends to the level of arcs. An arc a is called periodic if $\left(F^{*}\right)^{l}(\alpha) \supset \alpha$ for some $l \in \mathbb{Z}_{+}$. (Of course, the period $l$ must be a multiple of $p$.)

Lemma 4.4 (see [P], Theorem 5.8). Periodic horizontal arcs do not exist.

Proof. Assume that $\alpha$ is such a periodic horizontal arc with period $l$. Let us endow $\mathbb{C} \backslash \mathcal{O}$ with hyperbolic metric, and let us represent $\alpha$ be a geodesic path $\gamma \subset \mathbb{C} \backslash \mathcal{D}$. It is the shortest representative of $\alpha$.

Since $\left(F^{*}\right)^{l} \alpha=\alpha$, the lift $F^{*} \gamma$ represents the same arc $\alpha$. so that, it is no shorter than $\gamma$. On the other hand, the Schwarz Lemma implies that it is shorter than $\gamma$ - contradiction.
4.5. Expansion property. Let $F^{-1} \mathbf{H}$ stand for the family of paths in $H \backslash F^{-1} \mathcal{D}$ with endpoints on $F^{-1} \overline{\mathcal{D}}$. Each path of $F^{-1} \mathbf{H}$ is homeomorphically mapped by $F$ onto some path of $\mathbf{H}$.

Let us consider a $|\mathbf{H}| \times|\mathbf{H}|$ matrix $M=M(F)$ defined as follows: $M_{\gamma \delta}=n$ if $\gamma$ contains $n$ disjoint segments $\gamma_{i} \in F^{-1} \mathbf{H}$ such that $F \gamma_{i}=\delta$. One can easily check that the matrix $M^{k}$ can obtained by applying the same construction to the map $F^{k}$. .

Lemma 4.5. For any $\gamma \in \mathbf{H}$, we have $\sum_{\delta} M_{\gamma \delta}^{p} \geqslant 2$.
Proof. Not that the path $F^{p} \gamma$ has the same endpoints as $\gamma$ does. If $\sum_{\delta} M_{\gamma \delta}^{p}=1$ then $F^{p} \gamma$ would not cross any disks of $\mathcal{D}$. It would follow that $f^{p} \gamma=\gamma$, which is impossible by Lemma 4.4.
4.6. Periodic vertical arcs. Let $\mathbf{H}^{\perp}$ be the arc diagram consisting of vertical arcs in $\mathbb{C} \backslash \mathcal{D}$ that do not intersect the Hubbard tree $H$ (up to homotopy).

Lemma 4.6. If a vertical arc $\beta$ is periodic then $\beta \in \mathbf{H}^{\perp}$.
Proof. Let $\left(F^{*}\right)^{l} \beta \supset \beta$ for some $l \in \mathbb{Z}_{+}$. Let $\mathcal{G}$ be the family of paths $\gamma$ in $\mathbb{C} \backslash \mathcal{D}$ representing $\beta$. For a path $\gamma \in \mathcal{G}$, let $\gamma^{\prime}$ be its maximal initial segment whose endpoints belong to the Julia set $K(F)$. Let $\mathcal{G}^{\prime}$ be the family of all paths $\gamma^{\prime}$ that can be obtained in this way.

Since $\left(F^{*}\right)^{l} \beta \supset \beta$, any path $\gamma \in \mathcal{G}$ can be lifted by $F^{l}$ to a path $\delta \in \mathcal{G}$. Then $\delta^{\prime}$ is a lift of $\gamma^{\prime}$, and thus, we obtain the lifting map $L: \mathcal{G}^{\prime} \rightarrow \mathcal{G}^{\prime}$.

Let us endow the punctured plane $S=\mathbb{C} \backslash \mathcal{O}$ with the hyperbolic metric. Let $|\gamma|$ stand for the hyperbolic length of a path $\gamma$, and let

$$
\mu=\inf _{\gamma^{\prime} \in \mathcal{G}^{\prime}}\left|\gamma^{\prime}\right| .
$$

This infimum is realized by some hyperbolic geodesic $\gamma^{\prime}$. If $\mu>0$ then by the Schwarz Lemma, $\left|L\left(\gamma^{\prime}\right)\right|<\left|\gamma^{\prime}\right|$, contradicting minimality of $\gamma^{\prime}$. Hence $\mu=0$ and so, $\mathcal{G}$ contains a curve $\gamma$ that does not intersect the Julia set $K(F)$ (except for the initial point in $\mathcal{D}$ ). This is the desired vertical representative of $\beta$.
4.7. Pullbacks of vertical arcs. Given two families $\boldsymbol{\alpha}$ and $\boldsymbol{\beta}$ of proper arcs (or paths) in a Riemann surface $S$, the "inner product" $\langle\boldsymbol{\alpha}, \boldsymbol{\beta}\rangle \equiv\langle\boldsymbol{\alpha}, \boldsymbol{\beta}\rangle_{S}$ is the minimal number of intersection points between path representatives of $\boldsymbol{\alpha}$ and $\boldsymbol{\beta}$.

For instance, $\alpha \in \mathbf{H}^{\perp}$ if and only if $\langle\alpha, \mathbf{H}\rangle=0$. The dual statement is also valid: $\alpha \in \mathbf{H}$ if and only if $\left\langle\alpha, \mathbf{H}^{\perp}\right\rangle=0$.

Lemma 4.7. For any vertical arc $\beta$, we have: $\bigcup_{n=0}^{\infty}\left(F^{n}\right)^{*}(\beta) \supset \mathbf{H}^{\perp}$.
Proof. Let $\gamma$ be a vertical path representing $\beta$, and let $\gamma_{1}$ be a component of $F^{*} \gamma$. Since $F: \gamma_{1} \rightarrow \gamma$ is a homeomorphism,

$$
\left\langle\gamma_{1}, F^{-1} \mathbf{H}\right\rangle_{\mathbb{C} \backslash F^{-1} \mathcal{D}}=\langle\gamma, \mathbf{H}\rangle_{\mathbb{C} \backslash \mathcal{D}} .
$$

Since $H \subset F^{-1} H$,

$$
\begin{equation*}
\left\langle\gamma_{1}, \mathbf{H}\right\rangle \leqslant\langle\gamma, \mathbf{H}\rangle . \tag{4.8}
\end{equation*}
$$

Let us now consider a chain of vertical $\operatorname{arcs} \beta_{n} \subset\left(F^{n}\right)^{*}(\beta)$. By (4.8), the intersection number $\left\langle\beta_{n}, \mathbf{H}\right\rangle$ does not increase with $n$, and hence eventually stabilizes at some value $k$. But for a given $k$, there are only finitely many different vertical $\operatorname{arcs} \beta$ such that $\langle\beta, \mathbf{H}\rangle=k$. Hence $\beta_{n}=\beta_{n+l}$ for some $l>0$. By Lemma 4.6, $\beta_{n} \in \mathbf{H}^{\perp}$.

But then all further lifts of $\beta_{n}$ by the iterates of $F$ belong to $\mathbf{H}^{\perp}$ as well. Lifting $\beta_{n}$ to the critical disk $D_{0}$, we obtain two symmetric arcs, $\sigma_{1}$ and $\sigma_{2}$, landing at $D_{0}$. By symmetry, these arcs are different. But by Lemma 4.3, for any disk $D_{k}$, there exist at most two different arcs of $\mathbf{H}^{\perp}$ landing on $D_{k}$. Thus, $\sigma_{1}$ and $\sigma_{2}$ are the only arcs of $\mathbf{H}^{\perp}$ landing on $D_{0}$. Lifting these arcs further to all the domains $D_{k}$, we obtain all the arcs of $\mathbf{H}^{\perp}$.
4.8. Trees $H^{l}$ and the associated objects. For any $l \in \mathbb{Z}_{\geqslant}$, we can introduce the following objects:

- $\mathcal{D}^{l}=F^{-l}(\mathcal{D}) ; D_{k}^{l}$ are the components of $\mathcal{D}$;
- $H^{l}=F^{-l}(H)$; it is a disked tree with disks $D_{k}^{l}$;
- $\mathbf{H}^{l}$ is the family of the paths/arcs aligned with $H^{l}$;
- $\mathrm{G}^{l}$ is a graph whose vertices are the disks $D_{k}^{l}$ and edges are the arcs of $\mathbf{H}^{l}$; as for $l=0$ (Lemma 4.1), it is a tree of complete graphs $\mathbf{G}_{k}^{l}$.

Notice that $F^{l}$ maps $\mathbf{G}^{l}$ onto $\mathbf{G}$, so that each complete graph $\mathbf{G}_{k}^{l}$ is mapped onto some complete graph $\mathbf{G}_{j}$.

By Lemma 11.11 applied to $\mathbf{G}^{l}$, for any edge $\left[D_{m}, D_{n}\right]$ of $\mathbf{H}$, there is a unique chain of $\mathbf{G}^{l}$,

$$
\begin{equation*}
D_{m}=D_{k(1)}^{l}, D_{k(2)}^{l}, \ldots, D_{k(d)}^{l}=D_{n}, \tag{4.9}
\end{equation*}
$$

connecting $D_{m}$ to $D_{n}$. Here $d=d_{\mathbf{G}^{l}}\left(D_{m}, D_{n}\right)$ is the distance between $D_{m}$ and $D_{n}$ in $\mathbf{G}^{l}$. Lemma 4.5 implies:

$$
\begin{equation*}
d_{\mathbf{G}^{r p}}\left(D_{m}, D_{n}\right) \geqslant 2^{r} . \tag{4.10}
\end{equation*}
$$

In fact, we have:
Lemma 4.11. At least $2^{r-1}$ of the disks $D_{k(i)}^{r p}$ belong to $\mathcal{D}^{r p} \backslash \mathcal{D}^{(r-1) p}$.
Proof. For $r=1$, the assertion follows from Lemma 4.5 and a remark that all the intermediate disks in the path (4.9) do not belong to $\mathcal{D}$. Since each path of $\mathbf{H}^{p}$ is mapped homeomorphically under $F^{p}$ onto a path of $\mathbf{H}$ (transferring chains in $\mathbf{G}^{r p}$ to chains in $\left.\mathbf{G}^{(r-1) p}\right)$, the assertion follows by induction in $r$.

## 5. Restrictions of WAD's

5.1. Restriction of the domains. Let us consider a quadratic-like map $f: U^{1} \rightarrow U^{0}$. We let $U^{n}=f^{-n} U^{0}$, and observe that $U_{n+1} \subset U_{n}$, and $f: U_{n+1} \rightarrow U_{n}$ is quadratic-like.

Assume $f$ is primitively renormalizable with some period $p$, and let $\mathcal{K}$ be the cycle of the corresponding little (filled) Julia sets. Let $\tilde{\mathcal{K}}=$ $f^{-1}(\mathcal{K})$. Then we have an embedding $i: U^{n+1} \backslash \mathcal{K} \rightarrow U^{n} \backslash \mathcal{K}$, a covering $f: U^{n+1} \backslash \tilde{\mathcal{K}} \rightarrow U^{n} \backslash \mathcal{K}$, and an embedding $U^{n+1} \backslash \tilde{\mathcal{K}} \subset U^{n+1} \backslash \mathcal{K}$, that together form a triangle diagram. We will properly mark the Riemann surfaces $U^{n} \backslash \mathcal{K}$ and $U^{n} \backslash \tilde{\mathcal{K}}$ by declaring the ends corresponding to the little Julia sets $K_{i}$ to be proper. Then the maps in the above diagram are partially proper. Moreover, the embedding $U^{n+1} \backslash \tilde{\mathcal{K}} \subset U^{n+1} \backslash \mathcal{K}$ is also proper on $\partial U^{n+1}$.
5.2. Domination relations. Let us select a non-decreasing sequence of numbers $q_{n} \in \mathbb{R}_{+}$such that

$$
\begin{equation*}
q_{n+1} \geqslant 3 p\left(q_{n}+2\right) \tag{5.1}
\end{equation*}
$$

Let $X^{n}=W_{\text {can }}^{\mathrm{h}}\left(U^{n} \backslash \mathcal{K}\right)$ and $\hat{X}^{n}=X^{n}-q_{n}$.
Proposition 5.2. We have:
(i) $\hat{X}^{n+1} \leqslant i^{*} \hat{X}^{n}$;
(ii) $f^{*} \hat{X}^{n} \multimap \hat{X}^{n+1}$.

Proof. (i) Since the embedding $i$ is proper on the ends corresponding to the little Julia sets $K_{j}$, we have $X^{n+1} \leqslant i^{*} X^{n}$ by Property C. This yields the desired inequality since the sequence $\left(q_{n}\right)$ is non-decreasing.
(ii) Let us properly mark the ends of $U^{n} \backslash \mathcal{K}$ and $U^{n+1} \backslash \tilde{\mathcal{K}}$ corresponding to the sets $\mathcal{K}$ and $\tilde{\mathcal{K}}$ respectively. Since the covering $f$ maps $\mathbf{U}^{n+1}$ to $\mathbf{U}^{n}$ and maps $\tilde{\mathcal{K}}$ to $\mathcal{K}$, horizontal and vertical arcs of $U^{n} \backslash \mathcal{K}$ lift respectively to horizontal and vertical arcs of $U^{n+1} \backslash \tilde{\mathcal{K}}$. Hence $f^{*} X^{n}=W_{\text {can }}^{\mathrm{h}}\left(U^{n+1} \backslash \tilde{\mathcal{K}}\right)$.

By Property D,

$$
W_{\text {can }}\left(U^{n+1} \backslash \tilde{\mathcal{K}}\right) \multimap W_{\text {can }}^{\mathrm{h}}\left(U^{n+1} \backslash \mathcal{K}\right)-6 p
$$

But since the embedding $U^{n+1} \backslash \tilde{\mathcal{K}} \subset U^{n+1} \backslash \mathcal{K}$ is proper on $\partial U^{n+1}$, the itinerary of any horizontal path $\gamma$ in $U^{n+1} \backslash \mathcal{K}$ consists only of horizontal arcs of $U^{n+1} \backslash \tilde{\mathcal{K}}$. It follows that

$$
W_{\text {can }}^{\mathrm{h}}\left(U^{n+1} \backslash \tilde{\mathcal{K}}\right) \multimap W_{\text {can }}^{\mathrm{h}}\left(U^{n+1} \backslash \mathcal{K}\right)-6 p .
$$

Thus, $f^{*} X^{n} \multimap X^{n+1}-6 p$. Taking into account Lemma 2.9 and (5.1), we conclude:

$$
f^{*} \hat{X}^{n}=f^{*} X^{n}-q_{n} \multimap X^{n+1}-6 p-3 p q_{n} \geqslant \hat{X}^{n+1}
$$

In the remainder of Sections 5 and 6 , we will reason using only combinatorics and topology, combined with (i) and (ii) above, to conclude our main goal, Lemma 6.6.
5.3. Topological arrow. Let us consider two Riemann surfaces, $U \subset$ $V$, and let $\boldsymbol{\alpha}$ and $\boldsymbol{\beta}$ be multiarcs on $U$ and $V$ respectively. We say that $\boldsymbol{\alpha}$ topologically arrows $\boldsymbol{\beta}, \boldsymbol{\alpha} \rightsquigarrow \boldsymbol{\beta}$, if for any $\operatorname{arc} \beta \in \boldsymbol{\beta}$ there is a sequence $\left(\alpha_{k}\right)$ of arcs with $\alpha_{k} \in \boldsymbol{\alpha}$ for each $k$, and $\left(\alpha_{k}\right) \longrightarrow \beta$.

