FIXED POINT THEOREMS FOR PLANE CONTINUA WITH APPLICATIONS

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Dedicated to Harold Bell
Contents

List of Figures ix
Preface xi
Chapter 1. Introduction 1

Part 1. Basic Theory 11

Chapter 2. Preliminaries and outline of Part 1 13
  2.1. Index 13
  2.2. Variation 14
  2.3. Classes of maps 15
  2.4. Partitioning domains 17

Chapter 3. Tools 19
  3.1. Stability of Index 19
  3.2. Index and variation for finite partitions 20
  3.3. Locating arcs of negative variation 23
  3.4. Crosscuts and bumping arcs 25
  3.5. Index and Variation for Carathéodory Loops 27
  3.6. Prime Ends 28
  3.7. Oriented maps 30
  3.8. Induced maps of prime ends 32

Chapter 4. Partitions of domains in the sphere 35
  4.1. Kulkarni-Pinkall Partitions 35
  4.2. Hyperbolic foliation of simply connected domains 38
  4.3. Schoenflies Theorem 40
  4.4. Prime ends 41

Part 2. Applications of basic theory 47

Chapter 5. Description of main results of Part 2 49
  5.1. Outchannels 49
  5.2. Fixed points in invariant continua 50
  5.3. Fixed points in non-invariant continua – the case of dendrites 50
  5.4. Fixed points in non-invariant continua – the planar case 51
  5.5. The polynomial case 52

Chapter 6. Outchannels and their properties 55
  6.1. Outchannels 55
Abstract

In this memoir we present proofs of basic results, including those developed so far by Harold Bell, for the plane fixed point problem: does every map of a non-separating plane continuum have a fixed point? Some of these results had been announced much earlier by Bell but without accessible proofs. We define the concept of the variation of a map on a simple closed curve and relate it to the index of the map on that curve: \( \text{Index} = \text{Variation} + 1 \). A prime end theory is developed through hyperbolic chords in maximal round balls contained in the complement of a non-separating plane continuum \( \mathcal{X} \). We define the concept of an outchannel for a fixed point free map which carries the boundary of \( \mathcal{X} \) minimally into itself and prove that such a map has a unique outchannel, and that outchannel must have variation \(-1\). Also Bell’s Linchpin Theorem for a foliation of a simply connected domain, by closed convex subsets, is extended to arbitrary domains in the sphere.

We introduce the notion of an oriented map of the plane and show that the perfect oriented maps of the plane coincide with confluent (that is composition of monotone and open) perfect maps of the plane. A fixed point theorem for positively oriented, perfect maps of the plane is obtained. This generalizes results announced by Bell in 1982.

A continuous map of an interval \( I \subset \mathbb{R} \) to \( \mathbb{R} \) which sends the endpoints of \( I \) in opposite directions has a fixed point. We generalize this to maps on non-invariant continua in the plane under positively oriented maps of the plane (with appropriate boundary conditions). Similar methods imply that in some cases non-invariant continua in the plane are degenerate. This has important applications in complex dynamics. E.g., a special case of our results shows that if \( \mathcal{X} \) is a non-separating invariant subcontinuum of the Julia set of a polynomial \( P \) containing no fixed Cremer points and exhibiting no local rotation at all fixed points, then \( \mathcal{X} \) must be a point. It follows that impressions of some external rays to polynomial Julia sets are degenerate.
## List of Figures

3.1 Replacing $f : S \to \mathbb{C}$ by $f_1 : S \to \mathbb{C}$ with one less subarc of nonzero variation. 22

3.2 Bell’s Lollipop. 24

3.3 $\text{var}(f, A) = -1 + 1 - 1 = -1$. 27

4.1 Maximal balls have disjoint hulls. 36

6.1 The strip $\mathcal{S}$ from Lemma 6.1.2 56

6.2 Uniqueness of the negative outchannel. 60

7.1 Replacing the links $[a_{n(1)} - 1, a_{n(1)}], \ldots, [a_{m(1)} - 1, a_{m(1)}]$ by a single link $[a_{n(1)} - 1, a_{m(1)}]$. 73

7.2 Illustration to the proof of Theorem 7.4.7. 78

7.3 Illustration to the proof of Lemma 7.4.9. 83

7.4 A general puzzle-piece 86
Preface

By a continuum we mean a compact and connected metric space and by a non-separating continuum $X$ in the plane $\mathbb{C}$ we mean a continuum $X \subset \mathbb{C}$ such that $\mathbb{C} \setminus X$ is connected. Our work is motivated by the following long-standing problem [Ste35] in topology.

Plane Fixed Point Problem: “Does a continuous function taking a non-separating plane continuum into itself always have a fixed point?”

To give the reader perspective we would like to make a few brief historical remarks (see [KW91, Bin69, Bin81] for much more information).

Borsuk [Bor35] showed in 1932 that the answer to the above question is yes if $X$ is also locally connected. Cartwright and Littlewood [CL51] showed in 1951 that a map of a non-separating plane continuum $X$ to itself has a fixed point if the map can be extended to an orientation-preserving homeomorphism of the plane. It was 27 years before Harold Bell [Bel78] extended this result to the class of all homeomorphisms of the plane. Then Bell announced in 1982 (see also Akis [Aki99]) that the Cartwright-Littlewood Theorem can be extended to the class of all holomorphic maps of the plane. For other partial results in this direction see, e.g., [Ham51, Hag71, Bel79, Min90, Hag96, Min99].

In this memoir the Plane Fixed Point Problem is addressed. We develop and further generalize tools, first introduced by Bell, to elucidate the action of a fixed point free map (should one exist). We are indebted to Bell for sharing his insights with us. Some of the results in this memoir were first obtained by him. Unfortunately, many of the proofs were not accessible. Since there are now multiple papers which rely heavily upon these tools (e.g., [OT07, BO09, BCLOS08]) we believe that they deserve to be developed in a coherent fashion. We also hope that by making these tools available to the mathematical community, other applications of these results will be found. In fact, we include in Part 2 of this text new applications which illustrate their usefulness.

Part 1 contains the basic theory, the main ideas of which are due to Bell. We introduce Bell’s notion of variation and prove his theorem that index equals variation increased by 1 (see Theorem 3.2.2). Bell’s Linchpin Theorem 4.2.5 for simply connected domains is extended to arbitrary domains in the sphere and proved using an elegant argument due to Kulkarni and Pinkall [KP94]. Our version of this theorem (Theorem 4.1.5) is essential for the results later in the paper.

Building upon these ideas, we will introduce in Part 1 the class of oriented maps of the plane and show that it decomposes into two classes, one of which preserves and the other of which reverses local orientation. The extension from holomorphic to positively oriented maps is important since it allows for simple local perturbations of the map (see Lemma 7.5.1) and significantly simplifies further usage of the developed tools.
In Part 2 new applications of these results are considered. A Zorn’s Lemma argument shows, that if one assumes a negative solution to the Plane Fixed Point Problem, then there is a subcontinuum $X$ which is minimal invariant. It follows from Theorem 6.1.4 that for such a minimal continuum, $f(X) = X$. We recover Bell’s result [Bel67] (see also Sieklucki [Sie68], and Iliadis [Ili70]) that the boundary of $X$ is indecomposable with a dense channel (i.e., there exists a prime end $E_\ell$ such that the principal set of the external ray $R_\ell$ is all of $\partial X$).

As the first application we show in Chapter 6 that $X$ has a unique outchannel (i.e., a channel in which points basically map farther and farther away from $X$) and this outchannel must have variation $-1$ (i.e., as the above mentioned points map farther and farther away from $X$, they are “flipped with respect to the center line of the channel”).

The next application of the tools developed in Part 1 directly relates to the Plane Fixed Point Problem. We introduce the class of oriented maps of the plane (i.e., all perfect maps of the plane onto itself which are the compositions of monotone and branched covering maps of the plane). The class of oriented maps consists of two subclasses: positively oriented and negatively oriented maps. In Theorem 7.1.3 we show that the Cartwright-Littlewood Theorem can be extended to positively oriented maps of the plane.

These results are used in [BO09]. There we consider a branched covering map $f$ of the plane. It follows from the above that if $f$ has an invariant and fixed point free continuum $Z$, then $f$ must be negatively oriented. We show in [BO09] that if, moreover, $f$ is an oriented map of degree 2, then $Z$ must contain a continuum $X$ such that $X$ is fully invariant (so that $X$ contains the critical point and $f|_X$ is not one-to-one). Thus, $X$ bears a strong resemblance to a connected filled in Julia set of a quadratic polynomial.

The rest of Part 2 is devoted to extending the existence of a fixed point in planar continua under positively oriented maps established in Theorem 7.1.3. We extend this result to non-invariant planar continua. First the result is generalized to dendrites; moreover, it is strengthened by showing that in certain cases the map must have infinitely many periodic cutpoints.

The above results on dendrites have applications in complex dynamics. For example, they are used in [BCO08] to give a criterion for the connected Julia set of a complex polynomial to have a non-degenerate locally connected model. That is, given a connected Julia set $J$ of a complex polynomial $P$, it is shown in [BCO08] that there exists a locally connected topological Julia set $J_{top}$ and a monotone map $m : J \to J_{top}$ such that for every monotone map $g : J \to X$ from $J$ onto a locally connected continuum $X$, there exists a monotone map $f : J_{top} \to X$ such that $g = f \circ m$. Moreover, the map $m$ has a dynamical meaning. It semi-conjugates the map $P|_J$ to a topological polynomial $P_{top} : J_{top} \to J_{top}$. In general, $J_{top}$ can be a single point. In [BCO08] a necessary and sufficient condition for the non-degeneracy of $J_{top}$ is obtained. These results extend Kiwi’s fundamental result [Kiw04] on the semi-conjugacy of polynomials without Cremer or Siegel points to all polynomials with connected Julia set.

Finally the results on the existence of fixed points in invariant planar continua under positively oriented maps are extended to non-invariant planar continua. We introduce the notion of “scrambling of the boundary” of a plane continuum $X$ under a positively oriented map and extend the fixed point results to non-invariant
continua on which the map scrambles the boundary. These conclusions are strength-
ened by showing that, under additional assumptions, a non-degenerate continuum
must either contain a fixed point in its interior, or must contain a fixed point near
which the map “locally rotates”. Hence, if neither of these is the case, then the
continuum in question must be a point. This latter result is used to show that in
certain cases impressions of external rays to connected Julia sets are degenerate.

These last named results have had other applications in complex dynamics. In
[BCLOS08] these results were used to generalize the well-known Fatou-Shishikura
inequality in the case of a polynomial $P$ (in general, the Fatou-Shishikura inequality
holds for rational functions, see [Fat20, Shi87]). For polynomials this inequality
limits the number of attracting and irrationally neutral periodic cycles by the num-
er of critical points of $P$. The improved count involves classes of (weakly recurrent)
critical points and wandering subcontinua in the Julia set.

The results in Part 1 of this memoir were mostly obtained in the late 1990’s. Most of the applications in Part 2, including the results on non-invariant plane
continua and the applications in dynamics, have been obtained during 2006–2009.
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1.0.1. Notation and the main problem. We denote the plane by \( \mathbb{C} \), the Riemann sphere by \( \mathbb{C}^\infty = \mathbb{C} \cup \{\infty\} \), the real line by \( \mathbb{R} \) and the unit circle by \( \mathbb{S}^1 = \mathbb{R}/\mathbb{Z} \). Let \( X \) be a plane compactum. Since \( \mathbb{C} \) is locally connected and \( X \) is closed, complementary domains of \( X \) are open. By \( T(X) \) we denote the topological hull of \( X \) consisting of \( X \) union all of its bounded complementary domains. Thus, \( U^\infty = U^\infty(X) = \mathbb{C}^\infty \setminus T(X) \) is the unbounded complementary component of \( X \) containing infinity. Observe that if \( X \) is a continuum, then \( U^\infty(X) \) is simply connected. The Plane Fixed Point Problem, attributed to [Ste35], is one of the central long-standing problems in plane topology. It serves as a motivation for our work and can be formulated as follows.

**Problem 1.0.1 (Plane Fixed Point Problem).** Does a continuous function taking a non-separating plane continuum into itself always have a fixed point?

1.0.2. Historical remarks. To give the reader perspective we would like to make a few historical remarks concerning the Plane Fixed Point Problem (here we cover only major steps towards solving the problem).

In 1912 Brouwer [Bro12] proved that any orientation preserving homeomorphism of the plane, which keeps a bounded set invariant, must have a fixed point (though not necessarily in that set). This fundamental result has found many important applications. It was recognized early on that the location of a fixed point should be determined if the invariant set is a non-separating continuum (in that case a fixed point should be located in the invariant continuum) and many papers have been devoted to obtaining partial solutions to the Plane Fixed Point Problem.

Borsuk [Bor35] showed in 1932 that the answer is yes if \( X \) is also locally connected. Cartwright and Littlewood [CL51] showed in 1951 that a continuous map of a non-separating continuum \( X \) to itself has a fixed point in \( X \) if the map can be extended to an orientation-preserving homeomorphism of the plane. (See Brown [Bro77] for a very short proof of this theorem based on the above mentioned result by Brouwer). The proof by Cartwright-Littlewood Theorem made use of the index of a map on a simple closed curve and this idea has remained the basic approach in many partial solutions.

The most general result was obtained by Bell [Bel67] in the early 1960’s. He showed that any counterexample must contain an invariant indecomposable subcontinuum. Hence the Plane Fixed Point Problem has a positive solution for hereditarily decomposable plane continua (i.e., for continua \( X \) which do not contain indecomposable subcontinua). Bell’s result was also based on the notion of the index of a map, but he introduced new ideas to determine the index of a simple closed curve which runs tightly around a possible counterexample. Unfortunately, these ideas were not transparent and were never fully developed. Alternative proofs

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of Bell’s result appeared soon after Bell’s announcement (see [Sie68, Ili70]). Regrettably these results did not develop Bell’s ideas.

In 1978 Bell [Bel78] used his earlier result to extend the result by Cartwright and Littlewood to the class of all homeomorphisms of the plane. Then Bell announced in 1982 (see also Akis [Aki99] where a wider class of differentiable functions was used) that the Cartwright-Littlewood Theorem can be extended to the class of all holomorphic maps of the plane. The existence of fixed points for orientation preserving homeomorphisms of the entire plane under various conditions was also considered in [Bro84, Fat87, Fra92, Gui94], and the existence of a point of period two for orientation reversing homeomorphisms in [Bon04].

As indicated above, positive results require an additional hypothesis either on the continuum $X$ (as in Borsuk’s result where the assumption is that $X$ is locally connected) or on the map (as in Bell’s case where the assumption is that $f$ is a homeomorphism of the plane). Other positive results of the first type include results by Hamilton [Ham51] ($X$ is chainable), Hagopian [Hag71] ($X$ is arcwise connected) Minc [Min90] ($X$ is the continuous image of the pseudo arc) and [Hag96] ($X$ is simply connected). Positive results of the second type require the map to be either a homeomorphism [CL51, Bel78], holomorphic (as announced by Bell) or smooth with non-negative Jacobian and isolated singularities [Aki99].

David Bellamy [Bel79] produced an important related counterexample. He showed that there exists a tree-like continuum $X$, whose every proper subcontinuum is an arc and which admits a fixed point free homeomorphism. It is not known if examples of this type can be embedded in the plane. Minc [Min99] constructed a tree-like continuum which is the continuous image of the pseudo arc and admits a fixed point free map.

**1.0.3. Major tools.** In this subsection we describe the major tools developed in Part 1.

**1.0.3.1. Finding fixed points with index and variation.** It is easy to see that a map of a plane continuum to itself can be extended to a perfect map of the plane. We study the slightly more general question, “Is there a plane continuum $Z$ and a perfect continuous function $f : \mathbb{C} \to \mathbb{C}$ taking $Z$ into $T(Z)$ with no fixed points in $T(Z)$?” A Zorn’s Lemma argument shows that if one assumes that the answer is “yes,” then there is a subcontinuum $X \subset Z$, minimal with respect to these properties. It will follow from Theorem 6.1.4 that for such a minimal continuum, $f(X) = X = \partial T(X)$ (though it may not be the case that $f(T(X)) \subset T(X)$). Here $\partial T(X)$ denotes the boundary of $T(X)$.

Many fixed point results make use of the notion of the index $\text{ind}(f, S)$, which counts the number of revolutions of the vector connecting $z$ with $f(z)$ for $z \in S$ running along a simple closed curve $S$ in the plane. As is well-known, if $f : \mathbb{C} \to \mathbb{C}$ is a map and $\text{ind}(f, S) \neq 0$, then $f$ must have a fixed point in $T(S)$ (for completeness we prove this in Theorem 3.1.4). In order to establish fixed points in invariant plane continua $X$, one often approximates $X$ by a simple closed curve $S$ such that $X \subset T(S)$. If $\text{ind}(f, S) \neq 0$ and $S$ is sufficiently tight around $X$, one can conclude that $f$ must have a fixed point in $T(X)$. Hence the main work is in showing that $\text{ind}(f, S) \neq 0$ for a suitable simple closed curve around $X$.

Bell’s fundamental idea was to replace the count of the number of rotations of the vector $zf(z)$ with respect to a fixed axis (say, the $x$-axis) by a count which involves the moving frame of external rays. Consider, for example, the unit circle.
Let $K$ be a maximal round closed ball (or a half plane) such that $\text{int}(B)$ has the property that any two distinct chords meet in at most a common endpoint in $E$. The Linchpin Theorem (see Theorem 4.2.5 and the remark following it) states that the collection $\text{conv}_E L(E)$ contains a fixed point as desired. Observe that the connection between index and variation is essential for all the applications in Part 2.

1.0.3.2. Other tools, such as foliations and oriented maps. Let $X$ be a non-separating plane continuum and let $f : C \to C$ be a map such that $f(X) \subset T(X)$ and $f$ has no fixed points in $T(X)$. In general we need more control of the simple closed curve $S$ around $X$ (and of the action of the map on $S \setminus X$). Bell originally accomplished this by partitioning the complement of $T(X)$ in the Euclidean convex hull $\text{conv}_E(X)$ of $X$ by Euclidean convex sets. Suppose that $B$ is a maximal round closed ball (or a half plane) such that $\text{int}(B) \cap T(X) = \emptyset$ and $|B \cap X| \geq 2$, and consider the set $\text{conv}_E(B \cap X)$. For any two such balls $B_1, B_2$ either $K_{1,2} = \text{conv}_E(B_1 \cap X) \cap \text{conv}_E(B_2 \cap X)$ is empty, or $K_{1,2}$ is a single point in $X$, or this intersection is a common chord contained in both of their boundaries. Bell’s Linchpin Theorem (see Theorem 4.2.5 and the remark following it) states that the collection $\text{conv}_E(B \cap X)$ over all such maximal balls covers all of $\text{conv}_E(X) \setminus T(X)$.

Hence the collection $\text{conv}_E(B \cap X)$ over all such balls provides a partition of $\text{conv}_E(X) \setminus X$ into Euclidean convex sets contained in maximal round balls. The collection of chords in the boundaries of the sets $\text{conv}_E(B \cap X)$ for all such balls have the property that any two distinct chords meet in at most a common endpoint in $X$. In other words, this set of chords is a lamination in the sense of Thurston [Thu09] even though in Thurston’s paper laminations appear in a very different, namely complex dynamical, context. This Linchpin Theorem can be used to extend the map $f|_{T(X)}$ over $\text{conv}_E(X) \setminus T(X)$ (first linearly over all the chords in the lamination and then over all remaining components of the complement). We will illustrate the usefulness of Bell’s partition by showing that the well-known Schoenflies Theorem follows immediately. It can also be used to obtain a particular simple closed curve $S$ around $X$ so that every component of $S \setminus X$ is a chord in the lamination.

In our version of Bell’s Linchpin Theorem we consider arbitrary open and connected subsets of the sphere $U$ and we use round balls in the spherical metric on $S^1$ and a fixed point free map $f : S^1 \to C$. For each $z = e^{2 \pi i \theta} \in S^1$ let $R_z$ be the external ray $\{re^{2 \pi i \theta} : r > 1\}$. Now count the number of times the point $f(z)$, $z \in S^1$, crosses the external ray $R_z$, taking into account the direction of the crossing. Call this count the variation $\text{var}(f, S^1)$. It is easy to see that in this case $\text{ind}(f, S^1) = \text{var}(f, S^1) + 1$.
the sphere. Moreover, the Euclidean geodesics are replaced by hyperbolic geodesics (in either the hyperbolic metric in each ball or, if $U$ is simply connected, in the hyperbolic metric on $U$). This way we get a lamination of all of $U^\infty(X) = C^\infty \setminus T(X)$ (and not just $\text{conv}_X(X) \setminus T(X)$) and the resulting lamination is easier to apply in other settings. We give a proof of this theorem using an elegant argument due to Kulkarni and Pinkall $[\text{KP}94]$ which also allows for the extension over arbitrary open and connected subsets of the sphere (see Theorem 4.1.5; this later theorem is used in Chapter 6). Bell’s Linchpin Theorem follows as a corollary.

This new partition of a complementary domain of a continuum can also be used in other settings to extend a homeomorphism on the boundary of a planar domain over the entire domain. In $[\text{OT}07]$ this is used to show that an isotopy of a planar continuum, starting at the identity, extends to an isotopy of the plane. This extends a well-known result regarding the extension of a holomorphic motion $[\text{ST}86, \text{Slo}91]$. In $[\text{OV}09]$ this partition is used to give necessary and sufficient conditions to extend a homeomorphism, of an arbitrary planar continuum, over the plane.

The development of the necessary tools in Part 1 is completed by introducing the notion of (positively or negatively) oriented maps and studying their properties. Holomorphic maps are prototypes of positively oriented maps but in general positively oriented maps do not have to be differentiable, light, open or monotone. Locally at non-critical points, positively oriented maps behave like orientation-preserving homeomorphisms in the sense that they preserve local orientation. Compositions of open, perfect and of monotone, perfect surjections of the plane are confluent (i.e., such that components of the preimage of any continuum map onto the continuum) and naturally decompose into two classes, one of which preserves and the other of which reverses local orientation. We show that any confluent map of the plane is itself a composition of a monotone and a light-open map of the plane. It is shown that an oriented map of the plane induces a map from the circle of prime ends of a component of the pre-image of an acyclic plane continuum to the circle of prime ends of that continuum.

1.0.4. Main applications. Part 2 contains applications of the tools developed in Part 1. Directly or indirectly, these applications deal with the Plane Fixed Point Problem. We describe them below.

1.0.4.1. Outchannel and hypothetical minimal continua without fixed points. The first application is in Chapter 6 where we establish the existence of a unique outchannel. Let us consider this in more detail. If there exists a counterexample to the Plane Fixed Point Problem, then there exists a continuum $X$ which is minimal with respect to $f(X) \subset T(X)$ and $f$ has no fixed point in $T(X)$. Bell has shown $[\text{Bel}67]$ (see also $[\text{Sie}68, \text{Ili}70]$) that such a continuum has at least one dense outchannel of negative variation. Since $X$ is minimal with a dense outchannel, $X$ is an indecomposable continuum and $f(X) = X$.