A basic example comes from two WAD's, $X \in W(U)$ and $Y \in W(V)$, such that $X \multimap Y$. Then $\operatorname{supp} X \rightsquigarrow \operatorname{supp} Y$, as follows immediately from the definitions.
5.4. Invariant horizontal arc diagram. Given a partially proper embedding $U \subset V$ such that $U$ is a deformation retract of $V$, we can view horizontal arcs on $V$ as horizontal arcs on $U$ (by retracting them from $V$ to $U$ ). When we'd like to emphasize this point of view, we will use notation $\alpha \mid U$.

Let us say that a horizontal $\operatorname{arc} \operatorname{diagram} \boldsymbol{\alpha}$ on $U^{n} \backslash \mathcal{K}$ is invariant if

$$
f^{*} \boldsymbol{\alpha} \rightsquigarrow \boldsymbol{\alpha}
$$

Proposition 5.3. There exists $n \leqslant 3 p$ such that the horizontal arc diagram $\boldsymbol{\alpha}^{n}=\operatorname{supp} \hat{X}^{n}$ is invariant.
Proof. Since $\hat{X}^{n+1} \leqslant i^{*} \hat{X}^{n}\left(\right.$ by Lemma 5.2 (i)), $\boldsymbol{\alpha}^{n+1} \subset \boldsymbol{\alpha}^{n} \mid U^{n+1} \backslash \mathcal{K}$. Since $\left|\boldsymbol{\alpha}^{0}\right| \leqslant 3 p$, there exists an $n \leqslant 3 p$ such that $\boldsymbol{\alpha}^{n+1}=\boldsymbol{\alpha}^{n} \mid U^{n+1} \backslash \mathcal{K}$. Since by Lemma 5.2 (ii), $f^{*} \boldsymbol{\alpha}^{n} \rightsquigarrow \boldsymbol{\alpha}^{n+1}$, we are done.
5.5. Alignment with the Hubbard tree. We say that a horizontal arc diagram a is aligned with the Hubbard tree if it is so for the superattracting realization of $\mathbf{a}$.

Lemma 5.4. Any invariant horizontal arc diagram $\boldsymbol{\alpha}$ is aligned with the Hubbard tree.

Proof. Let $\boldsymbol{\alpha}$ be an invariant horizontal arc diagram for a superattracting polynomial $F$. It can be realized by a "disked graph" $A$ whose "vertices" are the disks $\bar{D}_{k}$ and edges are paths representing the arcs of $\boldsymbol{\alpha}$. Let $\Delta$ be the unbounded component of $\mathbb{C} \backslash A$. Then one of the disks $\bar{D}_{k}$ intersects $\partial \Delta$. A path $\gamma$ in $\Delta$ landing on $\bar{D}_{k}$ represents a vertical arc $\beta$ such that $\langle\boldsymbol{\alpha}, \beta\rangle=0$.

Then $F^{*} \gamma$ does not intersect $F^{-1}(A)$ (except for the landing points). But by invariance, the arcs of $\boldsymbol{\alpha}$ can be realized as paths in $F^{-1}(A)$. Hence $F^{*} \gamma$ does not intersect $\boldsymbol{\alpha}$.

Iterating, we see that $\boldsymbol{\alpha}$ does not intersect $\left(F^{n}\right)^{*} \gamma$ for all $n=0,1, \ldots$. By Lemma 4.7, it does not intersect $\mathbf{H}^{\perp}$ either. So, $\left\langle\boldsymbol{\alpha}, \mathbf{H}^{\perp}\right\rangle=0$, and we are done.

Putting this together with Proposition 5.3 and Proposition 5.2 (i), we obtain:

Corollary 5.5. For any $n \geqslant 3 p$, the horizontal arc diagram $\boldsymbol{\alpha}^{n}=$ supp $\hat{X}^{n}$ is aligned with the Hubbard tree $H$.

## 6. Entropy argument

6.1. Domination and electric circuits. Let us consider two topological disked trees $H \subset H^{\prime}$, with the families of disks $\left\{D_{j}\right\} \subset\left\{D_{i}^{\prime}\right\}$. Then the corresponding aligned arc diagrams are related by the topological arrow: $\mathbf{H}^{\prime} \rightsquigarrow \mathbf{H}$. Let $\mathbf{G}$ and $\mathbf{G}^{\prime}$ be the corresponding trees of complete graphs.

We say that WAD $Y$ is aligned with $H$ if $\operatorname{supp} Y$ is aligned with $H$. Such a diagram induces an unplugged electric circuit $\mathcal{C}_{Y}$ based on $\mathbf{G}$ by letting the conductance of the edge $e \in \mathbf{H}$ be $Y(e)$. The following lemma relates the domination relation between aligned WAD's to the domination relation between the corresponding electric circuits (see §11.6).
Lemma 6.1. Let WAD's $Y$ and $Y^{\prime}$ be aligned with trees $H \subset H^{\prime}$ as above. If $Y^{\prime} \multimap Y$ then $\mathcal{C}_{Y^{\prime}} \multimap \mathcal{C}_{Y}$.
Proof. By definition of domination, there exist edges $\beta_{i} \in \mathbf{H}$ concatenated by paths $\left(\alpha_{i j}\right)_{j}$ in $\mathbf{H}^{\prime}$, and positive numbers $w_{i}, v_{i j}$ such that

$$
\begin{equation*}
Y=\sum_{i} w_{i} \beta_{i}, \quad Y^{\prime} \geqslant \sum_{i} \sum_{j} v_{i j} \alpha_{i j}, \tag{6.2}
\end{equation*}
$$

and for any $i$,

$$
\begin{equation*}
w_{i} \leqslant \bigoplus_{j} v_{i j} . \tag{6.3}
\end{equation*}
$$

Let $B(e)$ be the family of edges $\beta_{i}$ equal to $e$. For any $\beta_{i} \in B(e)$, let us consider an auxiliary electric circuit $\mathcal{C}_{i}^{\prime}$ whose resistors are $\left(\alpha_{i j}\right)_{j}$ with conductances $\left(v_{i j}\right)_{j}$ plugged in series with battery $\partial e$. Then (6.3) translates into an estimate: $w_{i} \leqslant \mathbf{W}\left(\mathcal{C}_{i}^{\prime}\right)$, where $\mathbf{W}\left(\mathcal{C}_{i}^{\prime}\right)=\oplus_{j} v_{i j}$ is the conductance of $\mathcal{C}_{i}^{\prime}$.

Each electric circuit $\mathcal{C}_{i}^{\prime}$ admits a natural projection into the circuit $\mathcal{C}_{Y}^{\prime}$. By Lemma 11.7 and the second inequality of (6.2), $\sum \mathbf{W}\left(\mathcal{C}_{i}\right) \leqslant$ $\mathbf{W}\left(\mathcal{C}_{Y^{\prime}}(e)\right)$. Hence

$$
Y(e)=\sum_{\beta_{i} \in B(e)} w_{i} \leqslant \sum_{\beta_{i} \in B(e)} \mathbf{W}\left(\mathcal{C}_{i}\right) \leqslant \mathbf{W}\left(\mathcal{C}_{Y^{\prime}}(e)\right)
$$

and we are done.
6.2. Dynamical setting. Let us now consider the Hubbard tree $H$ of a superattracting map $F$ with the basin $\mathcal{D}=\cup D_{j}$, where $j \in \mathbb{Z} / p Z$. If $Y$ is a WAD aligned with $H$, we let

$$
Y \mid D_{j}=\sum_{\partial \alpha \ni D_{j}} Y(\alpha)
$$

be the local conductances of the associated electric circuit (see §11.5).
Lemma 6.4. If $Y$ and $Z$ are WAD's aligned with the Hubbard tree $H$ and such that $F^{*} Y \multimap Z$, then for any two disk $D_{j}$ we have:

$$
Z\left|D_{j} \leqslant\left(\operatorname{deg} F \mid D_{j}\right) \cdot Y\right| D_{j+1}
$$

Proof. We have $F^{*} Y\left|D_{j}=\operatorname{deg}\left(F \mid D_{j}\right) \cdot Y\right| D_{j+1}$. By Lemma 6.1, $\mathcal{C}_{F * Y} \multimap \mathcal{C}_{Y}$. Hence $F^{*} Y\left|D_{j} \geqslant Z\right| D_{j}$ by Lemma 11.10. Thus,

$$
Z\left|D_{j} \leqslant\left(\operatorname{deg} F \mid D_{j}\right) \cdot Y\right| D_{j+1}
$$

Lemma 6.5. Let $Y$ and $Z$ be WAD's aligned with the Hubbard trees H. Assume $\left(F^{r p}\right)^{*} Y \multimap Z$ for some $r \in \mathbb{Z}_{+}$and some WAD Z. Then for any edge $\alpha$ of $\mathbf{H}$, we have:

$$
Z(\alpha) \leqslant \frac{1}{2^{r-2}} \max _{j}\left(Y \mid D_{j}\right)
$$

Proof. The WAD $\left(F^{r p}\right)^{*} Y$ is aligned with the tree $H^{r p}=F^{-r p} H \supset H$ (see $\S 4.8$ ). Let $\mathcal{C}^{l}$ be the associated electric circuit. Recall that $\mathcal{C}^{l}(\alpha)$ stands for the restriction of this circuit to the connected component of $\mathbf{G}^{l} \backslash \mathcal{D}$ attached to $\partial \alpha$. By Lemma 6.1, $\mathcal{C}^{l} \multimap \mathcal{C}$, so that,

$$
Z(\alpha) \leqslant \mathbf{W}\left(\mathcal{C}^{l}(\alpha)\right)
$$

Let us consider the path of disks $\left(D_{k}^{r p}\right)_{k}$ connecting $\partial \alpha$ in the tree of graphs $\mathbf{G}^{r p}$ (see Lemma 11.11). By Lemma 11.12

$$
\mathbf{W}\left(\mathcal{C}^{l}(\alpha)\right) \leqslant \bigoplus_{k}\left(F^{r p}\right)^{*} Y \mid D_{k}^{r p}
$$

If $D_{k}^{r p} \in \mathcal{D}^{r p} \backslash \mathcal{D}^{(r-1) p}$ then $F^{r p}$ maps $D_{k}^{r p}$ onto some disk $D_{j}$ with degree at most 2. Hence

$$
\left(F^{r p}\right)^{*} Y\left|D_{k}^{r p} \leqslant 2 Y\right| D_{j}
$$

Since by Lemma 4.11, there are at least $2^{r-1}$ such disks, we have:

$$
\bigoplus_{k}\left(F^{r p}\right)^{*} Y \left\lvert\, D_{k}^{r p} \leqslant \frac{1}{2^{r-2}} \max _{j}\left(Y \mid D_{j}\right) .\right.
$$

Putting the above estimates together, we obtain the desired.
6.3. Loss of the horizontal weight. Recall now the WAD's $\hat{X}^{n}=$ $X^{n}-q_{n}$ from §5.2. By Lemma 4.2, they can be viewed as WAD's of the superattracting model $F$. We will keep the same notation for these diagrams.

Lemma 6.6. Let $p$ be the period of 0 under $F$. Then

$$
\left\|\hat{X}^{10 p}\right\|_{1} \leqslant \frac{1}{4}\|\hat{X}\|_{1} .
$$

Proof. By Corollary 5.5, the WAD's $\hat{X}^{n}$ are aligned with the Hubbard tree $H$ for $n \geqslant 3 p$. By Proposition 5.2 (ii), $\left(F^{6 p}\right)^{*} \hat{X}^{4 p} \multimap \hat{X}^{10 p}$. Hence by Lemma 6.5,

$$
\hat{X}^{10 p}(\alpha) \leqslant \frac{1}{16} \max _{j}\left(\hat{X}^{4 p} \mid D_{j}\right)
$$

By Lemma 6.4, $\left(\hat{X}^{4 p} \mid D_{j}\right) \leqslant 2\left(\hat{X}^{3 p} \mid D_{k}\right)$ for any two disks $D_{j}$ and $D_{k}$. Hence

$$
\max \left(\hat{X}^{4 p} \mid D_{j}\right) \leqslant \frac{2}{p}\left\|\hat{X}^{3 p}\right\|_{1} \leqslant \frac{2}{p}\|\hat{X}\|_{1}
$$

where the last estimate comes from Proposition 5.2(i).
Putting the above estimates together and summing it up over all $\alpha \in \operatorname{supp} X^{10 p}$ (taking into account that $\operatorname{supp} X^{10 p}$ contains at most at most $3 p$ arcs), we obtain:

$$
\left\|\hat{X}^{10 p}\right\|_{1} \leqslant \frac{3}{8}\left\|\hat{X} \mid D_{j}\right\|_{1}
$$

Let us now go back to the original map $f$. The following result shows that after an appropriate restriction of the domain of $f$, there is a definite loss of the horizontal weight of the associated canonical WAD.

Corollary 6.7. Let $f$ be a renormalizable $\psi$-ql map with period $p$. Then there exists $M=M(p)$ such that

$$
\left\|W_{\mathrm{can}}^{h}\left(\mathbf{U}^{10 p} \backslash \mathcal{K}\right)\right\|_{1} \leqslant \frac{1}{2}\left\|W_{\mathrm{can}}^{h}(\mathbf{U} \backslash \mathcal{K})\right\|_{1},
$$

provided $\left\|W_{\text {can }}^{h}(\mathbf{U} \backslash \mathcal{K})\right\|_{1}>M(p)$.
Proof. The last lemma immediately yields

$$
\left\|W_{\mathrm{can}}^{h}\left(\mathbf{U}^{10 p} \backslash \mathcal{K}\right)\right\|_{1} \leqslant \frac{1}{4}\left\|W_{\mathrm{can}}^{h}(\mathbf{U} \backslash \mathcal{K})\right\|_{1}+3 p q_{10 p}
$$

which implies the desired.

## 7. Push-forward Argument

In this section we will derive Theorem 3.1 from the loss of horizontal weight (Corollary 6.7), using some simple observations and the Covering Lemma[KL1].
7.1. Vertical weight from the loss of horizontal. We first observe

Lemma 7.1. For any $n \in \mathbb{Z}_{\geqslant 0}$, we have:

$$
\left\|W_{\mathrm{can}}^{\mathrm{v}+\mathrm{h}}(U \backslash \mathcal{K})\right\|_{1} \leqslant\left\|W_{\mathrm{can}}^{\mathrm{v}+\mathrm{h}}\left(U^{n} \backslash \mathcal{K}\right)\right\|_{1}+C(p) .
$$

Proof. This follows immediately from Theorem 2.17, because $U^{n} \backslash \mathcal{K} \subset$ $U \backslash \mathcal{K}$, and the inclusion is homotopic to a homeomorphism.