A dense outchannel is a prime end so that its principal set is all of $X$ and if $\{C_i\}$ is a defining sequence of crosscuts, then $f$ maps these crosscuts essentially “out of the channel” (i.e., closer to infinity) for $i$ sufficiently large. The latter statement is accurately reflected by the fact that $\text{var}(f, C_i) \neq 0$. In case that the complement of $X$ is invariant, this can be described by saying that the crosscut $f(C_i)$ separates $C_i$ from infinity in $U^\infty(X)$ (and hence, in this case, crosscuts do really map out of the channel). As a new result, the main steps of the proof of which were outlined by
1. INTRODUCTION

Bell, we show that there always exists exactly one outchannel and that its variation is $-1$, while all other prime ends must have variation 0. Using these results it is shown in [BO09] that if $f$ is a negatively oriented branched covering map of degree 2 which has a non-separating invariant continuum $Z$ without fixed points, then the minimal subcontinuum $X \subset Z$ has the following additional properties:

1. $X$ is indecomposable,
2. the unique critical point $c$ of $f$ belongs to $X$, $f(X) = X$, and $f(\mathbb{C} \setminus X) = \mathbb{C} \setminus X$ (so that $X$ is fully invariant),
3. $f$ induces a covering map $F: S^1 \to S^1$ from the circle of prime ends of $X$ to itself of degree $-2$.
4. $F$ has three fixed points one of which corresponds to the unique dense outchannel whereas the remaining two fixed points correspond to dense inchannels (i.e., for a defining sequence of crosscuts $\{C_i\}$, $C_i$ separates $f(C_i)$ from infinity in $U^\infty$).

Moreover, as part of the argument, the map $f$ is modified in $\mathbb{C} \setminus X$ so that the new map $g$ keeps the tail of the external ray, which runs down the outchannel, invariant and maps the points on them closer to infinity.

1.0.4.2. Fixed points in invariant continua for positively oriented maps. Other applications of the tools developed in Part 1 are obtained in Chapter 7. These are also related to the Plane Fixed Point Problem. As we will see below, the corresponding results can be in turn further applied in complex dynamics, leading to some structural results in the field, such as constructing finest locally connected models for connected Julia sets or studying wandering continua inside Julia sets and an extension of the Fatou-Shishikura inequality so that it includes counting wandering branch-continua (see Section 7.5).

The first application in Chapter 7 is the most straightforward of them all: in Theorem 7.1.3 from Section 7.1 we prove that a positively oriented map $f$ which takes a continuum $X$ into the topological hull $T(X)$ of $X$ must have a fixed point in $T(X)$. In other words, in Theorem 7.1.3 the Plane Fixed Point Problem is solved in the affirmative for positively oriented maps. As we will see, the extension from holomorphic maps to positively oriented maps is important since the latter class allows for easy local perturbations. This will allow us to deal with parabolic points in a Julia set (see Lemma 7.5.1, Theorem 7.5.2 and Corollary 7.5.4).

The idea of the proof is as follows. First we prove in Corollary 7.1.2 that if a crosscut $C$ of $X$ is mapped off itself by $f$ then the variation on $C$ is non-negative. This is done by completing the crosscut $C$ to a very tight simple closed curve $S$ around $X$ and observing that in fact the variation in question can be computed by computing the winding number of $f$ on $S$. Notice that versions of this idea are used later on when we prove the existence of fixed points in non-invariant continua satisfying certain additional conditions.

To prove Theorem 7.1.3, we first assume by way of contradiction that $T(X)$ contains no fixed points. In this case there are no fixed points in the closure $\overline{U}$ of a sufficiently small neighborhood $U$ of $T(X)$. Using this, we construct a simple closed curve which goes around $X$ inside such a neighborhood $U$ and “touches” $X$ at a sufficiently dense set of points so that arcs between consecutive points of $S \cap X$ are very small. Since there are no fixed points in $\overline{U}$, we can guarantee that the images of these arcs are disjoint from themselves. Hence by the above described Corollary 7.1.2 the variations of all these arcs are non-negative. By Theorem 3.2.2
1. INTRODUCTION

this implies that the index of $f$ on $S$ is not equal to zero and hence, by Theorem 3.1.4 there must exist a fixed point inside $T(S)$, a contradiction.

1.0.4.3. Fixed points in non-invariant continua: the case of dendrites. Our generalizations of Theorem 7.1.3 are inspired by a simple observation. The most well-known particular case for which the Plane Fixed Point Problem is solved is that of a map of a closed interval $I = [a, b]$, $a < b$ into itself in which case there must exist a fixed point in $I$. However, in this case a more general result can easily be proven, of which the existence of a fixed point in an invariant interval is a consequence.

Namely, instead of considering a map $f : I \rightarrow I$ consider a map $f : I \rightarrow \mathbb{R}$ such that either (a) $f(a) \geq a$ and $f(b) \leq b$, or (b) $f(a) \leq a$ and $f(b) \geq b$. Then still there must exist a fixed point in $I$ which is an easy corollary of the Intermediate Value Theorem applied to the function $f(x) - x$. Observe that in this case $I$ need not be invariant under $f$. Observe also that without the assumptions on the endpoints, the conclusion on the existence of a fixed point inside $I$ cannot be made because, e.g., a shift map on $I$ does not have fixed points at all. The conditions (a) and (b) above can be thought of as boundary conditions imposing restrictions on where $f$ maps the boundary points of $I$ in $\mathbb{R}$.

Our main aim in the remaining part of Chapter 7 is to consider some other cases for which the Plane Fixed Point Problem can be solved in the affirmative (i.e., the existence of a fixed point in a continuum can be established) despite the fact that the continuum $X$ in question is not invariant. We proceed with our studies in two directions. Considering $X$, we replace the invariantness of the continuum by boundary conditions in the spirit of the above “interval version” of the Plane Fixed Point Problem. We also show that there must exist a fixed point of “rotational type” in the continuum (and hence, if it is known that such a point does not exist, then the continuum in question is a point).

Since we now deal with continua significantly more complicated than an interval, inevitably the boundary conditions become rather intricate. Thus we postpone the precise technical statement of the results until Chapter 7 and use here a more descriptive approach. Observe that particular cases for which the Plane Fixed Point Problem is solved so far can be divided into two categories: either $X$ has additional properties, or $f$ has additional properties. In the first category the above considered “interval case” is the most well-known. A direct extension of it is the following well-known theorem (which follows from Borsuk’s theorem [Bor35], see [Nad92] for a direct proof): recall that a dendrite is a locally connected continuum containing no simple closed curves.

Theorem 1.0.2. If $f : D \rightarrow D$ is a continuous map of a dendrite into itself then it has a fixed point.

Here $f$ is just a continuous map but the continuum $D$ is very nice. In Section 7.2 Theorem 1.0.2 is generalized to the case when $f : D_1 \rightarrow D_2$ maps a dendrite $D_1$ into a dendrite $D_2 \supset D_1$ and certain conditions on the behavior of the points of the set $E = D_2 \backslash D_1 \cap D_1$ under the map $f$ are fulfilled (observe that $E$ may be infinite). This presents a “non-invariant” version of Plane Fixed Point Problem for dendrites and can be done in the spirit of the interval case described earlier. Moreover, with some additional conditions it has consequences related to the number of periodic points of $f$.

More precisely, we introduce the notion of boundary scrambling for dendrites in the situation above. It simply means that for each non-fixed point $e \in E$, $f(e)$
is contained in a component of $D_2 \setminus \{e\}$ which intersects $D_1$ (see Definition 5.3.1). Observe that if $D_1$ is invariant then $f$ automatically scrambles the boundary. We prove the following theorem.

**Theorem 7.2.2.** Suppose that $f : D_1 \rightarrow D_2$ is a map between dendrites, where $D_1 \subset D_2$, which scrambles the boundary. Then $f$ has a fixed point.

Next in Section 7.2 we define *weakly repelling periodic points*. Basically, a point $a \in D_1$ is a *weakly repelling periodic point* (for $f^n$) if there exists $n \geq 1$ and a component $B$ of $D_1 \setminus \{a\}$ such that $f^n(a) = a$ and arbitrarily close to $a$ in $B$ there exist cutpoints of $D_1$ fixed under $f^n$ or points $x$ separating $a$ from $f^n(x)$. Note that a fixed point $a$ of $f$ can be a weakly repelling periodic point for $f^n$ while it is not weakly repelling for $f$. We use this notion to prove Theorem 7.2.6 where we show that if $D$ is a dendrite and $f : D \rightarrow D$ is continuous and all its periodic points are weakly repelling, then $f$ has infinitely many periodic cutpoints. Then we rely upon Theorem 7.2.6 in Theorem 7.2.7 where it is shown that if $g : J \rightarrow J$ is a *topological polynomial* on its dendritic Julia set (e.g., if $g$ is a complex polynomial with a dendritic Julia set) then it has infinitely many periodic cutpoints.

1.0.4.1. Fixed points in non-invariant continua: the planar case. In Sections 7.3 and 7.4 we draw a parallel with the interval case for planar maps and extend Theorem 7.1.3 to non-invariant continua under positively oriented maps such that certain “boundary” conditions are satisfied. Namely, suppose that $f : C \rightarrow C$ is a positively oriented map and $X \subset C$ is a non-separating continuum. Since we are interested in fixed points of $f|_X$, it makes sense to assume that at least $f(X) \cap X \neq \emptyset$. Thus, we can think of $f(X)$ as a new continuum which “grows” from $X$ at some places. We assume that the “pieces” of $f(X)$ which grow outside $X$ are contained in disjoint non-separating continua $Z_i$ so that $f(X) \setminus X \subset \bigcup Z_i$.

We also assume that places at which the growth takes places - i.e., sets $Z_i \cap X = K_i$ - are non-separating continua for all $i$. Finally, the main assumption here is the following restriction upon where the continua $K_i$ map under $f$: we assume that for all $i$, $f(K_i) \cap \big[ Z_i \setminus K_i \big] = \emptyset$. If this is all that is satisfied, then the map $f$ is said to *scramble the boundary* (of $X$). A stronger version of that is when for all $i$, either $f(K_i) \subset K_i$, or $f(K_i) \cap Z_i = \emptyset$; then we say that $f$ *strongly scrambles the boundary* (of $X$) (see Definition 5.4.1). The continua $K_i$ are called *exit continua* (of $X$). The main result of Section 7.3 is:

**Theorem 7.3.3.** Suppose that $f$ is positively oriented and strongly scrambles the boundary of $X$, then $f$ has a fixed point in $X$.

As an illustration, consider the case when $X \cup (\bigcup Z_i)$ is a dendrite and all sets $K_i$ are singletons. Then it is easy to see that both scrambling and strong scrambling of the boundary in the sense of dendrites mean the same as in the sense of the planar definition. Of course, in the planar case we deal with a much more narrow class of maps, namely positively oriented maps, and with a much wider variety of continua, namely all non-separating planar continua. This fits into the “philosophy” of our approach: whenever we obtain a result for a wider class of continua, we have to consider a more specific class of maps.

For the family of positively oriented maps with isolated fixed points we specify this result as follows. We introduce the notion of the map $f$ *repelling outside $X$ at a fixed point $p$* (see Definition 7.4.5; basically, it means that there exists an invariant external ray of $X$ which lands at $p$ and along which the points are repelled away
1. INTRODUCTION

From \( p \). Then in Theorem 7.4.7 we show that if \( f \) is a positively oriented map with isolated fixed points and \( X \subset \mathbb{C} \) is a non-separating continuum or a point such that \( f \) scrambles the boundary of \( X \) and for every fixed point \( a \) the winding number at \( a \) equals 1 and \( f \) repels at \( a \), then \( X \) must be a point.

1.0.4.5. **Fixed points in non-invariant continua for polynomials.** These theorems apply to polynomials \( P \), allowing us to obtain a few corollaries dealing with the existence of periodic points in certain parts of the Julia set of a polynomial and degeneracy of certain impressions. To discuss this we assume knowledge of the standard definitions such as Julia sets \( J_P \), filled-in Julia sets \( K_P = \overline{T(J_P)} \), Fatou domains, parabolic periodic points etc which are formally introduced in Section 5.5 and further discussed in Section 7.5 (see also [Mil00]). Recall that the set \( U^\infty(J_P) \) (called in this context the *basin of attraction of infinity*) is partitioned by the (conformal) external rays \( R_\alpha \) with arguments \( \alpha \in \mathbb{S}^1 \). If \( J_P \) is connected, all rays \( R_\alpha \) are smooth and pairwise disjoint while if \( J_P \) is not connected limits of smooth external rays must be added. Still, given an external ray \( R_\alpha \) of \( K_P \), its principal set \( R_\alpha \setminus R_\alpha \) can be introduced as usual.