Combining Corollary 6.7 with Lemma 7.1, we obtain:
Lemma 7.2. Let $f$ be a renormalizable quadratic-like map with period $p$. Then there exists $M=M(p)$ such that

$$
\left\|W_{\text {can }}^{\mathrm{v}}\left(U^{10 p} \backslash \mathcal{K}\right)\right\|_{1} \geqslant \frac{1}{2}\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}(U \backslash \mathcal{K})\right\|_{1},
$$

provided $\left\|W_{\text {can }}^{h}(U \backslash \mathcal{K})\right\|_{1}>M(p)$.
Our plan in the next two sections is to use the Covering Lemma to "push forward" the $W_{\text {can }}^{\mathrm{v}}\left(U^{10 p} \backslash \mathcal{K}\right)$ to get a lower bound on $W_{\text {can }}^{\mathrm{v}}(U \backslash$ $\mathcal{K})$.
7.2. Covering Lemma. Let us state a suitable version of the Covering Lemma[KL1] in the language of the canonical weighted arc diagram.

Theorem 7.3. Suppose that $U$ and $V$ are simply connected Riemann surfaces, $A=\bigcup_{i=1}^{p} A_{i}$ is a union of $p$ disjoint compact connected full subsets of $U, B=\bigcup_{i=1}^{n} B_{i}$ is the same on $V$, and $f: U \rightarrow V$ is a holomorphic branched cover of degree $D$, such that all the critical values of $f$ lie in $B$, and for each $i, A_{i}$ is a component of $f^{-1} B_{i}$, and the degree
of $\left.f\right|_{A_{i}}$ is d (in the sense that $A_{i}$ and $B_{i}$ have open neighborhoods $A_{i}^{\prime}$ and $B_{i}^{\prime}$ such that $f: A_{i}^{\prime} \rightarrow B_{i}^{\prime}$ is a degree d branched cover). Suppose further that

$$
\begin{equation*}
\left\|W_{\text {can }}^{\mathrm{v}}(U \backslash A)\right\|_{1} \geqslant \epsilon\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}(V \backslash B)\right\|, \tag{7.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|W_{\text {can }}^{\mathrm{v}}(U \backslash A)\right\|_{1} \geqslant M(\epsilon, D, p) \tag{7.5}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left\|W_{\mathrm{can}}^{\mathrm{v}}(U \backslash A)\right\|_{1} \leqslant \frac{12 d^{2}}{\epsilon}\left\|W_{\mathrm{can}}^{\mathrm{v}} V \backslash A\right\|_{1} . \tag{7.6}
\end{equation*}
$$

Before proving this theorem, let us review the notation of [KL1] and related it to $W_{\text {can }}(U \backslash A)$. Accordingly, let $U$ and $A$ be as in the statement of Theorem 7.3. We then define three numbers determined by $U$ and $A$ :

$$
\begin{gather*}
X=\mathcal{W}\left(S, \bigcup_{j=1}^{N} A_{j}\right) \\
Y=\sum_{j=1}^{N} \mathcal{W}\left(S, A_{j}\right),  \tag{7.7}\\
Z=\sum_{j=1}^{N} \mathcal{W}\left(S \backslash \bigcup_{k \neq j} A_{k}, A_{j}\right) .
\end{gather*}
$$

It is easy to see that $X \leqslant Y \leqslant Z$. We further claim that

$$
\begin{equation*}
W_{\text {can }}^{\mathrm{v}}(U \backslash A) \leqslant X \leqslant W_{\text {can }}^{\mathrm{v}}(U \backslash A)+4 p \tag{7.8}
\end{equation*}
$$

and

$$
\begin{equation*}
W_{\mathrm{can}}^{2 \mathrm{~h}+\mathrm{v}}(U \backslash A) \leqslant Z \leqslant W_{\mathrm{can}}^{2 \mathrm{~h}+\mathrm{v}}(U \backslash A)+4 p^{2} \tag{7.9}
\end{equation*}
$$

To see the lower bound for $X$ in (7.8), we need only observe that every vertical leaf in the canonical lamination for $U \backslash A$ contributes to the defining width for $X$, so $X$ is at least the total width of this vertical lamination. To see the upper bound, we consider the extremal foliation $\mathcal{F}_{X}$ for the defining path family for $X$. Then the leaves of $\mathcal{F}_{X}$ can lie in at most $2 p$ homotopy classes $\alpha_{i}$, and

$$
W\left(\left.\mathcal{F}_{X}\right|_{\alpha_{i}}\right) \leqslant W_{\text {can }}^{U \backslash A}\left(\alpha_{i}\right)+2
$$

by maximality. The derivation of (7.9) is similar and left as an exercise for the reader.

We can now reduce Theorem 7.3 to the version of the Quasi-Invariance Law stated in Section 3.1.1 of [KL1].

Proof of Theorem 7.3. We wish to think of the setting of our Theorem as an instance of the setting of the QI Law stated in Section 3 of [KL1], and specialized in Section 3.1.1. So we apply the latter, letting $U$ and $V$ in [KL1] be our $U$ and $V$, letting $\Lambda_{j}^{\prime}=\Lambda_{j}$ be our $A_{j}$, and letting $B_{j}^{\prime}=B_{j}$ be our $B_{j}$. Then the $Y_{U}$ in Section 3.1.1 of [KL1] is our $Y(U, A)$; we will use only that it is at least $X(U, A)$. Because the critical values of $f$ lie in $B$, the $Z_{V}=Z_{V}\left\{B_{j}, B_{j}^{\prime}, C V\right\}$ in Section 3.1.1 is simply our $Z(V, B)$, and of course $X_{V}$ in [KL1] is our $X(V, B)$.

The context and direction of the two inequalities in Section 3.1.1 of [KL1] are such that we can replace $Y_{U}$ with $X_{U}$ in them. We can then derive, assuming $\left\|W^{\mathrm{v}}(U \backslash A)\right\| \geqslant 4 p^{\epsilon}$,

$$
\begin{align*}
Z_{V} & \leqslant\left\|W_{\mathrm{can}}^{2 \mathrm{~h}+\mathrm{v}}(V \backslash B)\right\|_{1}+4 p^{2} & & \text { by }(7.9) \\
& \leqslant 2\left\|W_{\mathrm{can}}^{\mathrm{h}+\mathrm{v}}(V \backslash B)\right\|_{1}+4 p^{2} & & \\
& \leqslant \frac{3}{\epsilon}\left\|W_{\mathrm{can}}^{\mathrm{v}}(U \backslash A)\right\|_{1} & & \text { by }(7.4)  \tag{7.4}\\
& \leqslant \frac{3}{\epsilon} X_{U} & & \text { by }(7.8) . \tag{7.8}
\end{align*}
$$

This is the separation assumption (with $\xi:=3 / \epsilon$ ) in Section 3.1.1 of [KL1]; we can then conclude, assuming (7.5) (which implies $X_{U} \geqslant$ $K(\xi, p, D)$ ),

$$
X_{U} \leqslant \frac{6 d^{2}}{\epsilon} X_{V}
$$

which implies

$$
\left\|W_{\mathrm{can}}^{\mathrm{v}}(U \backslash A)\right\|_{1} \leqslant \frac{6 d^{2}}{\epsilon}\left(\left\|W_{\mathrm{can}}^{\mathrm{v}}(V \backslash B)\right\|_{1}+4 p\right)
$$

Assuming then that $W_{\text {can }}^{\mathrm{v}}(U \backslash A) \geqslant 48 d^{2} p / \epsilon$, we obtain (7.6).
7.3. Vertical part has a definite weight. We can now prove Theorem 3.1:
Proof of Theorem 3.1. If $\left\|W_{\text {can }}^{\mathrm{v}}(U \backslash \mathcal{K})\right\|_{1} \geqslant\left\|W_{\text {can }}^{\mathrm{h}}(U \backslash \mathcal{K})\right\|_{1}$, there is nothing to prove. So, we assume that the opposite inequality holds, and then

$$
\left\|W_{\mathrm{can}}^{\mathrm{h}}(U \backslash \mathcal{K})\right\|_{1}>\frac{1}{2}\left\|W_{\mathrm{can}}^{\mathrm{v}+\mathrm{h}}(U \backslash \mathcal{K})\right\|_{1}
$$

Then taking $M=M(p)$ as in Corollary 7.2, we conclude that

$$
\begin{equation*}
\left\|W_{\text {can }}^{\mathrm{v}}\left(U^{10 p} \backslash \mathcal{K}\right)\right\|_{1} \geqslant \frac{1}{2}\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}(U \backslash \mathcal{K})\right\|_{1} \tag{7.10}
\end{equation*}
$$

provided $\left\|W_{\text {can }}(U \backslash \mathcal{K})\right\|_{1}>2 M$.
We are now ready to apply Theorem 7.3 to the map $f^{10 p}: U^{10 p} \rightarrow U$, where we let $A=B=\mathcal{K}$. We observe that all the hypotheses are
satisfied, with $d=2^{10}, D=2^{10 p}$, and $\epsilon=1 / 2$. Moreover, by Lemma 7.1 and our assumption, we have

$$
\left\|W_{\mathrm{can}}^{\mathrm{h}+\mathrm{v}}\left(U^{10 p} \backslash \mathcal{K}\right)\right\| \geqslant\left\|W_{\mathrm{can}}^{\mathrm{h}+\mathrm{v}}(U \backslash \mathcal{K})\right\|-C(p) \geqslant M^{\prime}(p),
$$

so we obtain, as our form of (7.6),

$$
\left\|W_{\text {can }}^{\mathrm{v}}(U \backslash \mathcal{K})\right\| \geqslant\left(24 \cdot 2^{20}\right)^{-1}\left\|W_{\text {can }}^{\mathrm{v}}\left(U^{10 p} \backslash \mathcal{K}\right)\right\|
$$

The Lemma then follows from this and Lemma 7.2.

## 8. Pseudo-quadratic-Like maps and canonical RENORMALIZATION

Definition 8.1. A pseudo-polynomial-like map $^{8}$ (or $\psi$-polynomial-like map) of degree $d$ is a pair of holomorphic maps $(i, f): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$ and an FJ-set $K \subset \mathbf{U}$ such that
(1) $\mathbf{U}^{\prime}$ and $\mathbf{U}$ are disks,
(2) $i$ is an immersion,
(3) $f$ is a degree $d$ branched cover,
(4) $i^{-1}(K)=f^{-1}(K)$, and, denoting this set by $K^{\prime}$,
(5) $K^{\prime}$ is connected.

We will call $K$ the filled-Julia set of $(i, f)$. Note that $K^{\prime}$ is connected if and only if all the branch values of $f$ lie in $K$. Also note that, because $f$ is proper, $K^{\prime}$ is an FJ-set. We will see that $K \subset \mathbf{U}$ is uniquely determined by $(i, f): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$. We will also see that these conditions imply that $i: K^{\prime} \rightarrow K$ is a bijection.
8.1. Finding a polynomial-like restriction. We can now prove that every $\psi$-polynomial-like map $F$ has a polynomial-like restriction, with modulus controlled by the modulus for $F$.

Theorem 8.1. Let $F=(i, f): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$ be a $\psi$-polynomial-like map of degree $d$ with filled Julia set $K$. Then we can find $U \subset \mathbf{U}$ open, with $K \subset U$ such that, letting $U^{\prime} \equiv f^{-1} U \subset \mathbf{U}^{\prime}$,

- $f: U^{\prime} \rightarrow U$ is a degree d holomorphic cover,
- $\left.i\right|_{U^{\prime}}$ is an embedding,
- $i\left(U^{\prime}\right) \subset U$, and
- $\bmod \left(U \backslash i\left(U^{\prime}\right)\right) \geqslant \mu(d, \bmod (\mathbf{U} \backslash K))$, where $\mu(d, \nu)>0$ for $\nu>0$.

[^5]Proof. Let us show that $i$ is an embedding in a neighborhood of $K^{\prime}$. To this end let us consider the annuli $\mathbf{U} \backslash K$ and $\mathbf{U}^{\prime} \backslash K^{\prime}$, and let us uniformize them by the round annuli:

$$
\phi: \mathbb{A}(1, r) \rightarrow \mathbf{U} \backslash K, \quad \phi^{\prime}: \mathbb{A}\left(1, r^{\prime}\right) \rightarrow \mathbf{U}^{\prime} \backslash K^{\prime}, \quad \text { where } r^{\prime}=r^{1 / d} .
$$

Let $I=\phi^{-1} \circ i \circ \phi^{\prime}: \mathbb{A}\left(1, r^{\prime}\right) \rightarrow \mathbb{A}(1, r)$. Since $i\left(K^{\prime}\right) \subset K$, the map $I$ is proper near the unit circle $\mathbb{T}$. By the Reflection Principle, it admits an analytic extension to a map $\mathbb{A}\left(1 / r^{\prime}, r^{\prime}\right) \rightarrow \mathbb{A}(1 / r, r)$ that restricts to a covering $\mathbb{T} \rightarrow \mathbb{T}$ of some degree $k>0$. (We will use the same notation, $I$, for this extension).

Let us consider the geodesic $\omega^{\prime}=\mathbb{T}_{\sqrt{r^{\prime}}}$ in $\mathbb{A}\left(1, r^{\prime}\right)$, and let $\omega=$ $I\left(\omega^{\prime}\right)$. Since for given $r, r^{\prime}>0$, the family of homotopically nontrivial holomorphic immersions $\mathbb{A}\left(1, r^{\prime}\right) \rightarrow \mathbb{A}(1, r)$ is compact, the (Euclidean) distance from $\omega$ to $\mathbb{T}$ is greater than $\rho=\rho\left(r^{\prime}, r\right)>0$. Let $\Lambda=\mathbb{A}(1 / \rho, \rho)$ and let $\Lambda^{\prime}$ be the component of $I^{-1}(\Lambda)$ containing $\mathbb{T}$. By the Argument Principle, the map $I: \Lambda^{\prime} \rightarrow \Lambda$ is a covering of degree $k$.

In fact, $k=1$. Indeed, let $\gamma$ and $\gamma^{\prime}$ be the outer boundaries of $\Lambda$ and $\Lambda^{\prime}$ respectively. Let $\Gamma=\phi(\gamma), \Gamma^{\prime}=\phi^{\prime}\left(\gamma^{\prime}\right)$, and let $\Omega$ and $\Omega^{\prime}$ be the disks bounded by $\Gamma$ and $\Gamma^{\prime}$ respectively. Since $I: \gamma^{\prime} \rightarrow \gamma$ is a covering of degree $k$, so is $i: \Gamma^{\prime} \rightarrow \Gamma$. Hence $i: \Omega^{\prime} \rightarrow \Omega$ is a branched covering of degree $k$. If $k>1, i$ would necessarily have a critical point, contradicting the assumption that it is an immersion.

Furthermore, let $F=\phi^{-1} \circ f \circ \phi^{\prime}: \mathbb{A}\left(1, r^{\prime}\right) \rightarrow \mathbb{A}(1, r)$. Since it is a covering map of degree $d, F(z)=z^{d}$. Hence $F^{-1}(\mathbb{A}(1, \rho))=\mathbb{A}\left(1, \rho^{\prime}\right)$, where $\rho^{\prime}=\rho^{1 / d}$. Notice that the hyperbolic distance from $\mathbb{T}_{\rho}$ to $\mathbb{T}$ in $\mathbb{A}(1 / r, r)$ is equal to the hyperbolic distance from $\mathbb{T}_{\rho^{\prime}}$ to $\mathbb{T}$ in $\mathbb{A}\left(1 / r^{\prime}, r^{\prime}\right)$. Since the map $I: \mathbb{A}\left(1 / r^{\prime}, r^{\prime}\right) \rightarrow \mathbb{A}(1 / r, r)$ contracts the respective hyperbolic metrics by a uniform factor, we have: $I\left(\mathbb{T}_{\rho^{\prime}}\right) \subset \mathbb{A}(1, \lambda \rho)$, where $\lambda=\lambda(r, d, \rho)<1$. (Actually $\lambda=\lambda(r, d)$ because $\rho=\rho\left(r^{1 / d}, r\right)$.) Putting all these together, we conclude that

$$
W^{\prime} \equiv I \circ F^{-1}(\mathbb{A}(1, \rho)) \subset \mathbb{A}(1, \lambda \rho) \equiv W .
$$

and moreover, the map $F \circ I^{-1}: W^{\prime} \rightarrow W$ is a covering of degree $d$.
Letting

$$
U^{\prime}=K \cup \phi\left(W^{\prime}\right), \quad U=K \cup \phi(W)
$$

we see that the map $f \circ i^{-1}: U^{\prime} \rightarrow U$ is polynomial-like of degree $d$ and $\bmod \left(U \backslash U^{\prime}\right) \geqslant \frac{1}{2 \pi} \log \lambda$, where $\lambda$ depends only on $d$ and $\bmod (\mathbf{U} \backslash$ $K)$.

In the proof we have that $i: U^{\prime} \rightarrow i(U)$ is a homemorphism, and hence $i: K^{\prime} \rightarrow K$ is a bijection.

The above estimates with poincare metric also show that if $\hat{K} \subset \mathbf{U}$ is compact, and $K \subset \hat{K}$, and $\left(\hat{K}_{i}\right)_{i=0}^{\infty}$ is defined by $\hat{K}_{0}=\hat{K}$ and $\hat{K}_{i+1}=$ $i\left(f^{-1}(\hat{K})\right)$, then $K=\bigcap_{i=0}^{\infty} \hat{K}_{i}$. Therefore $K$ is uniquely determined by $F$.
8.2. Iterating a $\psi$-polynomial-like map. It is not immediately obvious how to iterate a $\psi$-polynomial-like map. We will briefly review a few simple constructions in the categories of topological spaces and Riemann surfaces, and then use this to construct the iterations.
8.2.1. A few simple constructions. Given topological spaces $X, Y, Z$ and continuous maps $f: X \rightarrow Z$ and $g: Y \rightarrow Z$, we can form the pullback $E=\{(x, y) \in X \times Y \mid f(x)=g(y)\}$, along with maps $\pi_{X}: E \rightarrow X$ and $\pi_{Y}: E \rightarrow Y$. With this notation, we observe that $\pi_{Y}$ can inherit certain properties from $f$ (and $\pi_{X}$ from $g$ ).

Lemma 8.2. We have
(1) $f \circ \pi_{X}=g \circ \pi_{Y}$.
(2) If $f$ is proper, then $\pi_{Y}$ is proper.
(3) If $g$ is a local homeomorphism, then $\pi_{X}$ is a local homeomorphism.
(4) If $g$ is injective (or surjective, or bijective), then $\pi_{X}$ is injective (or surjective, or bijective).

Proof. Statement (1) is immediate.
For Statement (2), suppose $B \subset Y$ is compact. Then

$$
\begin{aligned}
\pi_{Y}^{-1}(B) & =\{(x, y) \in X \times B \mid g(x)=h(y)\} \\
& =\left\{(x, y) \in f^{-1}(g(X)) \times B \mid g(x)=h(y)\right\}
\end{aligned}
$$

which is a closed subset of a compact set.
For Statement (3), take $(x, y) \in E$. Then we can find $V \subset Y$ such that $y \in V$ and $\left.g\right|_{V}: V \rightarrow g(V)$ is a homeomorphism. Let $\hat{V}=$ $\pi_{Y}^{-1} V$. Then $\left.\pi_{X}\right|_{\hat{V}}$ is a homeomorphism, because its inverse is $x \mapsto$ $\left(x,(g \mid V)^{-1}(f(x))\right)$.

Statement (4) is straightforward (and belongs to the category of sets) and is left to the reader.

It will also help to have a few similar observations about composition.
Lemma 8.3. Suppose that $U, V, W$ are Riemann surfaces and $g: U \rightarrow$ $V$ is continuous and $f: V \rightarrow W$ are holomorphic. Then
(1) if $f$ and $f \circ g$ are holomorphic immersions, then $g$ is a holomorphic immersion, and
(2) if $f \circ g$ is proper, then $g$ is proper.

Proof. For (1), we can locally write $g=(f \circ g)^{-1} \circ f$ for a suitable local branch of $(f \circ g)^{-1}$.

For (2), we suppose $A \subset V$ is compact. Then

$$
g^{-1}(A) \subset g^{-1}\left(f^{-1}(f(A))=(f \circ g)^{-1}(f(A))\right.
$$

is compact.
We can now prove
Lemma 8.4. Suppose that $U_{f}, U_{i}$, and $U_{i}$ are Riemann surfaces, and $f: U_{f} \rightarrow U$ is holomorphic and proper, and $i: U_{i} \rightarrow U$ is a holomorphic immersion. Then we can find a Riemann surface $E$ and $\hat{i}: E \rightarrow U_{f}$ and $\hat{f}: E \rightarrow U_{i}$ such that
(1) $\hat{f} \circ i=\hat{i} \circ f$,
(2) $\hat{i}$ is holomorphic immersion,
(3) $\hat{f}$ is holomorphic and proper, and $\operatorname{deg} \hat{f}=\operatorname{deg} f$.

Moreover, these objects are unique, in that if $E^{\prime}, \hat{f}^{\prime}$ and $\hat{i}^{\prime}$ satisfy Condition 1-3 above, then there is a holomorphic isomorphism $h: E^{\prime} \rightarrow E$ such that $\hat{i}^{\prime}=\hat{i} \circ h$ and $\hat{f}^{\prime}=\hat{f} \circ h$.

Proof. We can construct $E$ as a topological space as the pullback described above; we relabel $\pi_{U_{i}}$ and $\pi_{U_{f}}$ as $\hat{f}$ and $\hat{i}$. Then Condition (1) holds, and $\hat{f}$ is proper, and $\hat{i}$ is a local homeomorphism, by Lemma 8.2. We can then pull back the Riemann surface structure on $U_{f}$ to $E$ by $\hat{i}$ to make $\hat{i}$ holomorphic. Then $\hat{f}$ is holomorphic, because we can write it locally as $\hat{i} \circ f \circ i^{-1}$, for a suitable local branch of $i^{-1}$.

For a generic $x \in U_{i}$, the equation $f(y)=i(x)$ has $\operatorname{deg} f$ solutions in $y$. Therefore $\operatorname{deg} \hat{f}=\operatorname{deg} f$.

Now suppose that $E^{\prime}$, $\hat{f}^{\prime}$, and $\hat{i}^{\prime}$ satisfy Conditions 1-3 above. Then there is a continuous map $h: E^{\prime} \rightarrow E$ such that $\hat{i}^{\prime}=\hat{i} \circ h$ and $\hat{f}^{\prime}=\hat{f} \circ h$. We can apply Lemma 8.3 to see that $h$ is a holomorphic immersion, and $h$ is proper. Therefore $h$ is a finite-sheeted covering, and $\operatorname{deg} \hat{f}^{\prime}=$ $(\operatorname{deg} \hat{f})(\operatorname{deg} h)$, so $\operatorname{deg} h=1$, and $h$ is a holomorphic isomorphism.
Remark 8.1. While the pullback of $(f, g):(U, V) \rightarrow W$ always exists, we have chosen to construct it only in a simple case where we not only have universality, but a strong form of uniqueness.
8.2.2. Iteration. We can now construct our iterations of a $\psi$-polynomiallike map.

Theorem 8.5. Let $(i, f):\left(\mathbf{U}^{\prime}, K^{\prime}\right) \rightarrow(\mathbf{U}, K)$ be a degree d $\psi$-polynomiallike map. Then there exists an essentially unique degree d $\psi$-polynomial-like-map $\left(i_{1}, f_{1}\right):\left(\mathbf{U}^{\prime \prime}, K^{\prime \prime}\right) \rightarrow\left(\mathbf{U}^{\prime}, K^{\prime}\right)$ such that $i \circ f_{1}=f \circ i_{1}$.

Proof. By Lemma 8.4. there is a Riemann surface $\mathbf{U}^{\prime \prime}$ with a holomorphic immersion $i_{1}: \mathbf{U}^{\prime \prime} \rightarrow \mathbf{U}^{\prime}$ and a degree $d$ holomorphic branched cover $f_{1}: \mathbf{U}^{\prime \prime} \rightarrow \mathbf{U}^{\prime}$ such that $i \circ f_{1}=f \circ i_{1}$. (Lemma 8.4 also says that these data are essentially unique.) Moreover, for $z \in \mathbf{U}^{\prime \prime}$,
$i_{1}(z) \in K^{\prime} \Longleftrightarrow\left(f \circ i_{1}\right)(z) \in K \Longleftrightarrow\left(i \circ f_{1}\right)(z) \in K \Longleftrightarrow f_{1}(z) \in K^{\prime}$.
So $i_{1}^{-1}\left(K^{\prime}\right)=f_{1}^{-1}\left(K^{\prime}\right)$; we denote this set by $K^{\prime \prime}$. We observe that $\left(f_{1}, i_{1}\right): K^{\prime \prime} \rightarrow K^{\prime}$ is the topological pullback of $i: K^{\prime} \rightarrow K$ and $f: K^{\prime} \rightarrow$ $K$, and hence $i_{1}$ is bijective by Lemma 8.2.

We claim that $\mathbf{U}^{\prime \prime}$ is a disk. To see this, we observe that $z \in \mathbf{U}^{\prime \prime}$ is a critical point of $f_{1}$ if and only if it is a critical point of $i \circ f_{1}=f \circ i_{1}$, and this occurs exactly when $i_{1}(z)$ is a critical point of $f$. Now $f$ has exactly $d-1$ critical points (by Riemann-Hurwitz) and they all lie in $K^{\prime}$, and $i_{1}: K^{\prime \prime} \rightarrow K^{\prime}$ is a bijection, so $f_{1}$ has exactly $d-1$ critical points and hence $\mathbf{U}^{\prime \prime}$ is a disk (by Riemann-Hurwitz).

The branch values of $f_{1}$ map by $i$ to branch values of $f$, so all these branch values lie in $K^{\prime}$. Therefore $K^{\prime \prime}$ is connected, and we have verified that $\left(i_{1}, f_{1}\right):\left(\mathbf{U}^{\prime \prime}, K^{\prime \prime}\right) \rightarrow\left(\mathbf{U}^{\prime}, K^{\prime}\right)$ is a degree $d \psi$-polynomial-like map.

Let us relabel $(i, f)$ as $\left(i_{0}, f_{0}\right):\left(\mathbf{U}^{1}, K^{1}\right) \rightarrow\left(\mathbf{U}^{0}, K^{0}\right)$. We can then inductively apply Theorem 8.5 to construct degree $d \psi$-polynomial-like maps $\left(i_{n}, f_{n}\right):\left(\mathbf{U}^{n+1}, K^{n+1}\right) \rightarrow\left(\mathbf{U}^{n}, K^{n}\right)$ such that $i_{n} \circ f_{n+1}=f_{n} \circ i_{n+1}$ for all $n$. By a very slight abuse of notation ${ }^{9}$, we will often just denote $i_{n}, f_{n}$, and $K_{n}$ as $i, f$, and $K$. This then allows us to refer to the iterations $i^{n}: \mathbf{U}^{m+n} \rightarrow \mathbf{U}^{m}$ and $f^{n}: \mathbf{U}^{m+n} \rightarrow \mathbf{U}^{m}$ (where we usually take $m=0$ ).
8.3. Canonical Renormalization. We now describe a form of renormalization in the category of $\psi$-quadratic-like maps that has the crucial property of preserving certain moduli and hyperbolic lengths.
8.3.1. Localization. Suppose that $S$ is a Riemann surface, and $A \subset S$ is an FJ-set. The localization of $S$ at $A$ is a disk $X$ with a holomorphic immersion $\pi: X \rightarrow S$ such that $\pi: \pi^{-1}(A) \rightarrow A$ is a bijection and $\pi: X \backslash$ $\pi^{-1}(A) \rightarrow S \backslash A$ is a covering.

To construct $X$ (and $\pi$ ), we find a neighborhood $U$ of $A$ that is a disk, and take the cover $\pi^{\prime}: X^{\prime} \rightarrow S \backslash A$ corresponding to the fundamental group of the annulus $U \backslash A$. Then there is a natural lift of $U \backslash A$ to $X^{\prime}$; we let $X$ be the quotient of the disjoint union of $U$ and $X^{\prime}$ by identifying points in $U \backslash A$ with their image by this lift. It is clear that $X$ is a Riemann surface, and it is a disk because it is the identification

[^6]of a disk with an annulus along a sub-annulus that is a deformation retract. We let $\pi(x)=\pi^{\prime}(x)$ for $x \in X^{\prime}$, and $\pi(x)=x$ for $x \in U$. As a general rule of notation, we will let $\hat{A}=\pi^{-1}(A)$.