We then define a *general puzzle-piece* of a filled-in Julia set \( K_P \) as a continuum \( X \) which is cut from \( K_P \) by means of choosing a few exit continua \( E_i \subset X \) each of which contains the principal sets of more than one external ray. We then assume that there exists a component \( C_X \) of the complement in \( \mathbb{C} \) to the union of all such exit continua \( E_i \) and their external rays such that \( X \subset (C_X \cap K_P) \cup (\bigcup E_i) \). The external rays accumulating inside an exit continuum \( E_i \) cut the plane into wedges one of which, denoted by \( W_i \), contains points of \( X \). The "degenerate" case when there are no exit continua is also included and simply means that \( X \) is an invariant subcontinuum of \( K_P \).

The main assumption on the dynamics of a general puzzle piece \( X \) which we make is that \( P(X) \cap C_X \subset X \) and for each exit continuum \( E_i \) we have \( P(E_i) \subset W_i \). It is easy to see that this essentially means that \( P \) scrambles the boundary of \( X \) (where the role of the "boundary" is played by the union of exit continua).

The conclusion, obtained in Theorem 7.5.2, is based upon the above described results, in particular on Theorem 7.4.7. It states that for a general puzzle-piece either \( X \) contains an invariant parabolic Fatou domain, or \( X \) contains a fixed point which is neither repelling nor parabolic, or \( X \) contains a repelling or parabolic fixed point \( a \) at which the local rotation number is not 0. Let us now list the main dynamical applications of this result.

1.0.4.6. **Further dynamical applications.** There are a few ways Theorem 7.5.2 applies in complex (polynomial) dynamics. First, it is instrumental in studying wandering cut-continua for polynomials with connected Julia sets. A continuum/point \( L \subset J_P \) is a cut-continuum (of valence \( \text{val}(L) \)) if the cardinality \( \text{val}(L) \) of the set of components of \( J_P \setminus L \) is greater than 1. A collection of disjoint cut-continua (it might, in particular, consist of one continuum) is said to be wandering if their forward images form a family of pairwise disjoint sets. The main result of [BCLOS08] in the case of polynomials with connected Julia sets is the following generalization of the Fatou-Shishikura inequality.

**Theorem 1.0.3.** Let \( P \) be a polynomial with connected Julia set, let \( N \) be the sum of the number of distinct cycles of its bounded Fatou domains and the number of cycles of its Cremer points, and let \( \Gamma \neq \emptyset \) be a wandering collection of
cut-continua $Q_i$ with valences greater than 2 which contain no preimages of critical points of $P$. Then $\sum_i (\text{val}(Q_i) - 2) + N \leq d - 2$.

In [BCLOS08] a partition of the plane into pieces by rays with rational arguments landing at periodic cutpoints of $J_P$ and their preimages is used. Theorem 7.5.2 plays a significant role in the proof of the fact that wandering cut-continua do not enter the pieces containing Cremer or Siegel periodic points which is an important ingredient of the arguments in [BCLOS08] proving Theorem 1.0.3.

Another application of Theorem 7.5.2 can be found in [BCO08] where Kiwi’s fundamental result [Kiw04] on the semiconjugacy of polynomials on their Julia sets without Cremer or Siegel points is extended to all polynomials with connected Julia sets; in both cases topological polynomials on their topological Julia sets serve as locally connected models. Denote the monotone semiconjugacy in question by $\varphi$.

In showing in [BCO08] that if $x$ is a (pre)periodic point and $\varphi(x)$ is not equal to a $\varphi$-image of a Cremer or Siegel point or its preimage then $J_P$ is locally connected at $x$, Theorem 7.5.2 plays a crucial role.

Finally, our results concerning dendrites (such as Theorem 7.2.6 and Theorem 7.2.7) are used in [BCO08] where a criterion for the connected Julia set to have a non-degenerate locally connected model is obtained. We also rely on Theorem 7.2.6 and Theorem 7.2.7 to show in [BCO08] that if such model exists, and is a dendrite, then the polynomial must have infinitely many bi-accessible periodic points in its Julia set.

1.0.5. Concluding remarks and acknowledgments. All of the positive results on the existence of fixed points in this memoir are either for simple continua (i.e., those which do not contain indecomposable subcontinua) or for positively oriented maps of the plane. Hence the following special case of the Plane Fixed Point Problem is a major remaining open problem:

Problem 1.0.4. Suppose that $f : \mathbb{C} \to \mathbb{C}$ is a negatively oriented branched covering map, $|f^{-1}(y)| \leq 2$ for all $y \in \mathbb{C}$ and $Z$ is a non-separating plane continuum such that $f(Z) \subset Z$. Must $f$ have a fixed point in $Z$?

Suppose that $c$ is the unique critical point of $f$ and that $X \subset Z$ is a minimal continuum such that $f(X) \subset T(X)$. Then, as was mentioned above, the answer is yes if there exists $y \in X \setminus \{f(c)\}$ such that $|f^{-1}(y) \cap X| < 2$. In particular the answer to Problem 1.0.4 is yes if $f|_Z$ is one-to-one.

Finally let us express, once again, our gratitude to Harold Bell for sharing his insights with us. His notion of variation of an arc, his index equals variation plus one theorem and his linchpin theorem of partitioning a complementary domain of a planar continuum into convex subsets are essential for the results we obtain here. Theorem 6.2.1 (Unique Outchannel) is a new result the main steps of which were outlined by Bell. Complete proofs of the following results by Bell: Theorems 3.2.2, 3.3.1, 4.2.5 and 6.2.1, appear in print for the first time. For the convenience of the reader we have included an index at the end of the paper.
Part 1

Basic Theory
CHAPTER 2

Preliminaries and outline of Part 1

In this chapter we give the formal definitions and describe the results of part 1 in more detail. By a map \( f : X \to Y \) we will always mean a continuous function.

Let \( p : \mathbb{R} \to S^1 \) denote the covering map \( p(x) = e^{2\pi i x} \). Let \( g : S^1 \to S^1 \) be a map. By the degree of the map \( g \), denoted by \( \text{degree}(g) \), we mean the number \( \hat{g}(1) - \hat{g}(0) \), where \( \hat{g} : \mathbb{R} \to \mathbb{R} \) is a lift of the map \( g \) to the universal covering space \( \mathbb{R} \) of \( S^1 \) (i.e., \( p \circ \hat{g} = g \circ p \)). It is well-known that \( \text{degree}(g) \) is independent of the choice of the lift.

2.1. Index

Let \( g : S^1 \to \mathbb{C} \) be a map and \( f : g(S^1) \to \mathbb{C} \) a fixed point free map. Define the map \( v : S^1 \to S^1 \) by

\[
v(t) = f(g(t)) - g(t).
\]

Then the map \( v : S^1 \to S^1 \) lifts to a map \( \tilde{v} : \mathbb{R} \to \mathbb{R} \). Define the index of \( f \) with respect to \( g \), denoted \( \text{ind}(f, g) \) by

\[
\text{ind}(f, g) = \tilde{v}(1) - \tilde{v}(0) = \text{degree}(v).
\]

Note that \( \text{ind}(f, g) \) measures the net number of revolutions of the vector \( f(g(t)) - g(t) \) as \( t \) travels through the unit circle one revolution in the positive direction.

Remark 2.1.1. The following basic facts hold.

(a) If \( g : S^1 \to \mathbb{C} \) is a constant map with \( g(S^1) = c \) and \( f(c) \neq c \), then \( \text{ind}(f, g) = 0 \).

(b) If \( f \) is a constant map and \( f(\mathbb{C}) = w \) with \( w \notin g(S^1) \), then \( \text{ind}(f, g) = \text{win}(g, S^1, w) \), the winding number of \( g \) about \( w \). In particular, if \( f : S^1 \to T(S^1) \setminus S^1 \) is a constant map, then \( \text{ind}(f, \text{id}|_{S^1}) = 1 \), where \( \text{id}|_{S^1} \) is the identity map on \( S^1 \).

Note also, that for a simple closed curve \( S' \) and a point \( w \notin T(f(S')) \) we have \( \text{win}(f, S', w) = 0 \). Suppose \( S \subset \mathbb{C} \) is a simple closed curve and \( A \subset S \) is a subarc of \( S \) with endpoints \( a \) and \( b \). Then we write \( A = [a, b] \) if \( A \) is the arc obtained by traveling in the clockwise direction from the point \( a \) to the point \( b \) along \( S \). In this case we denote by \( < \) the linear order on the arc \( A \) such that \( a < b \). We will call the order \( < \) the clockwise order on \( A \). Note that \( [a, b] \neq [b, a] \).

More generally, for any arc \( A = [a, b] \subset S^1 \), with \( a < b \) in the clockwise order, define the fractional index \([\text{Bro90}]\) of \( f \) on the sub-path \( g|[a,b] \) by

\[
\text{ind}(f, g|[a,b]) = \tilde{v}(b) - \tilde{v}(a).
\]
While, necessarily, the index of $f$ with respect to $g$ is an integer, the fractional index of $f$ on $g_{[a,b]}$ need not be. We shall have occasion to use fractional index in the proof of Theorem 3.2.2.

**Proposition 2.1.2.** Let $g: S^1 \to \mathbb{C}$ be a map with $g(S^1) = S$, and suppose $f: S \to \mathbb{C}$ has no fixed points on $S$. Let $a \neq b \in S^1$ with $[a, b]$ denoting the counterclockwise subarc on $S^1$ from $a$ to $b$ (so $[a, b]$ and $(b, a)$ are complementary arcs and $S^1 = [a, b] \cup [b, a]$). Then $\text{ind}(f, g) = \text{ind}(f, g_{[a,b]}) + \text{ind}(f, g_{[b,a]}).

**2.2. Variation**

In this section we introduce the notion of variation of a map on an arc and relate it to winding number.

**Definition 2.2.1 (Junctions).** The standard junction $J_O$ is the union of the three rays $J^+_O = \{z \in \mathbb{C} \mid z = re^{i\pi/2}, \ r \in [0, \infty)\}$, $J^+_O = \{z \in \mathbb{C} \mid z = re^{-i\pi}, \ r \in [0, \infty)\}$, $J^+_O = \{z \in \mathbb{C} \mid z = re^{i\pi}, \ r \in [0, \infty)\}$, having the origin $O$ in common. A junction (at $v$) $J_v$ is the image of $J_O$ under any orientation-preserving homeomorphism $h: \mathbb{C} \to \mathbb{C}$ where $v = h(O)$. We will often suppress $h$ and refer to $h(J_O)$ as $J_v$, and similarly for the remaining rays in $J_v$. Moreover, we require that for each bounded neighborhood $W$ of $v$, $d(J^+_v \setminus W, J^-_v \setminus W) > 0$.

**Definition 2.2.2 (Variation on an arc).** Let $S \subseteq \mathbb{C}$ be a simple closed curve, $f: S \to \mathbb{C}$ a map and $A = [a, b]$ a subarc of $S$ such that $f(a), f(b) \in T(S)$ and $f(A) \cap A = \emptyset$. We define the variation of $f$ on $A$ with respect to $S$, denoted $\text{var}(f, A, S)$, by the following algorithm:

1. Let $v \in A$ and let $J_v$ be a junction with $J_v \cap S = \{v\}$.
2. Counting crossings: Consider the set $M = f^{-1}(J_v) \cap [a, b]$. Each time a point of $f^{-1}(J_v^+) \cap [a, b]$ is immediately followed in $M$, in the counterclockwise order $< \{a, b\} \subset S$, by a point of $f^{-1}(J_v^-)$, count $+1$ and each time a point of $f^{-1}(J_v^-) \cap [a, b]$ is immediately followed in $M$ by a point of $f^{-1}(J_v^+)$, count $-1$. Count no other crossings.
3. The sum of the crossings found above is the variation $\text{var}(f, A, S)$.