Now suppose that $S^{\prime}$ and $S$ are Riemann surfaces, $A^{\prime} \subset S^{\prime}$ and $A \subset S$ are FJ-sets, and $g:\left(S^{\prime}, A^{\prime}\right) \rightarrow(S, A)$ is holomorphic map such that $A^{\prime}$ is a component of $A$; we let $A^{\prime \prime}=g^{-1}(A) \backslash A^{\prime}$, and $S^{\prime \prime}=S^{\prime} \backslash A^{\prime \prime}$. We let $\pi^{\prime}: X^{\prime} \rightarrow S \backslash A^{\prime \prime}$ be localizations of $S^{\prime \prime}$ at $A^{\prime}$, and $\pi: X \rightarrow S$ and $S$ be the localization of $S$ at $A$. Then we leave it to the reader to verify that there is a unique map ("lift") $\hat{g}: X^{\prime} \rightarrow X$ such that $g \circ \pi=\pi \circ \hat{g}$. We will always have $\hat{g}^{-1}(\hat{A})=\hat{A}^{\prime}$. If $g$ is an immersion, then $\hat{g}$ is an immersion. The conditions on $g, A^{\prime}, A$ imply that $g: A^{\prime} \rightarrow A$ is a branched cover of finite degree $d_{A^{\prime}}$, in the sense that almost every point of $A$ has $d_{A^{\prime}}$ preimages in $A^{\prime}$, and there are simply connected neighborhoods $U^{\prime} \supset A^{\prime}$ and $U \supset A$ such that $g: U^{\prime} \rightarrow U$ is a degree $d_{A^{\prime}}$ branched cover. We then observe that if $g: S^{\prime} \rightarrow S$ is a branched cover ${ }^{10}$ branched only over $A$, then $\hat{g}: X^{\prime} \rightarrow X$ is a branched cover branched only over $A$, and hence has degree $d_{A^{\prime}}$. This is indeed the main idea of localization, in that we can localize a map of much higher degree to get a map of a lower degree, equal to a local degree.
8.3.2. Dynamical Localization. Now let us suppose that

- $(i, f):\left(\mathbf{U}^{\prime}, K^{\prime}\right) \rightarrow(\mathbf{U}, K)$ is a $\psi$-polynomial-like map of degree $d$,
- $A \subset K$ is an FJ-set (and let $A^{\prime}=i^{-1}(A)$ ),
- $B \subset K^{\prime}$ is connected and disjoint from $A$ (and let $B^{\prime}=i^{-1}(B)$ ),
- $A^{\prime}$ is a component of $f^{-1}(A)$,
- $B^{\prime} \subset f^{-1}(B)$, and
- all the branch values of $f$ lie in $A \cup B$.

We observe that $i\left(\mathbf{U}^{\prime} \backslash f^{-1}\left(B^{\prime}\right)\right) \subset \mathbf{U} \backslash B$, and $f: \mathbf{U}^{\prime} \backslash f^{-1}\left(B^{\prime}\right) \rightarrow \mathbf{U} \backslash B$ is a branched cover branched only over $A$. We let $\pi: X \rightarrow \mathbf{U} \backslash B$ be the localization of $\mathbf{U} \backslash B$ with respect to $A$, and $\pi^{\prime}: X^{\prime} \rightarrow \mathbf{U}^{\prime} \backslash f^{-1}\left(B^{\prime}\right)$ be the localization of $\mathbf{U}^{\prime} \backslash f^{-1}\left(B^{\prime}\right)$ with respect to $A^{\prime}$. Then, as observed above, we can lift $i$ and $f$ to $(\hat{i}, \hat{f}): X^{\prime} \rightarrow X$, and $\hat{i}$ is an immersion and $\hat{f}$ is a branched cover, with degree equal to the local degree $d_{A^{\prime}}$ of $f$ at $A^{\prime}$. So we have shown that $(\hat{i}, \hat{f}):\left(X^{\prime}, A^{\prime}\right) \rightarrow(X, A)$ is a $\psi$-polynomiallike map, of degree $d_{A^{\prime}}$. We call it the dynamical renormalization of $f$ at $A$ with the removal of $B$. Moreover, since $\pi: X \backslash \hat{A} \rightarrow \mathbf{U} \backslash B \backslash A$ is a covering, the poincare length of the core curve in $X \backslash \hat{A}$ is equal

[^7]to the poincare length of the closed geodesic in $\mathbf{U} \backslash(A \cup B)$ that goes around $A$.
8.3.3. Canonical Renormalization. Now suppose that $(i, f):\left(\mathbf{U}^{1}, K^{\prime}\right) \rightarrow$ $\left(\mathbf{U}^{0}, K\right)$ is a degree $2 \psi$-polynomial-like map (which we will call a $\psi$ -quadratic-like map). Suppose that $f$ has a quadratic-like restriction that is primitively renormalizable of period $p$. Then as before we let $K_{p} \subset K$ be the central small Julia set of the renormalization, and $\mathcal{K}$ be the union of small Julia sets. We let $K_{p}^{\prime}=i^{-1}\left(K_{p}\right)$ and so forth. We let $\mathcal{K}^{+}=\mathcal{K} \backslash \mathcal{K}_{p}$. We observe that $K_{p}, \mathcal{K}^{+}$satisfy all the conditions for $A, B$ in the previous section, when we replace $(i, f)$ with its iteration $\left(i^{p}, f^{p}\right): \mathbf{U}^{p} \rightarrow \mathbf{U}^{0}$. So we then let $(\hat{i}, \hat{f}):\left(X^{\prime}, \hat{K}_{p}^{\prime}\right) \rightarrow\left(X, \hat{K}_{p}\right)$ be the dynamical renormalization at $K_{p}$ with the removal of $\mathcal{K}_{p}$. We call this the canonical renormalization of $f$ (or really $(i, f)$ ) of period $p$.

As we did with polynomial-like maps, if $(i, f):\left(\mathbf{U}^{\prime}, K^{\prime}\right) \rightarrow(\mathbf{U}, K)$ is a $\psi$-quadratic-like map, then we let $\gamma(f)$ be the core curve of $\mathbf{U} \backslash K$. Moreover, if $f$ is $p$-renormalizable, we let $\gamma_{p}$ be the closed geodesic around $K_{p}$ in $\mathbf{U} \backslash \mathcal{K}_{p}$. We can now prove the a general version of the two "crucial properties" that we stated in Section 1.
Theorem 8.6. Suppose that $(i, f): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$ is a $\psi$-quadratic-like map that is primitively renormalizable with periods $p$ and also $p q$ (with $p, q>$ 1). Let $(j, g): \mathbf{V}^{\prime} \rightarrow \mathbf{V}$ be the canonical renormalization of $f$ of period p. Then
(1) $|\gamma(g)|=\left|\gamma_{p}(f)\right|$, and
(2) $\left|\gamma_{q}(g)\right| \geqslant\left|\gamma_{p q}(f)\right|$.

Proof. Statement (1) has already been observed in the definition of dynamical localization.

To prove (2), we consider the map $\pi: \mathbf{V} \rightarrow \mathbf{U}$ that maps $K_{g}$ to $\mathcal{K}_{p}(0)$ and $\mathbf{V} \backslash K_{g}$ to $\mathbf{U} \backslash \mathcal{K}_{p}$. We claim that $\pi\left(\mathbf{V} \backslash \mathcal{K}_{q}(g)\right) \subseteq \mathbf{U} \backslash \mathcal{K}_{p q}$; this follows because $\pi\left(\mathbf{V} \backslash K_{g}\right) \subseteq \mathbf{U} \backslash \mathcal{K}_{p} \subseteq \mathbf{U} \backslash \mathcal{K}_{p q}$ and $\pi\left(K_{g} \backslash \mathcal{K}_{q}(g)\right)=$ $\mathcal{K}_{p}(0) \backslash \mathcal{K}_{p q}$. Statement (2) then follows from the Schwarz Lemma.

## 9. A priori AND BEAU BOUNDS FOR BOUNDED-PRIMITIVE RENORMALIZATION

Let us first prove the generalization of Theorem 3.1 for $\psi$-quadraticlike maps:
Theorem 9.1. Let $(i, f): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$ be a primitively renormalizable $\psi$ -quadratic-like-map with period p, and let $\mathcal{K}$ be its cycle of little Julia sets. Then there exists $M=M(p)$ such that: If $\left\|W_{\text {can }}^{\mathrm{v}+\mathrm{h}}(\mathbf{U} \backslash \mathcal{K})\right\|_{1}>M$ then

$$
\left\|W_{\text {can }}^{\mathrm{v}}(\mathbf{U} \backslash \mathcal{K})\right\|_{1} \geqslant C^{-1}\left\|W_{\text {can }}^{\mathrm{h}}(\mathbf{U} \backslash \mathcal{K})\right\|_{1},
$$

with an absolute constant $C>0$.
Proof. We let $\mathbf{U}^{0}=\mathbf{U}, \mathbf{U}^{1}=\mathbf{U}^{\prime}, f_{0}=f$, and $i_{0}=i$. As described in Section 8, we can form disks $\mathbf{U}^{n}$ and $\psi$-quadratic-like maps $(i, f): \mathbf{U}^{n+1} \rightarrow \mathbf{U}^{n}$, such that all these $\psi$-quadratic-like maps restrict to the same quadratic-like germ. In particular, we the filled-Julia set $K^{n} \subset U^{n}$ and also the cycle $\mathcal{K}^{n} \subset K^{n}$ of small Julia sets; by a slight abuse of notation, we will just denote these by $K$ and $\mathcal{K}$. We then have (with the slight abuse of notation described at the end of Section 8.1),
(1) a covering $f: \mathbf{U}^{n+1} \backslash f^{-1} \mathcal{K} \rightarrow \mathbf{U}^{n} \backslash \mathcal{K}$,
(2) an inclusion $\mathbf{U}^{n+1} \backslash f^{-1} \mathcal{K} \hookrightarrow \mathbf{U}^{n+1} \backslash \mathcal{K}$, and
(3) an immersion $i: \mathbf{U}^{n+1} \backslash \mathcal{K} \rightarrow \mathbf{U}^{n} \backslash \mathcal{K}$.

The one significant change from the quadratic-like to the $\psi$-quadraticlike setting is the third map, which was previously an embedding, and is now an immersion.

Now we need only go through Sections 5, 6, and 7 and see what needs to be changed. The main observation for Sections 5 and 6 is that the results of these sections depend only on properties (i) and (ii) that are stated in Propostion 5.2; they use no other facts about the geometry of $f$. We therefore need only verify these two properties in our $\psi$-quadratic-like setting. Property (ii) follows just as before, because we use only the covering $f: \mathbf{U}^{n+1} \backslash f^{-1} \mathcal{K} \rightarrow \mathbf{U}^{n} \backslash \mathcal{K}$ and the inclusion $\mathbf{U}^{n+1} \backslash f^{-1} \mathcal{K} \hookrightarrow \mathbf{U}^{n+1} \backslash \mathcal{K}$. For Property (i) we need only observe that the immersion $i: \mathbf{U}^{n+1} \backslash \mathcal{K} \rightarrow \mathbf{U}^{n} \backslash \mathcal{K}$ is still proper on the ends corresponding to $\mathcal{K}$, and Lemma 2.5 applies to immersions as well as embeddings.

Finally, we consider the results of Section 7. For Lemma 7.1, we can replace the inclusion $U^{n} \backslash \mathcal{K} \subset U \backslash \mathcal{K}$ with the immersion $i_{n}: \mathbf{U}^{n} \backslash \mathcal{K} \rightarrow$ $\mathrm{U} \backslash \mathcal{K}$; this is still homotopic to a homeomorphism, so we can apply Theorem 2.17 as before. In a similar vein, examining the proof of Theorem 3.1 in Section 7.3, we find that we need only replace $f^{n}: U^{n} \rightarrow$ $U$ with $f^{n}: \mathbf{U}^{n} \rightarrow \mathbf{U}$, where now this $f^{n}$ is the iteration described at the end of Section 8.1.

We can also state and prove the version of Theorem 1.2 for a primitively $p$-renormalizable $\psi$-quadratic-like map $\mathbf{f} \equiv(f, i): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$, with $K, \mathcal{K}$ defined as before, and geodesics $\gamma$ around $K$ and $\gamma_{p}$ around the central small Julia set in $\mathbf{U} \backslash \mathcal{K}$ :

Theorem 9.2. There exists $\epsilon>0$, such that for all $p$ there exists M: Suppose that is a $\psi$-quadratic-like map that is primitively $p$ renormalizable. Then, with the notation as described just above,

$$
\begin{equation*}
\left|\gamma_{p}\right|>M \Longrightarrow|\gamma| \geqslant \epsilon p\left|\gamma_{p}\right| . \tag{9.3}
\end{equation*}
$$

Proof. This follows from Theorem 9.1 in exactly the same way that Theorem 1.2 follows from Theorem 3.1.

We can now prove our "if it's bad now, it was worse earlier" statement for geodesics around the central small Julia sets. We return to our original setting, where $f$ is a quadratic polynomial that is primitively renormalizable of periods $p_{1}\left|p_{2}\right| \ldots$, with $1<p_{k+1} / p_{k} \leqslant B$. We let $\gamma_{k}$ be the geodesic around $\mathcal{K}_{p_{k}}(0)$ in $\mathbb{C} \backslash \mathcal{K}_{p_{k}}$.

Theorem 9.4. There exists l such that for all B there exists $M$ : Suppose that $f$ is infinitely $B$-bounded primitively renormalizable (with notation as above). Then, for all $n$,

$$
\begin{equation*}
\left|\gamma_{n+l}\right|>M \Longrightarrow\left|\gamma_{n}\right|>2\left|\gamma_{n+l}\right| . \tag{9.5}
\end{equation*}
$$

Proof. Choose $l$ such that $2^{l} \epsilon>2$, where $\epsilon$ is the $\epsilon$ from Theorem 9.2. Let $P=B^{l}$, and choose $M$ such that $M$ is sufficiently large in Theorem 9.2 , for all $p \leqslant P$. We claim that (9.5) holds for these choices (and all $n)$. To see this, suppose that $\left|\gamma_{n+l}\right|>M$, and let $\mathbf{g} \equiv(i, g): \mathbf{U}^{\prime} \rightarrow \mathbf{U}$ be the canonical renormalization of $f$ of period $p_{n}$. Let $q=p_{n+l} / p_{n}$; then $\mathbf{g}$ is $q$-renormalizable, and by Theorem 8.6, we have

$$
\begin{equation*}
\left|\gamma_{n}\right|=|\gamma(g)| \text { and }\left|\gamma_{q}(g)\right| \geqslant\left|\gamma_{n+l}\right|>M \tag{9.6}
\end{equation*}
$$

So we can apply Theorem 9.2 to obtain

$$
|\gamma(g)| \geqslant \epsilon q\left|\gamma_{q}(g)\right| \geqslant \epsilon 2^{l}\left|\gamma_{q}(g)\right| \geqslant 2\left|\gamma_{q}(g)\right|,
$$

which combined with (9.6) yields (9.5).
We can also observe
Lemma 9.7. For all $P$ there is an $M \equiv M(P)$ such that if $f$ is a primitively renormalizable quadratic polynomial of period $p \leqslant P$, then $\left|\gamma_{p}\right| \leqslant M$.

Proof. This follows immediately from Theorem 1.2, taking the maximum over all $p \leqslant P$, and noting that $|\gamma(f)|=0$ for a polynomial.

Lemma 9.7 can also be proven by a standard compactness argument. The following is probably known, but we include it for completeness:

Theorem 9.8. Suppose $f$ is primitively p-renomalizable quadratic polynomial, and $\left|\gamma_{p}\right| \leqslant M$. Then $f$ has a (Douady-Hubbard) renormalization of period $p$ and modulus at least $m>0$, for $m \equiv m(M)$.
Proof. Let $V_{p} \subset \mathbb{C}$ be the (open) disk bounded by $\gamma_{p}$; and let $V_{p}^{\prime}$ be $V_{p}$ pulled back along $\mathcal{K}_{p}$, so $V_{p}^{\prime}$ is the unique component of $f^{-p} V_{p}$ that contains $K_{p}$. Then $f^{p}: V_{p}^{\prime} \rightarrow V_{p}$ is a degree 2 cover; we claim that $\overline{V_{p}^{\prime}} \subset V_{p}$, and that $\bmod \left(V_{p}, V_{p}^{\prime}\right)>m(M)$; this would then complete the
proof of the Theorem. We let $A_{p}=V_{p} \backslash K_{p}$ and $A_{p}^{\prime}=V_{p}^{\prime} \backslash K_{p}$; we need only show the outer boundary $\partial_{\text {out }} A_{p}^{\prime}$ of $A_{p}^{\prime}$ lies in $A_{p}$, with a definite modulus between $\partial_{\text {out }} A_{p}^{\prime}$ and $\partial_{\text {out }} A_{p}$.

Now let $(i, g): U_{p}^{\prime} \rightarrow U_{p}$ be the canonical renormalization of $f$ of period $p$. Then we will also denote by $\gamma_{p}$ and $A_{p}$ their lifts to $U_{p}$ and likewise for $\gamma_{p}^{\prime}$ and $A_{p}$ in $U_{p}^{\prime}$; then $g: A_{p}^{\prime} \rightarrow A_{p}$ and $g: \gamma_{p}^{\prime} \rightarrow \gamma_{p}$ are degree 2 coverings. We can give $H_{p} \equiv U_{p} \backslash K_{p}$ the poincare metric from doubling along its inner ideal boundary, and likewise $H_{p}^{\prime} \equiv U_{p}^{\prime} \backslash K_{p}^{\prime}$. We then observe that $i: H_{p}^{\prime} \rightarrow H_{p}$ is a contraction by the Schwarz lemma; it is strictly contracting at each point, because it is not a homeomorphism. The inner boundaries of $A_{p}^{\prime}$ and $A_{p}$ are the same of those of $H_{p}^{\prime}$ and $H_{p}$, and the outer boundaries for $A_{p}^{\prime}$ and $A_{p}$ (in $H_{p}^{\prime}$ and $H_{p}$ ) are the same constant distance from the inner boundaries (because $\gamma_{p}$ and $\gamma_{p}^{\prime}$ are core curves for $H_{p}$ and $H_{p}^{\prime}$ ). A simple compactness argument shows that $i$ is uniformly contracting on $A_{p}^{\prime}$, depending only on $\bmod \left(A_{p}\right)=\pi /\left(2\left|\gamma_{p}\right|\right)$. It follows that $\overline{i\left(A_{p}^{\prime}\right)} \subset A_{p}$, and $\bmod \left(A_{p}, i\left(A_{p}^{\prime}\right)\right) \geqslant m \equiv m(M)>0$. We then have the desired inclusion and bound on modulus for our original $A_{p}^{\prime}$ and $A_{p}$.

We can now conclude the a priori bounds for bounded-primitive type. We return to the same setting as in Theorem 9.4.

Theorem 9.9. For all $B$ there exists $M, m>0$ : Suppose $f$ is infinitely $B$-bounded primitively renormalizable, with periods $p_{1} \mid p_{2} \ldots$ Then, for all $n$,
(1) $\left|\gamma_{n}\right| \leqslant M$, and
(2) There is a renormalization of $f$ of period $p_{n}$ with modulus at least $m$.

Proof. Let us first prove Statement (1). We take $l$ as in Theorem 9.4, and given $B$, we take $M$ to be the greater value of $M(B)$ in Theorem 9.4 and $M\left(B^{l}\right)$ in Lemma 9.7. Now suppose there were an $n$ with $\left|\gamma_{p_{n}}\right|>M$; we take the least such value of $n$. If $n \leqslant l$ then $p_{n} \leqslant B^{l}$, and we contradict Lemma 9.7. If $n>l$, then $\left|\gamma_{p_{n-l}}\right| \geqslant 2\left|\gamma_{p_{n}}\right|$ by Theorem 9.4, so $n$ was not minimal. This concludes the proof of (1). Statement (2) then follows from (1) and Theorem 9.8.

## 10. Appendix A: Extremal length and width

There is a worth of sources containing background material on extremal length, see, e.g., the book of Ahlfors [A1]. We will briefly summarize the necessary minimum (see also the Appendix of [KL1]).
10.1. Definitions. Let $\mathcal{G}$ be a family of curves on a Riemann surface $U$. Given a (measurable) conformal metric $\mu=\mu(z)|d z|$ on $U$, let

$$
\mu(\mathcal{G})=\inf _{\gamma \in \mathcal{G}} \mu(\gamma)
$$

where $\mu(\gamma)$ stands for the $\mu$-length of $\gamma$. The length of $\mathcal{G}$ with respect to $\mu$ is defined as

$$
\mathcal{L}_{\mu}(\mathcal{G})=\frac{\mu(\mathcal{G})^{2}}{\operatorname{area}_{\mu}(U)},
$$

where area $_{\mu}$ is an area form of $\mu$. Taking the supremum over all conformal metrics $\mu$, we obtain the extremal length $\mathcal{L}(\mathcal{G})$ of the family $\mathcal{G}$.

The extremal width is the inverse of the extremal length:

$$
\mathcal{W}(\mathcal{G})=\mathcal{L}^{-1}(\mathcal{G})
$$

It can be also defined as follows. Consider all conformal metrics $\mu$ such that $\mu(\gamma) \geqslant 1$ for any $\gamma \in \mathcal{G}$. Then $\mathcal{W}(\mathcal{G})$ is the infimum of the areas $\operatorname{area}_{\mu}(U)$ of all such metrics.
10.2. Electric circuits laws. We say that a family $\mathcal{G}$ of curves overflows a family $\mathcal{G}^{\prime}$ if any curve of $\mathcal{G}$ contains some curve of $\mathcal{G}$. Let us say that $\mathcal{G}$ disjointly overflows two families, $\mathcal{G}_{1}$ and $\mathcal{G}_{2}$, if any curve of $\gamma \in \mathcal{G}$ contains the disjoint union $\gamma_{1} \sqcup \gamma_{2}$ of two curves $\gamma_{i} \in \mathcal{H}_{i}$.

The following crucial properties of the extremal length and width show that the former behaves like the resistance in electric circuits, while the latter behaves like conductance.

Series Law. Let $\mathcal{G}$ be a family of curves that disjointly overflows two other families, $\mathcal{G}_{1}$ and $\mathcal{G}_{2}$. Then

$$
\mathcal{L}(\mathcal{G}) \geqslant \mathcal{L}\left(\mathcal{G}_{1}\right)+\mathcal{L}\left(\mathcal{G}_{2}\right),
$$

or equivalently,

$$
\mathcal{W}(\mathcal{G}) \leqslant \mathcal{W}\left(\mathcal{G}_{1}\right) \bigoplus \mathcal{W}\left(\mathcal{G}_{2}\right)
$$

Parallel Law. For any two families $\mathcal{G}_{1}$ and $\mathcal{G}_{2}$ of curves we have:

$$
\mathcal{W}\left(\mathcal{G}_{1} \cup \mathcal{G}_{2}\right) \leqslant \mathcal{W}\left(\mathcal{G}_{1}\right)+\mathcal{W}\left(\mathcal{G}_{2}\right) .
$$

If $\mathcal{G}_{1}$ and $\mathcal{G}_{2}$ are contained in two disjoint open sets, then

$$
\mathcal{W}\left(\mathcal{G}_{1} \cup \mathcal{G}_{2}\right)=\mathcal{W}\left(\mathcal{G}_{1}\right)+\mathcal{W}\left(\mathcal{G}_{2}\right)
$$

10.3. Transformation rules. Both extremal length and extremal width are conformal invariants. More generally, we have:

Lemma 10.1. Let $f: U \rightarrow V$ be a holomorphic map between two Riemann surfaces, and let $\mathcal{G}$ be a family of curves on $U$. Then

$$
\mathcal{L}(f(\mathcal{G})) \geqslant \mathcal{L}(\mathcal{G}) .
$$

See Lemma 4.1 of [KL1] for a proof.
Corollary 10.2. Under the circumstances of the previous lemma, let $\mathcal{H}$ be a family of curves in $V$ satisfying the following lifting property: any curve $\gamma \in \mathcal{H}$ contains an arc that lifts to some curve in $\mathcal{G}$. Then $\mathcal{L}(\mathcal{H}) \geqslant \mathcal{L}(\mathcal{G})$.
Proof. The lifting property means that the family $\mathcal{H}$ overflows the family $f(\mathcal{G})$. Hence $\mathcal{L}(\mathcal{H}) \geqslant \mathcal{L}(f(\mathcal{G}))$, and the conclusion follows.
Corollary 10.3. Let $Q$ and $Q^{\prime}$ be two quadrilaterals, and let $e: Q \rightarrow$ $Q^{\prime}$ be a holomorphic map that maps the horizontal sides of $Q$ to the horizontal sides of $Q^{\prime}$. Then $\mathcal{W}(Q) \leqslant \mathcal{W}\left(Q^{\prime}\right)$.
Proof. Let $\mathcal{G}$ (resp. $\mathcal{G}^{\prime}$ ) be the family of horizontal curves in $Q$ (resp., in $Q^{\prime}$ ). Since the horizontal sides of $Q$ are mapped to the horizontal sides of $Q^{\prime}$, these families satisfy the lifting property of the previous Corollary. Hence $\mathcal{L}(\mathcal{G}) \leqslant \mathcal{L}\left(\mathcal{G}^{\prime}\right)$, and we are done.
Lemma 10.4. Let $f: U \rightarrow V$ be a branched covering between two compact Riemann surfaces with boundary. Let $A$ be an archipelago in $U$, $B=f(A)$, and assume that $f: A \rightarrow B$ is a branched covering of degree d. Then

$$
\bmod (V, B) \geqslant d \bmod (U, A)
$$

See Lemma 4.3 of [KL1] for a proof.
Given two compact subsets $A$ and $B$ in a Riemann surface $S$, let $\mathcal{W}_{S}(A, B)$ stand for the extremal width between them, i.e., the extremal width of the family of curves connecting $A$ to $B$.

Lemma 10.5. Let $S$ and $S^{\prime}$ be two compact Riemann surfaces with boundary. and let $f: S \rightarrow S^{\prime}$ be a branched covering of degree $D$. Let $S^{\prime}=A^{\prime} \sqcup B^{\prime}$, where $A^{\prime}$ and $B^{\prime}$ are closed, and let $A=f^{-1}\left(A^{\prime}\right), B=$ $f^{-1}\left(B^{\prime}\right)$. Then

$$
\mathcal{W}_{S}(A, B)=D \mathcal{W}\left(A^{\prime}, B^{\prime}\right)
$$

See [A1] for a proof. It makes use of the fact that the extremal width $\mathcal{W}(A, B)$ is achieved on the harmonic foliation $\mathcal{F}=\mathcal{F}_{S}(A, B)$ connecting $A$ and $B$, i.e., the gradient foliation of the harmonic function $\omega$ vanishing on 0 and equal to 1 on $B$. Hence $\mathcal{W}_{S}(A, B)$ is equal to the $l_{1}$-norm of the associated WAD $W_{\mathcal{F}}$.
10.4. Non-Intersection Principle. The following important principle says that two wide quadrilaterals cannot go non-trivially one across the other:

Lemma 10.6. Let us consider two quadrilaterals, $Q_{1}$ and $Q_{2}$, endowed with the vertical foliations $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$. If $\mathcal{W}\left(Q_{i}\right) \geqslant 1$ then

- either there exists a pair of disjoint leaves $\gamma_{i}$ of the foliations $\mathcal{F}_{i}$;
- or $\mathcal{W}\left(Q_{1}\right)=\mathcal{W}\left(Q_{2}\right)=1$, and the rectangles perfectly match in the sense that the vertical sides of one of them coincide with the horizontal sides of the other.

Proof. Assume that the first option is violated, so that, every leaf of $\mathcal{F}_{1}$ crosses every leaf of $\mathcal{F}_{2}$. Let us uniformize our rectangles by standard rectangles, $\phi_{i}: \mathbf{Q}\left(a_{i}\right) \rightarrow Q_{i}$, where $a_{i} \geqslant 1$. Let $\lambda_{i}$ be the Euclidean metrics on the $\mathbf{Q}\left(a_{i}\right)$, and let $\mu_{i}=\left(\phi_{i}\right)_{*}\left(\lambda_{i}\right)$. Since every vertical leave $\gamma$ of $Q_{2}$ crosses every vertical leaf of $Q_{1}, \mu_{1}(\gamma) \geqslant a_{1}$. Hence $a_{2}=\mathcal{W}\left(Q_{2}\right) \leqslant 1 / a_{1}$, which implies that $a_{1}=a_{2}=1$. Thus, both $\mathbf{Q}\left(a_{i}\right) \equiv \mathbf{Q}$ are the squares.

Moreover, $\mu_{1}$ must be the extremal metric for $\mathcal{F}_{2}$. Since the extremal metric is unique (up to scaling), we conclude that $\mu_{1}=\mu_{2}$. Hence $\phi_{2}^{-1} \circ \phi_{1}: \mathbf{Q} \rightarrow \mathbf{Q}$ is the isometry of the square, and the conclusion follows.

## 11. Appendix B: Elements of electric engineering

11.1. Potentials, currents and conductances. Electric circuit $\mathcal{C}$ is

- A connected graph $\Gamma$ with two marked vertices (battery). We let $\mathcal{E}=\mathcal{E}(\Gamma)$ be the set of edges of $\Gamma, \mathcal{V}=\mathcal{V}(\Gamma)$ be the set of its vertices, and $\mathcal{B}=\{a, b\} \subset \mathcal{V}$ be the battery.
- A conductance vector $W=\sum_{e \in \mathcal{E}} W(e) e$, where $W(e)>0$ for all $e \in \mathcal{E}$.

The edges of $\Gamma$ are interpreted as resistors of the circuit with conductances $W(e)$. The inverse numbers $R(e)=W(e)^{-1}$ as their resistances.

We write $x \sim y$ for two neighboring vertices of $\Gamma$. The vertices of $\mathcal{V} \backslash \mathcal{B}$ will also be called internal. If we forget the battery $\mathcal{B}$, we call the circuit "unplugged".