Note that $f^{-1}(J_v^+) \cap [a, b]$ and $f^{-1}(J_v^-) \cap [a, b]$ are disjoint closed sets in $[a, b]$. Hence, in (2) in the above definition, we count only a finite number of crossings and $\text{var}(f, A, S)$ is an integer. Of course, if $f(A)$ does not meet both $J_v^+$ and $J_v^-$, then $\text{var}(f, A, S) = 0$.

If $\alpha: S \to \mathbb{C}$ is any map such that $\alpha|_A = f|_A$ and $\alpha(S \setminus (a, b)) \cap J_v = \emptyset$, then $\text{var}(f, A, S) = \text{win}(\alpha, S, v)$. In particular, this condition is satisfied if $\alpha(S \setminus (a, b)) \subset T(S) \setminus \{v\}$. The invariance of winding number under suitable homotopies implies that the variation $\text{var}(f, A, S)$ also remains invariant under such homotopies. That is, even though the specific crossings in (2) in the algorithm may change, the sum remains invariant. We will state the required results about variation below without proof. Proofs can be obtained directly by using the fact that $\text{var}(f, A, S)$ is integer-valued and continuous under suitable homotopies.

**Proposition 2.2.3 (Junction Straightening).** Let $S \subseteq \mathbb{C}$ be a simple closed curve, $f: S \to \mathbb{C}$ a map and $A = [a, b]$ a subarc of $S$ such that $f(a), f(b) \in T(S)$ and $f(A) \cap A = \emptyset$. Any two junctions $J_w$ and $J_u$ with $u, v \in A$ and $J_w \cap S = \{w\}$ for $w \in \{u, v\}$ give the same value for $\text{var}(f, A, S)$. Hence $\text{var}(f, A, S)$ is independent of the particular junction used in Definition 2.2.2.
Let $J_v$ coming from a proper compact subarc of the open arc $(a,b)$. Consequently, $\text{var}(f, A, S)$ remains invariant under homotopies $h_t$ of $f|_{[a,b]}$. In the complement of $[v]$ such that $h_t(a), h_t(b) \notin J_v$ for all $t$. Moreover, the computation is stable under an isotopy $h_t : J_v \rightarrow A \cup [\mathbb{C} \setminus T(S)]$ that moves the entire junction $J_v$ (even off $A$), provided that during the isotopy $h_t(v) \notin f(A)$ and $f(a), f(b) \notin h_t(J_v)$ for all $t$.

In case $A$ is an open arc $(a,b) \subset S$ such that $\text{var}(f, \mathcal{A}, S)$ is defined, it will be convenient to denote $\text{var}(f, \mathcal{A}, S)$ by $\text{var}(f, A, S)$.

The following lemma follows immediately from the definition.

**Lemma 2.2.4.** Let $S \subset \mathbb{C}$ be a simple closed curve. Suppose that $a < c < b$ are three points in $S$ such that $\{f(a), f(b), f(c)\} \subset T(S)$ and $f([a,b]) \cap [a,b] = \emptyset$. Then $\text{var}(f, [a, b], S) = \text{var}(f, [a, c], S) + \text{var}(f, [c, b], S)$.

**Definition 2.2.5 (Variation on a finite union of arcs).** Let $S \subset \mathbb{C}$ be a simple closed curve and $A = [a,b]$ a subcontinuum of $S$ partitioned by a finite set $F = \{a = a_0 < a_1 < \cdots < a_n = b\}$ into subarcs. For each $i$ let $A_i = [a_i, a_{i+1}]$. Suppose that $f$ satisfies $f(a_i) \in T(S)$ and $f(A_i) \cap A_i = \emptyset$ for each $i$. We define the variation of $f$ on $A$ with respect to $S$, denoted $\text{var}(f, A, S)$, by

$$\text{var}(f, A, S) = \sum_{i=0}^{n-1} \text{var}(f, [a_i, a_{i+1}], S).$$

In particular, we include the possibility that $a_n = a_0$ in which case $A = S$.

By considering a common refinement of two partitions $F_1$ and $F_2$ of an arc $A \subset S$ such that $f(F_1) \cup f(F_2) \subset T(S)$ and satisfying the conditions in Definition 2.2.5, it follows from Lemma 2.2.4 that we get the same value for $\text{var}(f, A, S)$ whether we use the partition $F_1$ or the partition $F_2$. Hence, $\text{var}(f, A, S)$ is well-defined. If $A = S$ we denote $\text{var}(f, S, S)$ simply by $\text{var}(f, S)$.

The first main result in Part 1, Theorem 3.2.2 is that given a map $f : \mathbb{C} \rightarrow \mathbb{C}$, a simple closed curve $S \subset \mathbb{C}$ and a partition of $S$ into subarcs $A_i$, such that any two meet at most in a common endpoint, for each $i$ $f(A_i) \cap A_i = \emptyset$ and both endpoints map into $T(S)$,

$$\text{ind}(f, S) = \sum \text{var}(f, A_i) + 1.$$

In the first version of this theorem we partition $S$ into finitely many subarcs $A_i$. We extend this in Section 3.5 by allowing partitions of $S$ which consist of, possibly countably infinitely many subarcs. Since in our applications we often assume that we have an invariant continuum $X$ such that $f$ has no fixed point in $T(X)$ it follows from Theorem 3.1.4 that, for a sufficiently tight simple closed curve $S$ around $X$ with $X \subset T(S)$, we must have $\text{ind}(f, S) = 0$. It follows from the above theorem relating index and variation that for some subarc $A$ (which is the closure of a component of $S \setminus X$ and, hence a crosscut of $X$), $\text{var}(f, A) < 0$. In order to locate this crosscut of negative variation we establish Bell’s Lollipop Theorem in Section 3.3.

### 2.3. Classes of maps

Cartwright and Littlewood solved the Plane Fixed Point Problem for orientation preserving homeomorphism of the plane. In Section 3.7 we introduce and study (positively) oriented maps of the plane. We will show in Part 2 that the
Plane Fixed Point Problem has a positive solution for the class of positively oriented maps. We show in Section 3.7 that the class of oriented maps consists of all compositions of monotone and open perfect maps of the plane and that all such maps are confluent. In particular, analytic maps are confluent.

Let us begin by listing a few well-known definitions.

**Definition 2.3.1.** A perfect map is a closed continuous function each of whose point inverses is compact. We will assume in the remaining sections that all maps of the plane considered in this memoir are perfect. Let $X$ and $Y$ be topological spaces. A map $f : X \to Y$ is monotone provided for each continuum $K \subseteq Y$, $f^{-1}(K)$ is connected and $f$ is light provided for each point $y \in Y$, $f^{-1}(y)$ is totally disconnected. A map $f : X \to Y$ is confluent provided for each continuum $K \subseteq Y$ and each component $C$ of $f^{-1}(K)$, $f(C) = K$. Every map $f : X \to Y$ between compacta is the composition $f = l \circ m$ of a a monotone map $m : X \to Z$ and a light map $l : Z \to Y$ for some compactum $Z$ [Nad92, Theorem 13.3]. This representation is called the monotone-light decomposition of $f$.

Observe that any confluent map $f$ is onto. It is well-known that each homeomorphism of the plane is either orientation-preserving or orientation-reversing. We will establish an appropriate extension of this result for confluent perfect mappings of the plane (Theorem 3.7.4) by showing that such maps either preserve or reverse local orientation. As a consequence it follows that all perfect and confluent maps of the plane satisfy the Maximum Modulus Theorem. We will call such maps positively- or negatively oriented maps, respectively.

Complex polynomials $P : \mathbb{C} \to \mathbb{C}$ are prototypes of positively oriented maps, but positively oriented maps, unlike polynomials, do not have to be light or open. Observe that even though in some applications our maps are holomorphic (see Section 7.5), the notion of a positively oriented map is essential in Section 7.5 since it allows for easy local perturbations (see Lemma 7.5.1).

**Definition 2.3.2 (Degree of $f_p$).** Let $f : U \to \mathbb{C}$ be a map from a simply connected domain $U \subseteq \mathbb{C}$ into the plane. Let $S \subseteq \mathbb{C}$ be a positively oriented simple closed curve in $U$, and $p \in U \setminus f^{-1}(f(S))$ a point. Define $f_p : S \to \mathbb{S}^1$ by

$$f_p(x) = \frac{f(x) - f(p)}{|f(x) - f(p)|}.$$ 

Then $f_p$ has a well-defined degree, denoted degree($f_p$). Note that degree($f_p$) is the winding number $\text{wind}(f, S, f(p))$ of $f|_S$ about $f(p)$.

**Definition 2.3.3.** A map $f : U \to \mathbb{C}$ from a simply connected domain $U$ is positively oriented (respectively, negatively oriented) provided for each simple closed curve $S$ in $U$ and each point $p \in T(S) \setminus f^{-1}(f(S))$, we have that degree($f_p$) $> 0$ (degree($f_p$) $< 0$, respectively).

**Definition 2.3.4.** A perfect surjection $f : \mathbb{C} \to \mathbb{C}$ is oriented provided for each simple closed curve $S$ and each $x \in T(S)$, $f(x) \in T(f(S))$.

Clearly every positively oriented and each negatively oriented map is oriented. It will follow that all oriented maps satisfy the Maximum Modulus Theorem 3.7.4 (i.e., for every non-separating continuum $X$, $\partial f(X) \subset f(\partial X)$). In particular, every positively or negatively oriented map is oriented.
It is well-known that both open maps and monotone maps (and hence compositions of such maps) of continua are confluent. It will follow (Lemma 3.7.3) from a result of Lelek and Read [LR74] that each perfect, oriented surjection of the plane is the composition of a monotone map and a light open map.

2.4. Partitioning domains

In Chapter 4 we consider partitions of an open and connected subset $U$ of the sphere into convex subsets which are contained in round balls. Bell originally did this, using the Euclidean metric on the plane, for the complement of $X$ in its convex hull in the plane. Following Kulkarni and Pinkall [KP94], we will consider $U$ as a subset of the sphere and we will work with maximal round balls $B \subset U$ in the spherical metric (such balls correspond to either round balls in the plane or to half planes). We first specify Kulkarni and Pinkall’s result for our situation (see Theorem 4.1.5). It leads to a partition of $U$ into pairwise disjoint closed subsets $F_\alpha$ such for each $\alpha$ there exists a unique maximal closed round ball $B_\alpha$ with $\text{int}(B_\alpha) \cap \partial U = \emptyset$, $|B_\alpha \cap \partial U| \geq 2$ and $F_\alpha \subset B_\alpha$. In fact, $F_\alpha$ is the intersection of $U$ with the hyperbolic convex hull of $B_\alpha \cap \partial U$ in the hyperbolic metric on the ball $B_\alpha$. Note that every chord in the boundary of any partition element $F_\alpha$ is part of a round circle. This is the partition of $U$ which is used in Part 2, Chapter 6.

We show in Section 4.4 that the collection of all chords in the boundaries of all the sets $F_\alpha$, called KP-chords, is sufficiently rich for a satisfactory prime end theory. (Basically most prime ends can be defined through equivalence classes of crosscuts which are all KP-chords.)

However, even though in the above version of the Linchpin Theorem elements of the partition are closed in $U$ and pairwise disjoint and $U$ is an arbitrary connected open subset of the sphere, it has the disadvantage that chords in the boundary of the sets $F_\alpha$ are not naturally depending on $U$ (they depend only on $B_\alpha$, $B_\alpha \cap \partial U$ and the hyperbolic metric in $B_\alpha$). Moreover there may well be uncountably many distinct elements $F_\alpha$ which join the same two accessible points in $\partial U$. In order to avoid this problem we replace, when $U$ is simply connected, any chord in the boundary of any set $F_\alpha$ by the hyperbolic geodesic (in the hyperbolic metric on $U$) joining the same pair of points (see Theorem 4.2.5). We will show that the resulting set of hyperbolic geodesics is a closed lamination of $U$ in the sense of Thurston [Thu09]. This version of the Linchpin Theorem, which states that every point in $U$ is either contained in a unique hyperbolic geodesic $g$ in $U$, or in the interior of an unique hyperbolically convex gap $\mathfrak{G}$, both of which are contained in a maximal round ball, is used in [OT07, OV09] to extend a homeomorphism on the boundary of a simply connected domain over the entire domain. To illustrate the usefulness of these partitions we include a simple proof of the Schoenflies Theorem in Section 4.3. However, we will assume the Schoenflies Theorem throughout this paper.
CHAPTER 3

Tools

3.1. Stability of Index

Let \( f : \mathbb{C} \rightarrow \mathbb{C} \) be a map. All basic definitions of index of \( f \) on a simple closed curve and variation of \( f \) on an arc are contained in Chapter 2. The following standard theorems and observations about the stability of index under a fixed point free homotopy are consequences of the fact that index is continuous and integer-valued.