A potential distribution on $\mathcal{C}$ is a function $U: \mathcal{V} \rightarrow \mathbb{R}$.
Let $\mathcal{E}^{*}$ stand for the set of all possible oriented edges $\mathcal{E}$. An oriented edge $e^{*}$ can be written as $[x, y]$ where $x, y$ are its endpoints ordered according to the orientation of $e$. We also write $-e^{*}$ for the edge $e^{*}$ with the opposite orientation.

A current on $\mathcal{C}$ is an odd function $I: \mathcal{E}^{*} \rightarrow \mathbb{R}$, i.e., $I\left(-e^{*}\right)=-I\left(e^{*}\right)$.

A potential $U$ induces potential differences

$$
d U[x, y]=U(y)-U(x)
$$

on the oriented edges of $\Gamma$ (negative coboundary of $U$ ). It forces current

$$
I\left(e^{*}\right)=-W(e) d U e^{*}, \quad e^{*} \in \mathcal{E}^{*}
$$

to run through the resistors. The energy $E(e)$ of this current is equal to $I(e) d U(e)=W(e)(d U(e))^{2}$ (note that it is independent of the orientation of $e$ ), so that, the total energy of this potential distribution is equal to

$$
E=E(U)=\sum_{e \in \mathcal{E}} W(e)(d U(e))^{2} .
$$

The quantity $\mathbf{U}=U(a)-U(b)$ is called the battery potential.
11.2. Equilibrium. Let us define the boundary of the current $I$ as the following function on the vertices:

$$
\partial I(x)=\sum_{y \sim x} I[x, y] .
$$

We say that the potential distribution is in equilibrium, $U_{\text {eq }}$, if the current is conserved, i.e., $\partial I \equiv 0$ on $\mathcal{V} \backslash \mathcal{B}$. (In other words, a conserved current is a relative 1 -cycle on $(\Gamma, \mathcal{B})$ ). It is equivalent to saying that the potential $U$ is " $W$-harmonic" on $\Gamma \backslash \mathcal{B}$.

We say that a potential distribution $U$ is normalized if $U(a)=1$, $U(b)=0$.

Lemma 11.1. There exists a unique normalized equilibrium potential distribution on $\mathcal{C}$. This distribution is energy minimizing.

Proof. The restriction of the energy function $E$ to the affine subspace

$$
\mathcal{A}=\left\{U \in \mathbb{R}^{|\mathcal{E}|}: U(a)=1, U(b)=0\right\}
$$

is a positive quadratic function. One can easily check that $E(U) \rightarrow+\infty$ as $U \rightarrow \infty, U \in \mathcal{A}$, so that, $U$ has a global minimum $U_{\text {eq }}$ on $\mathcal{A}$. Moreover, $E$ is strictly convex, an hence can have at most one critical point. Hence $U_{e q}$ is the only critical point of $U$ on $\mathcal{A}$. Finally the criticality condition gives exactly the conservation law for the corresponding current.

Remark 1. Since the energy function $E(U)$ depends only on the potential differences, it is invariant under translations $U \mapsto U+t \mathbf{1}$, $t \in \mathbb{R}$, where $\mathbf{1} \equiv 1$. Thus, we can always normalize $U$ so that $U(b)=0$. Since $E(U)$ is homogeneous of order 2 in $U$, normalization $\mathbf{U}=\lambda$ would replace the equilibrium distribution $U_{e q}$ by $U_{e q}^{\lambda}=\lambda U_{e q}$. Then the equilibrium current would scale proportionally: $I_{e q}^{\lambda}=\lambda I_{e q}$.

Remark 2. The above lemma is just a solution of the Dirichlet boundary problem by minimizing the Dirichlet integral.

Lemma 11.2. For the equilibrium current $I_{\text {eq }}$, we have $\partial I_{\text {eq }}(a)=$ $-\partial I_{e q}(b)$.

Proof. We have:

$$
0=\sum_{e^{*} \in \mathcal{E}^{*}} I_{e q}\left(e^{*}\right)=\sum_{x \in \mathcal{V}} \partial I_{e q}(x)=\partial I_{e q}(a)+\partial I_{e q}(b),
$$

where the first equality is valid since $I$ is odd, the second is a rearrangement of terms, and the last one comes from the conservation law.

The total equilibrium current of the circuit $\mathcal{C}$ with battery potential $\mathbf{U}=\lambda$ is defined as $\mathbf{I}=I_{e q}^{\lambda}(a)=-I_{e q}^{\lambda}(b)$. It depends linearly on the battery potential, so we can define the total conductance of $\mathcal{C}$ as

$$
\mathbf{W}=\mathbf{W}(\mathcal{C})=\frac{\mathbf{I}}{\mathbf{U}}
$$

The total Resistance of the circuit is the inverse quantity: $\mathbf{R}=1 / \mathbf{W}$.
Let $\mathbf{E}=E\left(U_{e q}^{\lambda}\right)$ stand for the equilibrium energy of the circuit.
Lemma 11.3. At the normalized equilibrium, we have: $\mathbf{E}=\mathbf{W}$.
Proof. Let $\mathcal{A}_{0}=\left\{U \in \mathbb{R}^{|\mathcal{E}|}: U(b)=0\right\}=\mathbb{R}^{|\mathcal{E}|-1}$. Then the equilibrium state is the critical point of the quadratic form $E(U)=(Q U, U)$ on $\mathcal{A}_{0}$ subject of the restriction $U(a)=1$ (here $Q$ is the matrix of $E \mid \mathcal{A}_{0}$ ). By the Lagrange multipliers method, this stationary point satisfies the equation

$$
\frac{\partial E}{\partial U} \equiv 2 Q U=\lambda V,
$$

where $V$ is a basic vector in $R^{|\mathcal{E}|-1}$ with $V(a)=1, V(x)=0$ for $x \in \mathcal{V} \backslash \mathcal{B}$. From here we conclude:
(i) $2 E(U)=2(Q U, U)=\lambda(V, U)=\lambda U(a)=\lambda$;
(ii) $\frac{\partial E}{\partial U(a)}=\lambda V(a)=\lambda$.

But

$$
\frac{\partial E}{\partial U(a)}=\frac{\partial}{\partial U(a)} \sum_{x \sim a} W[a, x](U(a)-U(x))^{2}=2 \sum_{x \sim a} I[a, x]=2 \mathbf{W}
$$

and we are done.
Remark. So, in the normalized situation (when $\mathbf{U}=1$ ), we have

$$
\begin{equation*}
\mathbf{E}=\mathbf{I}=\mathbf{W} \tag{11.4}
\end{equation*}
$$

Since $\mathbf{I}$ is proportional to $\mathbf{U}, \mathbf{E}$ quadratically depends on $\mathbf{U}$, and $\mathbf{W}$ is independent of $\mathbf{U}$, we obtain by scaling the following physically obvious formulas:

$$
\mathbf{E}=\mathbf{U} \cdot \mathbf{I}, \quad \mathbf{W}=\frac{\mathbf{I}}{\mathbf{U}}
$$

11.3. Series and Parallel Laws. Given two circuits $\mathcal{C}_{1}$ and $\mathcal{C}_{2}$, we can put them in series, that is, to match the terminal battery pole of $\mathcal{C}_{1}$ to the initial battery pole of $\mathcal{C}_{2}$, and to declare the "free" battery poles a new battery of this "connected sum". More formally, let $\mathcal{B}_{i}=\left\{a_{i}, b_{i}\right\}$ be the battery of $\mathcal{C}_{i}=\left(\Gamma_{i}, W_{i}\right)$. Then we let $\mathcal{C}=C_{1} \sqcup_{a_{2}=b_{1}} C_{2}$ with the battery $\left\{a_{1}, b_{2}\right\}$ and the conductance vector: $W=W_{1}+W_{2}$. (direct sum).

Lemma 11.5 (Series law). If $\mathcal{C}$ is a series of two circuits $\mathcal{C}_{1}$ and $\mathcal{C}_{2}$ with conductances $\mathbf{W}, \mathbf{W}_{1}$ and $\mathbf{W}_{2}$ respectively, then

$$
\mathbf{W}=\mathbf{W}_{1} \oplus \mathbf{W}_{2}
$$

Proof. Let us consider the equilibrium state $(U, I)$ of the circuit $\mathcal{C}$. Then the conservation law is valid for both sub-circuits $\mathcal{C}_{1}$ and $\mathcal{C}_{2}$, so that they are both in the equilibrium state, with the battery potentials $\mathbf{U}_{1}=1-U(c)$ and $\mathbf{U}_{2}=U(c)$, where $c=b_{1}=a_{2}$. Let $\mathbf{I}_{i}$ be the corresponding equilibrium currents through the $\mathcal{C}_{i}$. By the current conservation at vertex $c, \mathbf{I}_{1}=\mathbf{I}_{2}$. In fact, $\mathbf{I}_{1}=\partial I\left(a_{1}\right)=\mathbf{I}$ and $\mathbf{I}_{2}=$ $-\partial I\left(b_{2}\right)=\mathbf{I}$. Hence

$$
\frac{1}{\mathbf{W}_{1}}=\frac{\mathbf{U}_{1}}{\mathbf{I}}, \quad \frac{1}{\mathbf{W}_{2}}=\frac{\mathbf{U}_{2}}{\mathbf{I}} .
$$

Summing these up, we obtain the desired.
Given two circuits as above, we can also put them in parallel, that is, to identify the two pairs of poles as follows: $a_{1} \sim a_{2}, b_{1} \sim b_{2}$.

Lemma 11.6 (Parallel Law). If $\mathcal{C}$ is a parallel of two circuits $\mathcal{C}_{1}$ and $\mathcal{C}_{2}$ as above, then

$$
\mathbf{W}=\mathbf{W}_{1}+\mathbf{W}_{2} .
$$

11.4. Quotients. Let us have a finite family of electric circuits $\mathcal{C}_{i}$ based on graphs $\Gamma_{i}$ with batteries $\mathcal{B}_{i}=\left\{a_{i}, b_{i}\right\}$ and conductance vectors $\mathcal{W}_{i}$. Let us identify certain vertices of the disjoint union $\sqcup \Gamma_{i}$ so that the batteries of different $\mathcal{C}_{i}$ 's get identified ( $a_{i}=a_{j}, b_{i}=b_{j}$ ) and no internal vertices get identified with a battery vertex. If in the quotient graph there are several edges $e_{i j} \in \mathcal{E}_{i} \equiv \mathcal{E}\left(\Gamma_{i}\right)$ connecting the same pair
of vertices, let us identify those edges, too. We obtain a graph $\Gamma$. For the edge $e$ of $\mathcal{E} \equiv \mathcal{E}(\Gamma)$ obtained this way, let

$$
W(e)=\sum_{e_{i j} \sim e} W\left(e_{i j}\right)
$$

In this way we obtain a new electric circuit $\mathcal{C}=(\Gamma, \mathcal{B}, \mathcal{W})$ called a quotient of $\sqcup \mathcal{C}_{i}$.

Let $\mathbf{W}=\mathcal{W}(\mathcal{C}), \mathbf{W}_{i}=\mathcal{W}\left(\mathcal{C}_{i}\right)$.
Lemma 11.7. If $\mathcal{C}$ is a quotient of $\sqcup \mathcal{C}_{i}$ then $\mathbf{W} \geqslant \sum \mathbf{W}_{i}$.
Proof. Let $U$ be the normalized equilibrium potential for $\mathcal{C}$.
By Lemma 11.3, its energy $E(U)$ is equal to $\mathbf{W}$.
The potential $U$ lifts to normalized potentials $U_{i}$ for the circuits $\mathcal{C}_{i}$ 's. By the same lemma, $\mathbf{W}_{i} \leqslant E\left(U_{i}\right)$. On the other hand,

$$
\sum_{i} E\left(U_{i}\right)=\sum_{i} \sum_{e_{i j} \in \mathcal{E}_{i}} W\left(e_{i j}\right) d U_{i}\left(e_{i j}\right)^{2}=\sum_{e \in \mathcal{E}} d U(e)^{2} \sum_{e_{i j} \sim e} W\left(e_{i j}\right)=E(U) .
$$

Putting these pieces together, we obtain the desired inequality.
Remark 11.1. The above inequality has a clear physical meaning: taking a quotient of a circuit gives more choices for the current to flow, which increases total conductance.

The following arithmetic inequality can be proven by interpreting it in terms of electric circuits:

Lemma 11.8. Let us consider finite sets of positive numbers $w_{i}$ and $v_{i j}, 1 \leqslant i, j \leqslant n$, such that $w_{i} \leqslant \oplus_{j} v_{i j}$. Let $w=\sum w_{i}$ and $v_{j}=\sum_{i} v_{i j}$. Then $w \leqslant \oplus v_{j}$.

Proof. Let us consider resistors $V_{i j}$ with conductances $v_{i j}$. Let $\mathcal{C}_{i}$ be an electric circuit obtained by plugging the resistors $V_{i j}$ in series. Then

$$
\mathbf{W}\left(\mathcal{C}_{i}\right)=\bigoplus_{j} v_{i j} \geqslant w_{i} .
$$

Let us also consider the quotient $\mathcal{C}$ of $\sqcup \mathcal{C}_{i}$ obtained by identifying the endpoints of the resistors $V_{i j}$ with the same $j$. In other words, we plug the $V_{i j}$ 's with the same $j$ in parallel obtaining circuits $V_{j}$, and then plug the $V_{i}$ 's in series. Then

$$
\mathbf{W}(\mathcal{C})=\oplus \mathbf{W}\left(V_{j}\right)=\oplus v_{j}
$$

By the previous lemma, $\sum \mathbf{W}\left(\mathcal{C}_{i}\right) \leqslant \mathbf{W}(\mathcal{C})$, which boils down to the desired estimate.

Remark 11.2. So, the signs of ordinary and harmonic sums can be interchanged following the rule:

$$
\sum_{i} \bigoplus_{j} v_{i j} \leqslant \bigoplus_{j} \sum_{i} v_{i j} .
$$

11.5. Local conductances. Given a vertex $x \in \mathcal{V}(\Gamma)$, we let

$$
W \mid x=\sum_{y \sim x} W[x, y]
$$

and call it the local conductance at $x$.
Lemma 11.9. We have: $\mathbf{W} \leqslant \min (W|a, W| b)$.
Proof. By the Maximum Principle for harmonic functions, $0 \leqslant U_{e q}(x) \leqslant$ 1 for any $x \in \Gamma$. Hence $\left|I_{e q}\left(e^{*}\right)\right| \leqslant W(e)$ for any $e^{*} \in \mathcal{E}^{*}$. Together with (11.4), this implies:

$$
\mathbf{W}=\mathbf{I}=\sum_{x \sim a} I_{e q}[x, a] \leqslant W \mid a,
$$

and similarly for $W \mid b$.
11.6. Domination. Let us consider two unplugged electric circuits $\mathcal{C}=(\Gamma, W)$ and $\mathcal{C}^{\prime}=\left(\Gamma^{\prime}, W^{\prime}\right)$ such that:

- $\Gamma^{\prime}$ is obtained from $\Gamma$ by replacing edges $e \in \mathcal{E}(\Gamma)$ with some graphs $\Gamma^{\prime}(e)$;
- Letting $\mathcal{C}^{\prime}(e)$ be the restriction of $\mathcal{C}^{\prime}$ to $\Gamma^{\prime}(e)$ with battery $\partial e$, we have: $\mathbf{W}\left(\mathcal{C}^{\prime}(e)\right) \geqslant W(e)$ for any $e \in \mathcal{E}(\Gamma)$.