**Theorem 3.1.1.** Let \( h_t : S^1 \rightarrow \mathbb{C} \) be a homotopy. If \( f : \cup_{t \in [0,1]} h_t(S^1) \rightarrow \mathbb{C} \) is fixed point free, then \( \text{ind}(f, h_0) = \text{ind}(f, h_1) \).

An embedding \( g : S^1 \rightarrow S \subset \mathbb{C} \) is **orientation preserving** if \( g \) is isotopic to the identity map \( \text{id}_{S^1} \). It follows from Theorem 3.1.1 that if \( g_1, g_2 : S^1 \rightarrow S \) are orientation preserving homeomorphisms and \( f : S \rightarrow \mathbb{C} \) is a fixed point free map, then \( \text{ind}(f, g_1) = \text{ind}(f, g_2) \). Hence we can denote \( \text{ind}(f, g_1) \) by \( \text{ind}(f, S) \) and if \([a, b] \) is a positively oriented subarc of \( S^1 \) we denote the fractional index \( \text{ind}(f, g_1|_{[a, b]}) \) by \( \text{ind}(f, g_1([a, b])) \), by some abuse of notation when the extension of \( g_1 \) over \( S^3 \) is understood.

**Theorem 3.1.2.** Suppose \( g : S^1 \rightarrow \mathbb{C} \) is a map with \( g(S^1) = S \), and \( f_1, f_2 : S \rightarrow \mathbb{C} \) are homotopic maps such that each level of the homotopy is fixed point free on \( S \). Then \( \text{ind}(f_1, g) = \text{ind}(f_2, g) \).

In particular, if \( S \) is a simple closed curve and \( f_1, f_2 : S \rightarrow \mathbb{C} \) are maps such that there is a homotopy \( h_t : S \rightarrow \mathbb{C} \) from \( f_1 \) to \( f_2 \) with \( h_t \) fixed point free on \( S \) for each \( t \in [0, 1] \), then \( \text{ind}(f_1, S) = \text{ind}(f_2, S) \).

**Corollary 3.1.3.** Suppose \( g : S^1 \rightarrow \mathbb{C} \) is an orientation preserving embedding with \( g(S^1) = S \), and \( f : S \rightarrow T(S) \) is a fixed point free map. Then \( \text{ind}(f, g) = \text{ind}(f, S) = 1 \).

**Proof.** Since \( f(S) \subset T(S) \) which is a disk with boundary \( S \) and \( f \) has no fixed point on \( S \), there is a fixed point free homotopy of \( f|_S \) to a constant map \( c : S \rightarrow \mathbb{C} \) taking \( S \) to a point in \( T(S) \setminus S \). By Theorem 3.1.2, \( \text{ind}(f, g) = \text{ind}(c, g) \). Since \( g \) is orientation preserving it follows from Remark 2.1.1 (b) that \( \text{ind}(c, g) = 1 \).

**Theorem 3.1.4.** Suppose \( g : S^1 \rightarrow \mathbb{C} \) is a map with \( g(S^1) = S \), and \( f : T(S) \rightarrow \mathbb{C} \) is a map such that \( \text{ind}(f, g) \neq 0 \), then \( f \) has a fixed point in \( T(S) \).

**Proof.** Notice that \( T(S) \) is a locally connected, non-separating, plane continuum and, hence, contractible. Suppose \( f \) has no fixed point in \( T(S) \). Choose point \( q \in T(S) \). Let \( c : S^1 \rightarrow \mathbb{C} \) be the constant map \( c(S^1) = \{q\} \). Let \( H \) be a homotopy from \( g \) to \( c \) with image in \( T(S) \). Since \( H \) misses the fixed point set of \( f \), Theorem 3.1.1 and Remark 2.1.1 (a) imply \( \text{ind}(f, g) = \text{ind}(f, c) = 0 \).
3.2. Index and variation for finite partitions

What links Theorem 3.1.4 with variation is Theorem 3.2.2 below, first announced by Bell in the mid 1980’s (see also Akis [Aki99]). Our proof is a modification of Bell’s unpublished proof. We first need a variant of Proposition 2.2.3. Let \( r : \mathbb{C} \to T(S^1) \) be radial retraction: \( r(z) = \frac{z}{|z|} \) when \(|z| \geq 1\) and \( r|_{T(S^1)} = id|_{T(S^1)} \).

**Lemma 3.2.1 (Curvature Straightening).** Suppose \( f : S^1 \to \mathbb{C} \) is a map with no fixed points on \( S^1 \). If \([a,b]\subset S^1\) is a proper subarc with \( f([a,b])\cap[a,b] = \emptyset\), \( f((a,b)) \subset \mathbb{C}\setminus T(S^1)\) and \( f((a,b)) \subset S^1\), then there exists a map \( \tilde{f} : S^1 \to \mathbb{C} \) such that \( \tilde{f}|_{S^1 \setminus (a,b)} = f|_{S^1 \setminus (a,b)} \), \( \tilde{f}|_{[a,b]} : [a,b] \to (\mathbb{C}\setminus T(S^1)) \cup \{f(a), f(b)\} \) and \( \tilde{f}|_{[a,b]} \) is homotopic to \( f|_{[a,b]} \) in \( (a,b) \cup \mathbb{C}\setminus T(S) \) relative to \([a,b]\), so that either \( r|_{\tilde{f}|_{[a,b]}} \) is locally one-to-one or a constant map. Moreover, \( \text{var}(f,[a,b],S^1) = \text{var}(\tilde{f},[a,b],S^1) \).

Note that if \( \text{var}(f,[a,b],S^1) = 0 \), then \( r \) carries \( \tilde{f}|_{[a,b]} \) one-to-one onto the arc (or point) in \( S^1 \setminus (a,b) \) from \( f(a) \) to \( f(b) \). If the \( \text{var}(f,[a,b],S^1) = m > 0 \), then \( r \circ \tilde{f} \) wraps the arc \([a,b]\) counterclockwise about \( S^1 \) so that \( \tilde{f}|_{[a,b]} \) meets each ray in \( J_\nu \) \( m \) times. A similar statement holds for negative variation. Note also that it is possible for index to be defined yet variation not to be defined on a simple closed curve \( S \). For example, consider the map \( z \to 2z \) with \( S \) the unit circle since there is no partition of \( S \) satisfying the conditions in Definition 2.2.2.

**Theorem 3.2.2 (Index = Variation + 1, Bell).** Suppose \( g : S^1 \to \mathbb{C} \) is an orientation preserving embedding onto a simple closed curve \( S \) and \( f : S \to \mathbb{C} \) is a fixed point free map. If \( F = \{a_0 < a_1 < \cdots < a_n\} \) is a partition of \( S \) and \( A_i = [a_i, a_{i+1}] \) for \( i = 0, 1, \ldots, n \) with \( a_{n+1} = a_0 \) such that \( f(F) \subset T(S) \) and \( f(A_i) \cap A_i = \emptyset \) for each \( i \), then

\[
\text{ind}(f,S) = \text{ind}(f,g) = \sum_{i=0}^{n} \text{var}(f, A_i, S) + 1 = \text{var}(f,S) + 1.
\]

**Proof.** By an appropriate conjugation of \( f \) and \( g \), we may assume without loss of generality that \( S = S^1 \) and \( g = id \). Let \( F \) and \( A_i = [a_i, a_{i+1}] \) be as in the hypothesis. Consider the collection of arcs

\[ \mathcal{K} = \{ K \subset S \mid K \text{ is the closure of a component of } f^{-1}(f(S) \setminus T(S)) \}. \]

For each \( K \in \mathcal{K} \), there is an \( i \) such that \( K \subset A_i \). Since \( f(A_i) \cap A_i = \emptyset \), it follows from the remark after Definition 2.2.2 that \( \text{var}(f, A_i, S) = \sum_{K \subset A_i, K \in \mathcal{K}} \text{var}(f, K, S) \). By the remark following Proposition 2.2.3, we can compute \( \text{var}(f, K, S) \) using one fixed junction for \( A_i \). It is now clear that there are at most finitely many \( K \in \mathcal{K} \) with \( \text{var}(f, K, S) \neq 0 \). Moreover, the images of the endpoints of each \( K \) lie on \( S \).

Let \( m \) be the cardinality of the set \( \mathcal{K}_F = \{ K \in \mathcal{K} \mid \text{var}(f, K, S) \neq 0 \} \). By the above remarks, \( m < \infty \) and \( \mathcal{K}_F \) is independent of the partition \( F \). We prove the theorem by induction on \( m \).

Suppose for a given \( f \) we have \( m = 0 \). Observe that from the definition of variation and the fact that the computation of variation is independent of the choice of an appropriate partition, it follows that,

\[ \text{var}(f, S) = \sum_{K \in \mathcal{K}} \text{var}(f, K, S) = 0. \]
We claim that there is a map \( f_1 : S \to \mathbb{C} \) with \( f_1(S) \subset T(S) \) and a homotopy \( H \) from \( f|_S \) to \( f_1 \) such that each level \( H_t \) of the homotopy is fixed point free and \( \text{ind}(f_1, id|_S) = 1 \).

To see the claim, first apply the Curve Straightening Lemma 3.2.1 to each \( K \in \mathcal{K} \) (if there are infinitely many, they form a null sequence) to obtain a fixed point free homotopy of \( f|_S \) to a map \( \tilde{f} : S \to \mathbb{C} \) such that \( \partial \tilde{f(K)} \) is locally one-to-one (or the constant map) on each \( f \) and \( \text{var}(\tilde{f}) = 0 \) for each \( K \in \mathcal{K} \). Let \( K \) be in \( \mathcal{K} \) with endpoints \( x, y \). Since \( \tilde{f}(K) \cap K = \emptyset \) and \( \text{var}(\tilde{f}, K, S) = 0 \) for each \( K \in \mathcal{K} \), we obtain the desired \( f_1 : S \to \mathbb{C} \) as the end map of a fixed point free homotopy from \( f \) to \( f_1 \). Since \( f_1 \) carries \( S \) into \( T(S) \), Corollary 3.1.3 implies \( \text{ind}(f_1, id|_S) = 1 \).

Since the homotopy \( f \simeq f_1 \) is fixed point free, it follows from Theorem 3.1.2 that \( \text{ind}(f, id|_S) = 1 \). Hence, the theorem holds if \( m = 0 \) for any \( f \) and any appropriate partition \( F \).

By way of contradiction suppose the collection \( \mathcal{F} \) of all maps \( f \) on \( S^1 \) which satisfy the hypotheses of the theorem, but not the conclusion is non-empty. By the above \( 0 < |K_f| < \infty \) for each. Let \( f \in \mathcal{F} \) be a counterexample for which \( m = |K_f| \) is minimal. By modifying \( f \), we will show there exists \( f_1 \in \mathcal{F} \) with \( |K_f| < m \), a contradiction.

Choose \( K \in \mathcal{K} \) such that \( \text{var}(f, K, S) \neq 0 \). Then \( K = [x, y] \subset A_i = [a_i, a_{i+1}] \) for some \( i \). By the Curve Straightening Lemma 3.2.1 and Theorem 3.1.2, we may suppose \( \partial f(K) \) is locally one-to-one on \( K \). Define a new map \( f_1 : S \to \mathbb{C} \) by setting \( f_1|_{S \setminus K} = f|_{S \setminus K} \) and setting \( f_1|_K = f|_K \). Then \( f_1 \) restricted to \( K \) is equal to the linear map taking \( [x, y] \) to \( f(x) \) to \( f(y) \) on \( S \) missing \( [x, y] \). Figure 3.1 (left) shows an example of a (straightened) \( f \) restricted to \( K \) and the corresponding \( f_1 \) restricted to \( K \) for a case where \( \text{var}(f, K, S) = 1 \), while Figure 3.1 (right) shows a case where \( \text{var}(f, K, S) = -2 \).

Since on \( S \setminus K \), \( f \) and \( f_1 \) are the same map, we have

\[
\text{var}(f, S \setminus K, S) = \text{var}(f_1, S \setminus K, S).
\]

Likewise for the fractional index,

\[
\text{ind}(f, S \setminus K) = \text{ind}(f_1, S \setminus K).
\]

By definition (refer to the observation we made in the case \( m = 0 \)),

\[
\text{var}(f, S) = \text{var}(f, S \setminus K, S) + \text{var}(f, K, S)
\]

\[
\text{var}(f_1, S) = \text{var}(f_1, S \setminus K, S) + \text{var}(f_1, K, S)
\]

and by Proposition 2.1.2,

\[
\text{ind}(f, S) = \text{ind}(f, S \setminus K) + \text{ind}(f, K)
\]

\[
\text{ind}(f_1, S) = \text{ind}(f_1, S \setminus K) + \text{ind}(f_1, K).
\]

Consequently,

\[
\text{var}(f, S) - \text{var}(f_1, S) = \text{var}(f, K, S) - \text{var}(f_1, K, S)
\]
Figure 3.1. Replacing \( f : S \to \mathbb{C} \) by \( f_1 : S \to \mathbb{C} \) with one less subarc of nonzero variation.

\[
\text{and} \quad \text{ind}(f, S) - \text{ind}(f_1, S) = \text{ind}(f, K) - \text{ind}(f_1, K).
\]

We will now show that the changes in index and variation, going from \( f \) to \( f_1 \) are the same (i.e., we will show that \( \text{var}(f, K, S) - \text{var}(f_1, K, S) = \text{ind}(f, K) - \text{ind}(f_1, K) \)). We suppose first that \( \text{ind}(f, K) = n + \alpha \) for some nonnegative \( n \in \mathbb{N} \) and \( 0 \leq \alpha < 1 \). That is, the vector \( f(z) - z \) turns through \( n \) full revolutions counterclockwise and \( \alpha \) part of a revolution counterclockwise as \( z \) goes from \( x \) to \( y \) counterclockwise along \( S \). (See Figure 3.1 (left) for the case \( n = 0 \) and \( \alpha \) about \( 0 \).)

Assume first that \( f(x) < x < y < f(y) \) in the circular order as illustrated in Figure 3.1 on the left. Then as \( z \) goes from \( x \) to \( y \) counterclockwise along \( S \), \( f_1(z) \) goes along \( S \) from \( f(x) \) to \( f(y) \) in the clockwise direction, so \( f_1(z) - z \) turns through \(-(1 - \alpha) = \alpha - 1 \) part of a revolution. Hence, \( \text{ind}(f_1, K) = \alpha - 1 \). It is easy to see that \( \text{var}(f, K, S) = n + 1 \) and \( \text{var}(f_1, K, S) = 0 \). Consequently,

\[
\text{var}(f, K, S) - \text{var}(f_1, K, S) = n + 1 - 0 = n + 1
\]

and

\[
\text{ind}(f, K) - \text{ind}(f_1, K) = n + \alpha - (\alpha - 1) = n + 1.
\]

We assumed that \( f(x) < x < y < f(y) \). The cases where \( f(y) < y < f(x) \) and \( f(x) = f(y) \) (and, hence, \( \alpha = 0 \)) are treated similarly. In this case \( f_1 \) still wraps around in the positive direction, but the computations are slightly different: \( \text{var}(f, K) = n, \text{ind}(f, K) = n + \alpha, \text{var}(f_1, K) = 0 \) and \( \text{ind}(f_1, K) = \alpha \).

Thus when \( n \geq 0 \), in going from \( f \) to \( f_1 \), the change in variation and the change in index are the same. However, in obtaining \( f_1 \) we have removed one \( K \in \mathcal{K}_f \), reducing the minimal \( m = |\mathcal{K}_f| \) for \( f \) by one, producing a counterexample \( f_1 \) with \( |\mathcal{K}_{f_1}| = m - 1 \), a contradiction.
The cases where \( \text{ind}(f,K) = n + \alpha \) for negative \( n \) and \( 0 < \alpha < 1 \) are handled similarly, and illustrated for \( n = -2, \alpha \) about 0.4 and \( f(y) < x < y < f(x) \) in Figure 3.1 (right). \( \square \)

### 3.3. Locating arcs of negative variation

The principal tool in proving Theorem 6.2.1 (unique outchannel) is the following theorem first obtained by Bell (unpublished). It provides a method for locating arcs of negative variation on a curve of index zero.

**Theorem 3.3.1** (Lollipop Lemma, Bell). Let \( S \subset \mathbb{C} \) be a simple closed curve and \( f : T(S) \to \mathbb{C} \) a fixed point free map. Suppose \( F = \{a_0 < \cdots < a_n < a_{n+1} < \cdots < a_m\} \) is a partition of \( S \), \( a_{n+1} = a_0 \) and \( A_i = [a_i, a_{i+1}] \) such that \( f(F) \subset T(S) \) and \( f(A_i) \cap A_j = \emptyset \) for \( i = 0, \ldots, m \). Suppose \( I \) is an arc in \( T(S) \) meeting \( S \) only at its endpoints \( a_0 \) and \( a_{n+1} \). Let \( J_{a_0} \) be a junction in \( (C \setminus T(S)) \cup \{a_0\} \) and suppose that \( f(I) \cap (I \cup J_{a_0}) = \emptyset \). Let \( R = T([a_0, a_{n+1}] \cup I) \) and \( L = T(\{a_{n+1}, a_{m+1}\} \cup I) \).

Then one of the following holds:

1. If \( f(a_{n+1}) \in R \), then
   \[
   \sum_{i\leq n} \text{var}(f, A_i, S) + 1 = \text{ind}(f, I \cup [a_0, a_{n+1}]).
   \]
2. If \( f(a_{n+1}) \in L \), then
   \[
   \sum_{i> n} \text{var}(f, A_i, S) + 1 = \text{ind}(f, I \cup [a_{n+1}, a_{m+1}]).
   \]

(Note that in (1) in effect we compute \( \text{var}(f, \partial R) \) but technically, we have not defined \( \text{var}(f, A_i, \partial R) \) since the endpoints of \( A_i \) do not have to map inside \( R \) but they do map into \( T(S) \). Similarly in Case (2).)

**Proof.** Suppose \( f(a_{n+1}) \in L \) (the case when \( f(a_{n+1}) \in R \) can be treated similarly). Consider the set \( C = [a_{n+1}, a_{m+1}] \cup I \) (so \( T(C) = L \)). We want to construct a map \( f' : C \to \mathbb{C} \), fixed point free homotopic to \( f|_C \), that does not change variation on any arc \( A_i \) in \( C \) and has the properties listed below.

1. \( f'(a_i) \in L, f'(A_i) \cap A_j = \emptyset \) for all \( n + 1 \leq i \leq m \) and \( f'(a_0) \in L \). Hence \( \text{var}(f', A_i, C) \) is defined for each \( i > n \).
2. \( \text{var}(f', A_i, C) = \text{var}(f, A_i, S) \) for all \( n + 1 \leq i \leq m \).
3. \( f'(I) \cap I = \emptyset \) and \( \text{var}(f', I, C) = 0 \).

Having such a map, it then follows from Theorem 3.2.2, that

\[
\text{ind}(f', C) = \sum_{i=n+1}^{m} \text{var}(f', A_i, C) + \text{var}(f', I, C) + 1.
\]

By Theorem 3.1.2 \( \text{ind}(f', C) = \text{ind}(f, C) \). By (2) and (3), \( \sum_{i> n} \text{var}(f', A_i, C) + \text{var}(f', I, C) = \sum_{i> n} \text{var}(f, A_i, S) \) and the Theorem would follow.

It remains to define the map \( f' : C \to \mathbb{C} \) with the above properties. For each \( i \) such that \( n + 1 \leq i \leq m + 1 \), choose an arc \( I_i \) joining \( f(a_i) \) to \( L \) as follows:

(a) If \( f(a_i) \in L \), let \( I_i \) be the degenerate arc \( \{f(a_i)\} \).
(b) If \( f(a_i) \in R \) and \( n + 1 < i < m + 1 \), let \( I_i \) be an arc in \( R \setminus \{a_0, a_{n+1}\} \) joining \( f(a_i) \) to \( I \).
(c) If \( f(a_0) \in R \), let \( I_0 \) be an arc joining \( f(a_0) \) to \( L \) such that \( I_0 \cap (L \cup J_{a_0}) \subset A_{n+1} \setminus \{a_{n+1}\} \).
Let $x_{n+1} = y_{n+1} = a_{n+1}$, $y_0 = y_{m+1} \in I \setminus \{a_0, a_{m+1}\}$ and $x_0 = x_{m+1} \in A_m \setminus \{a_m, a_{m+1}\}$. For $n+1 < i < m+1$, let $x_i \in A_{i-1}$ and $y_i \in A_i$ such that $y_{i-1} < x_i < a_i < y_i < x_{i+1}$. For $n+1 < i < m+1$ let $f'(a_i)$ be the endpoint of $I_i$ in $L$, $f'(x_i) = f'(y_i) = f(a_i)$ and extend $f'$ continuously from $[x_i, a_i] \cup [a_i, y_i]$ onto $I_i$ and define $f'$ from $[y_i, x_{i+1}] \subset A_i$ onto $f(A_i)$ by $f'|_{[y_i, x_{i+1}]} = f \circ h_i$, where $h_i : [y_i, x_{i+1}] \rightarrow A_i$ is a homeomorphism such that $h_i(y_i) = a_i$ and $h_i(x_{i+1}) = a_{i+1}$. Similarly, define $f'$ on $[y_0, a_{n+1}] \subset I$ to $f(I)$ by $f'|_{[y_0, a_{n+1}]} = f \circ h_0$, where $h_0 : [y_0, a_{n+1}] \rightarrow I$ is an onto homeomorphism such that $h(a_{n+1}) = a_{n+1}$ and extend $f'$ from $[x_{m+1}, a_0] \subset A_m$ and $[a_0, y_0] \subset I$ onto $I_0$ such that $f'(x_{m+1}) = f'(y_0) = f(a_0)$ and $f'(a_0)$ is the endpoint of $I_0$ in $L$. To define $f'|_{[a_{n+1}, x_{n+2}]}$ let $h_{n+1} : [y_{n+1}, x_{n+2}] \rightarrow [a_{n+1}, a_{n+2}]$ be a homeomorphism such that $h_{n+1}(y_{n+1}) = a_{n+1}$. Then define $f'(x)$ as $f \circ h_{n+1}(x)$ for $x \in [y_{n+1}, x_{n+2}]$ and $f'(x) = f(a_{n+1})$ if $x \in [a_{n+1}, y_{n+1}]$.
Note that \( f'(A_i) \cap A_i = \emptyset \) for \( i = n + 1, \ldots, m \) and \( f'(I) \cap [I \cup J_m] = \emptyset \). To compute the variation of \( f' \) on each of \( A_m \) and \( I \) we can use the junction \( J_m \). Hence \( \text{var}(f', I, C) = 0 \) and, by the definition of \( f' \) on \( A_m \), \( \text{var}(f', A_m, C) = \text{var}(f, A_m, S) \). For \( i = n + 1, \ldots, m - 1 \) we can use the same junction \( J_i \) to compute \( \text{var}(f', A_i, C) \) as we did to compute \( \text{var}(f, A_i, S) \). Since \( I_i \cup I_{i+1} \subset T(S) \setminus A_i \), we have that \( f'([a_i, y_i]) \cup f'([x_{i+1}, a_{i+1}]) \subset I_i \cup I_{i+1} \) misses that junction and, hence, make no contribution to variation \( \text{var}(f', A_i, C) \). Since \( f'^{-1}(J_n) \cap [y_i, x_{i+1}] \) is isomorphic to \( f'^{-1}(J_{n_i}) \cap A_i \), \( \text{var}(f', A_i, C) = \text{var}(f, A_i, S) \) for \( i = n + 1, \ldots, m \).