Under these circumstances we say that $\mathcal{C}^{\prime}$ dominates $\mathcal{C}, \mathcal{C}^{\prime} \multimap \mathcal{C}$.
Lemma 11.10. If $\mathcal{C}^{\prime} \multimap \mathcal{C}$ then $W^{\prime}|x \geqslant W| x$ for any $x \in \mathcal{V}(\Gamma)$.
Proof. Indeed, by Lemma 11.9, we have:

$$
W^{\prime}\left|x=\sum_{\partial e^{\prime} \ni x} W^{\prime}\left(e^{\prime}\right) \geqslant \sum_{\partial e \ni x} \mathbf{W}\left(\mathcal{C}^{\prime}(e)\right) \geqslant \sum_{\partial e \ni x} W(e)=W\right| x,
$$

where the summation is taken over $e \in \mathcal{E}(\Gamma), e^{\prime} \in \mathcal{E}\left(\Gamma^{\prime}\right)$.
11.7. Trees of complete graphs. In this section we will consider a special class of electric circuits based on trees of complete graphs (TCG). A TCG is an object that can be constructed inductively by the following rules:

- Any complete graph is a TCG;
- If $\Gamma, \Gamma^{\prime}$ are TCG's, $v \in \mathcal{V}(\Gamma), v^{\prime} \in \mathcal{V}\left(\Gamma^{\prime}\right)$, then $\Gamma \underset{v=v^{\prime}}{\sqcup} \Gamma^{\prime}$ is a tree of complete graphs as well.

A TCG is called an interval of complete graphs if any complete graph involved has a common vertex with at most two other complete graphs, and now three complete graphs have a common vertex.

Given three vertices $x, y, z$ in a graph $\Gamma$, we say that a vertex $y$ separates $x$ from $z$ if $x$ and $z$ belong to different components of $\Gamma \backslash\{y\}$.

We say that a sequence of vertices $\left(x_{0}, x_{1}, \ldots, x_{d}\right)$ form a chain in $\Gamma$ if $x_{i}$ separates $x_{i-1}$ from $x_{i+1}$ for each $i=1, \ldots, d-1$. Let $S S(x, y)$ stand for the set of vertices separating $x$ from $y$.

The reader can entertain himself by verifying the following fact:
Lemma 11.11. Let $\Gamma$ be a TCG. Then for any two vertices $x, y \in$ $\mathcal{V}(\Gamma)$, the set $S S(x, y) \cup\{x, y\}$ can be uniquely ordered to form a chain $\left(x=x_{0}, x_{1}, \ldots, x_{d}=y\right)$. Moreover, for any $i=0, \ldots, d-1$, the vertices $x_{i}$ and $x_{i+1}$ belong to the same complete graph $\Gamma_{k(i)}$, and these graphs form an interval of complete graphs.

We call it "the chain connecting $x$ to $y$ ", and we let $d_{\Gamma}(x, y)=d$.
Lemma 11.12. Let us consider an electric circuit $\mathcal{C}$ based on a tree of complete graphs, and let $\left(a=x_{0}, x_{1}, \ldots, x_{d}=b\right)$ be the chain connecting the poles of the battery. Then

$$
\mathbf{W} \leqslant \bigoplus_{i=1}^{d} W\left|x_{i} \leqslant \frac{1}{d_{\Gamma}(a, b)} \max _{x \in \Gamma} W\right| x .
$$

Proof. The second inequality is trivial, so we only need to prove the first one.

Let $G_{i}$ be the component of $\Gamma \backslash\left\{x_{i}, x_{i+1}\right\}$ containing the edge $\left(x_{i}, x_{i+1}\right)$. Let $\mathcal{C}_{i}$ be the restriction of the electric circuit $\mathcal{C}$ to $G_{i} \cup\left\{x_{i}, x_{i+1}\right\}$ with battery $\left\{x_{i}, x_{i+1}\right\}$. By the Series Law and Lemma 11.9,

$$
\mathbf{W} \leqslant \bigoplus_{i=1}^{d} \mathbf{W}\left(\mathcal{C}_{i}\right) \leqslant \bigoplus_{i=1}^{d} W \mid x_{i} .
$$

## 12. Appendix C: Estimates in hyperbolic geometry

In this appendix we will describe how to measure the geometry of a surface using transverse geodesic arcs. We will then show how to compute the lengths of peripheral closed geodesics of the surface using these measurements.

Suppose $T$ is a compact hyperbolic surface with geodesic boundary. The following lemma appears as the Corollary to section 3.3 of [ Ab ]:

Lemma 12.1. There is an $\epsilon_{0}>0$ such that any two distinct closed geodesics on $T$ of length at most $\epsilon_{0}$ are simple and disjoint.

Let $\mathbf{S}$ be a compact hyperbolic surface with geodesic boundary. Then a transverse geodesic arc for $\mathbf{S}$ is a proper path of minimal length in its proper homotopy class. If $\boldsymbol{\alpha}$ is a path on $\mathbf{S}$, it is a transverse geodesic arc if and only if it is a geodesic arc that meets $\partial \mathbf{S}$ orthogonally, or, equivalently, the double of $\boldsymbol{\alpha} \cup \overline{\boldsymbol{\alpha}}$ in $\mathbf{S} \cup \overline{\mathbf{S}}$ is a closed geodesic.

Let $S$ be a compact Riemann surface with boundary, and endow $\operatorname{Int} S$ with its Poincaré metric. The peripheral geodesics on Int $S$ bound a compact surface $\mathbf{S}$ with geodesic boundary, called the convex core of $S$. There is a homeomorphism $h: \mathbf{S} \rightarrow S$ that is isotopic through embeddings to the inclusion $\mathbf{S} \subset S$. We can then form a weighted arc diagram $M^{S}$ on $S$ as follows: for $\alpha \in \mathcal{A}(S)$, we find the transverse geodesic arc $\boldsymbol{\alpha}$ for $\mathbf{S}$ such that $h(\boldsymbol{\alpha}) \sim \alpha$. Then we let $M^{S}(\alpha)=$ $-\log L(\boldsymbol{\alpha})$ if $L(\boldsymbol{\alpha})<\epsilon_{0} / 2$, and $M^{S}(\alpha)=0$ otherwise. Then $M^{S}$ is supported on a set of disjoint arcs, so $M^{S}$ is a weighted arc-diagram for $S$. We let $\boldsymbol{M}^{S}=\bigcup\left\{\boldsymbol{\alpha}: L(\boldsymbol{\alpha})<\epsilon_{0} / 2\right\}$, so $\boldsymbol{M}^{S} \subset \mathbf{S}$ is the union of the short transverse geodesic arcs of $\mathbf{S}$.

We will call a non-peripheral simple closed geodesic a dividing geodesic. The following result appears in $[\mathrm{Ab}]$ :

Lemma 12.2. Let $T$ be a compact hyperbolic Riemann surface with bounded-length geodesic boundary. Then either $T$ is a pair of pants, or there is a bounded-length dividing geodesic on $T$.

We say that a hyperbolic surface $T$ is symmetric if it admits an isometric orientation-reversing involution, which we will denote by complex conjugation: " $z \mapsto \bar{z}$ ". Then we let $E_{T}=\{z \in T: z=\bar{z}\}$, and $E_{T}$ will be a union of (simple) closed geodesics and transverse geodesic $\operatorname{arcs.}$ ( $E_{T}$ depends implicitly on on the choice of involution. Whenever we say "symmetric hyperbolic surface" we will mean "symmetric hyperbolic surface and choice of involution.") Note that $T \backslash E_{T}$ has two components, call them $A_{T}$ and $\overline{A_{T}}$, which are mapped to each other by $z \mapsto \bar{z}$. We prove a symmetric version of Lemma 12.2:

Lemma 12.3. Every symmetric compact hyperbolic surface $T$ with bounded-length geodesic boundary has a bounded-length symmetric pair of pants decomposition.

Proof. It suffices to find a single bounded-length symmetric dividing geodesic on the surface, or a symmetric pair of disjoint bounded-length dividing geodesics, because then we can cut the surface $T$ along that geodesic or pair of geodesics, and repeat.

By Theorem 12.2, unless $T$ is a pair of pants, there is a dividing geodesic $\gamma$ of bounded length on $T$. If $\gamma \cap E_{T}=\varnothing$, then $\gamma \cap \bar{\gamma}=\varnothing$, and we are done. Likewise, if $\gamma \subset E_{T}$, then $\gamma$ is symmetric, and we
are done. Otherwise, let $\eta$ be a component of $\gamma \cap \mathrm{Cl} A_{T}$. Then $\eta \cup \bar{\eta}$ is a non-trivial non-peripheral simple closed curve so we let $\tau$ be the dividing geodesic homotopic to $\eta$. Then $L(\tau) \leqslant 2 L(\eta)<2 L(\gamma)$, so $\tau$ is the desired object.

We now prove two basic estimates on transverse geodesic arcs on pairs of pants. We denote by $[x, y, z, r, s, t]$ the right-angled hyperbolic hexagon with lengths $x, y, z, r, s, t$ in that order. We will omit lengths that are not specified, so for example $[a, b,, c$, ] denotes the rightangled hexagon with alternating side lengths $a, b$, and $c$. We first estimate the length of one side in a hyperbolic right-angled hexagon, in terms of the lengths of three alternating sides:

Lemma 12.4. Let $\left[a, c^{\prime}, b,, c\right.$, ] be a hyperbolic right-angled hexagon, and suppose that $a, b, c \leqslant r$. Then $c^{\prime}=-\log a-\log b+O(1 ; r) .{ }^{11}$

Proof. We use the formula (from $[\mathrm{F}]$ ):

$$
\begin{equation*}
\cosh c^{\prime}=\frac{\cosh c+\cosh a \cosh b}{\sinh a \sinh b}=e^{O(1 ; r)} \frac{1}{a b} \tag{12.5}
\end{equation*}
$$

and recall that $\cosh ^{-1} x=\log x+O(1 ; t)$ whenever $x \geqslant t>0$.
We let $\mathcal{P}(a, b, c)$ denote the pair of pants with cuff lengths $a, b$, and c.

Lemma 12.6. We have the following two estimates:
(1) If $P=\mathcal{P}(a, b, c)$ is a pair of pants, and $\gamma$ is the transverse geodesic arc that connects $a$ and $b$, then

$$
|\gamma|=-\log a-\log b+O(1)
$$

for $a, b, c \leqslant C_{0}$.
(2) If $\gamma$ is the transverse geodesic arc that connects $a$ and $a$ in $P(a, b, b)$, then

$$
|\gamma|=-2 \log a+O(1)
$$

for $a, b \leqslant C_{0}$.
Proof. We prove each of the above:
(1) We cut $P$ along the three pairwise transversals into two rightangled hexagons. By formula (12.5) these hexagons are equal, and hence each has type $[a / 2, \gamma, b / 2,, c / 2$,$] . Apply now Lemma$ 12.4 .

[^8](2) We let $\eta$ be the transversal from the length $a$ cuff to one other. Then we cut along $\eta$ and $\gamma$ to obtain the right-angled hexagon $[a / 4, \gamma, a / 4, \eta, b, \eta]$, and then apply Lemma 12.4.

Given a closed geodesic $\gamma$ and an arc $\alpha \in \mathcal{A}$, let $\langle\gamma, \alpha\rangle$ stand for the intersection number of $\gamma$ with $\alpha$, i.e., the minimal number of intersections of $\gamma$ with the paths representing $\alpha$. Given a weighted arc diagram $W=\sum W(\alpha) \alpha$, we can define the intersection number $\langle\gamma, W\rangle=\sum W(\alpha)\langle\gamma, \alpha\rangle$ by linearity. We can now prove the main theorem of this appendix:

Theorem 12.6.1. Let $S$ be a compact Riemann surface with boundary, and endow Int $S$ with its Poincaré metric. Suppose that $\gamma$ is a peripheral closed geodesic for $S$. Then

$$
L(\gamma)=2\left\langle M^{S}, \gamma\right\rangle+O(1 ; \chi(S))
$$

Proof. We let $\mathbf{S}$ be the convex core of $S$. We find a symmetric boundedlength pair of pants decomposition for $\mathbf{S} \cup \overline{\mathbf{S}}$ extending $\boldsymbol{M}^{S} \cup \overline{\boldsymbol{M}}^{S}$. Then we can write $\gamma=\bigcup t_{i}$, where the segments $t_{i}$ are interior-disjoint, and each is a transverse arc of one of the pairs of pants. Then $L\left(t_{i}\right)=$ $-\log a_{i}-\log b_{i}+O(1)$, where $a_{i}$ and $b_{i}$ are the lengths of the cuffs that $\gamma$ connects (possibly the same cuff). Therefore

$$
L(\gamma)=2 \sum-\log a_{i}+O(1)
$$

where the $a_{i}$ are the lengths of the cuffs that $\gamma$ crosses, counted with multiplicity. But

$$
\begin{equation*}
2 \sum-\log a_{i}=2\left\langle M^{S}, \gamma\right\rangle+O(1) . \tag{12.7}
\end{equation*}
$$

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[^0]:    ${ }^{1}$ An arc is a non-trivial homotopy class of properly embedded paths.

[^1]:    ${ }^{2}$ In what follows, all Riemann surfaces are assumed to be of this type, unless otherwise is explicitly stated.
    ${ }^{3}$ Note that this homotopy is automatically isotopy.
    ${ }^{4}$ We can also think of a WAD as a function $\mathcal{A}(S) \rightarrow \mathbb{R}_{+}$supported on some arc diagram.

[^2]:    ${ }^{5}$ In what follows, all laminations under consideration are assumed to be proper.

[^3]:    ${ }^{6}$ We will use abbreviated notations $W_{\text {can }}(S)$ or $W_{\text {can }}(\alpha)$ whenever it does not lead to confusion.

[^4]:    ${ }^{7}$ See $\S 11.7$ of Appendix B for the definition.

[^5]:    ${ }^{8}$ Strictly speaking, we are defining a " $\psi$-polynomial-like map with connected Julia set".

[^6]:    ${ }^{9}$ or clever definition of $i, f$, and $K$ on the disjoint union $\bigcup_{n \in \mathbb{N}} \mathbf{U}^{n}$

[^7]:    ${ }^{10}$ While we will use this only when $g$ is a finite degree branched cover, it also applies whenever $g$ has path lifting (which of course is non-unique at branch points).

[^8]:    ${ }^{11}$ Notation $O(1 ; r)$ stands for a quantity bounded in terms of $r$.