To see that \( f' \) is fixed point free homotopic to \( f|_{C_i} \), note that we can pull the image of \( A_i \) back along the arcs \( I_i \) and \( I_{i+1} \) in \( R \) without fixing a point of \( A_i \) at any level of the homotopy. \( \square \)

Note that if \( f \) is fixed point free on \( T(S) \), then \( \text{ind}(f, C) = 0 \) and the next Corollary follows.

**Corollary 3.3.2.** Assume the hypotheses of Theorem 3.3.1. Then if \( f(a_{n+1}) \in R \) there exists \( i \leq n \) such that \( \text{var}(f, A_i, S) < 0 \). If \( f(a_{n+1}) \in L \) there exists \( i > n \) such that \( \text{var}(f, A_i, S) < 0 \).

### 3.4. CROSSCUTS AND BUMPING ARCS

For the remainder of Chapter 3, we assume that \( f : C \to C \) takes the continuum \( X \) into \( T(X) \) with no fixed points in \( T(X) \), and \( X \) is minimal with respect to these properties.

**Definition 3.4.1 (Bumping Simple Closed Curve).** A simple closed curve \( S \) in \( C \) which has the property that \( S \cap X \) is nondegenerate and \( T(X) \subset T(S) \) is said to be a **bumping simple closed curve for \( X \)**. A subarc \( A \) of a bumping simple closed curve, whose endpoints lie in \( X \), is said to be a **bumping (sub)arc for \( X \)** or a **link of \( S \)**. Moreover, if \( S' \) is any bumping simple closed curve for \( X \) which contains \( A \), then \( S' \) is said to **complete \( A \)**. In fact, an arc \( A \) with endpoints in \( X \) which can be completed will be called a **bumping arc of \( X \)**.

Given a positively oriented simple closed curve \( S \), we can consider its positively oriented subarcs denoted by \([a, b]_S\), where \( a, b \) are the endpoints of the arc; if the curve is fixed, we simply write \([a, b]\). Similar notation is used for half-open or open subarcs of bumping simple closed curves. Often we will fix the choice of links into which we divide \( S \). In general a bumping arc of \( X \) may have points other than its endpoints which belong to \( X \) (e.g., if \( X \) is the closed unit disk and the unit circle is divided into several subarcs then each of them can be considered as a bumping arc of \( X \)). By definition, any bumping arc \( A \) of \( X \) can be extended to a bumping simple closed curve \( S \) of \( X \). Hence, every bumping arc has a well defined natural order \(<\) inherited from the positive circular order of a bumping simple closed curve \( S \) containing \( A \). If \( a < b \) are the endpoints of \( A \), then we will often write \( A = [a, b] \).

A **crosscut** of \( U^\infty = C^\infty \setminus T(X) \) is an open arc \( Q \) lying in \( U^\infty \setminus \{\infty\} \) such that \( Q \) is an arc with endpoints \( a \neq b \in T(X) \). In this case we will often write \( Q = (a, b) \). (As seems to be traditional, we use “crosscut of \( T(X) \)” interchangeably with “crosscut of \( U^\infty \).”) Evidently, a crosscut of \( U^\infty \) separates \( U^\infty \) into two disjoint domains, exactly one of which is unbounded. If \( S \) is a bumping simple closed curve so that \( X \cap S \) is nondegenerate, then each component of \( S \setminus X \) is a crosscut of \( T(X) \). A similar statement holds for a bumping arc \( A \). Given a non-separating
continuum $T(X)$, let $A \subset \mathbb{C}$ be a crosscut of $U^\infty(X) = \mathbb{C}^\infty \setminus T(X)$. Given a crosscut $A$ of $U^\infty(X)$ denote by $\text{Sh}(A)$, the shadow of $A$, the bounded component of $\mathbb{C} \setminus [T(X) \cup A]$. Moreover, suppose that $A$ is a bumping arc of $X$. Then by the shadow $\text{Sh}(A)$ of $A$, we mean the union of all bounded components of $\mathbb{C} \setminus (X \cup A)$ (since there may be more than endpoints of $A$ in $X \cap A$, we should talk about the union of all bounded components of $\mathbb{C} \setminus (X \cup A)$ here).

A variety of tools (such as index, variation, junction) have been described in previous sections. So far they have been applied to the properties of maps of the plane restricted to simple closed curves. Another application can be found in Theorem 3.1.4 where the existence of a fixed point in the topological hull of a simple closed curve is established. However we are mostly interested in studying continua as described in our Standing Hypotheses. The following construction shows how the above described tools apply in this situation.

Since $f$ has no fixed points in $T(X)$ and $X$ is compact, we can choose a bumping simple closed curve $S$ in a small neighborhood of $T(X)$ such that all crosscuts in $S \setminus X$ are small, have positive distance to their image and so that $f$ has no fixed points in $T(S)$. Thus, we obtain the following corollary to Theorem 3.1.4.

**Corollary 3.4.2.** Let $f : \mathbb{C} \to \mathbb{C}$ be a map and $X \subset \mathbb{C}$ a subcontinuum with $f(X) \subset T(X)$ and so that $f|_{T(X)}$ is fixed point free. Then there is a bumping simple closed curve $S$ for $X$ such that $f|_{T(S)}$ is fixed point free; hence, by 3.1.4, $\text{ind}(f,S) = 0$.

Moreover, any bumping simple closed curve $S'$ for $X$ such that $S' \subset T(S)$ has $\text{ind}(f,S') = 0$. Furthermore, any bumping arc $A$ of $T(X)$ for which $f$ has no fixed points in $T(X \cup A)$ can be completed to a bumping simple closed curve $S$ for $X$ for which $\text{ind}(f,S) = 0$.

The idea of the proofs of a few forthcoming results is that in some cases we can use the developed tools (e.g., variation) in order to compute out index and show, relying upon the properties of our maps, that index is not equal to zero thus contradicting Corollary 3.4.2. To implement such a plan we need to further study properties of variation in the setting described before Corollary 3.4.2.

**Proposition 3.4.3.** Let $f : \mathbb{C} \to \mathbb{C}$ be a map and $X \subset \mathbb{C}$ a subcontinuum with $f(X) \subset T(X)$ and so that $f|_{T(X)}$ is fixed point free. In the situation of Corollary 3.4.2, suppose $A$ is a bumping subarc for $X$. If $\text{var}(f,A,S)$ is defined for some bumping simple closed curve $S$ completing $A$, then for any bumping simple closed curve $S'$ completing $A$, $\text{var}(f,A,S) = \text{var}(f,A,S')$.

**Proof.** Since $\text{var}(f,A,S)$ is defined, $A = \cup_{i=1}^n A_i$, where each $A_i$ is a bumping arc with $A_i \cap f(A_i) = \emptyset$ and $|A_i \cap A_j| \leq 1$ if $i \neq j$. By the remark following Definition 2.2.5, it suffices to establish the desired result for each $A_i = A$. Let $S$ and $S'$ be two bumping simple closed curves completing $A$ for which variation is defined. Let $J_a$ and $J_{a'}$ be junctions whereby $\text{var}(f,A,S)$ and $\text{var}(f,A,S')$ are respectively computed. Suppose first that both junctions lie (except for $\{a,a'\}$) in $\mathbb{C} \setminus (T(S) \cup T(S'))$. By the Junction Straightening Proposition 2.2.3, either junction can be used to compute either variation on $A$, so the result follows. Otherwise, at least one junction is not in $\mathbb{C} \setminus (T(S) \cup T(S'))$. But both junctions are in $\mathbb{C} \setminus T(X \cup A)$. Hence, we can find another bumping simple closed curve $S''$ such that $S''$ completes $A$, and both junctions lie in $(\mathbb{C} \setminus T(S'')) \cup \{a,a'\}$. Then by the Propositions 2.2.3 and the definition of variation, $\text{var}(f,A,S) = \text{var}(f,A,S'') = \text{var}(f,A,S')$. □
It follows from Proposition 3.4.3 that variation on a crosscut $Q$, with $Q \cap f(Q) = \emptyset$, of $T(X)$ is independent of the bumping simple closed curve $S$ for $T(X)$ of which $Q$ is a subarc and is such that $\text{var}(f, Q, S)$ is defined. Hence, given a bumping arc $A$ of $X$, we can denote $\text{var}(f, A, S)$ by $\text{var}(f, A, X)$ or simply by $\text{var}(f, A)$ when $X$ is understood. The figure illustrates how variation is computed.

The following proposition follows from Corollary 3.4.2, Proposition 3.4.3 and Theorem 3.2.2.

**Proposition 3.4.4.** Let $f : \mathbb{C} \to \mathbb{C}$ be a map and $X$ a subcontinuum of $\mathbb{C}$ so that $f(X) \subset T(X)$ and $f$ has no fixed points in $T(X)$. Suppose $Q$ is a crosscut of $T(X)$ such that $f$ is fixed point free on $T(X \cup Q)$ and $f(Q) \cap Q = \emptyset$. Suppose $Q$ is replaced by a bumping subarc $A$ with the same endpoints such that $Q \cup T(X)$ separates $A \setminus X$ from $\infty$ and each component $Q_i$ of $A \setminus X$ is a crosscut such that $f(Q_i) \cap Q_i = \emptyset$. Then

$$\text{var}(f, Q, X) = \sum_i \text{var}(f, Q_i, X) = \text{var}(f, A, X).$$

### 3.5. Index and Variation for Carathéodory Loops

We extend the definitions of index and variation to *Carathéodory loops*. In particular, if $g : \mathbb{S}^1 \to g(\mathbb{S}^1) = S$ is a continuous extension of a Riemann map $\psi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus T(g(S^1))$, then $g$ is a Carathéodory loop, where $\mathbb{D}^\infty = \{ z \in \mathbb{C}^\infty \mid |z| > 1 \}$ is the “unit disk” about $\infty$.

**Definition 3.5.1 (Carathéodory Loop).** Let $g : \mathbb{S}^1 \to \mathbb{C}$ such that $g$ is continuous and has a continuous extension $\psi : \mathbb{C}^\infty \setminus T(\mathbb{S}^1) \to \mathbb{C}^\infty \setminus T(g(\mathbb{S}^1))$ such
that \( \psi|_{\subset T(g)} \) is an orientation preserving homeomorphism from \( \subset \setminus T(S^1) \) onto \( \subset \setminus T(g(S^1)) \). We call \( g \) (and loosely, \( S = g(S^1) \)), a Carathéodory loop.

Let \( g : S^1 \to \subset \) be a Carathéodory loop and let \( f : g(S^1) \to \subset \) be a fixed point free map. In order to define variation of \( f \) on \( g(S^1) \), we do the partitioning in \( S^1 \) and transport it to the Carathéodory loop \( S = g(S^1) \). An allowable partition of \( S^1 \) is a set \( \{a_0 < a_1 < \cdots < a_n\} \) in \( S^1 \) ordered counterclockwise, where \( a_0 = a_n \) and \( A_i \) denotes the counterclockwise interval \([a_i, a_{i+1}]\), such that for each \( i \), \( f(g(a_i)) \in T(g(S^1)) \) and \( f(g(A_i)) \cap g(A_i) = \emptyset \). Variation \( \text{var}(f, A_i, g(S^1)) = \text{var}(f, A_i) \) on each path \( g(A_i) \) is then defined exactly as in Definition 2.2.2, except that the junction (see Definition 2.2.1) is chosen so that the vertex \( v \in g(A_i) \) and \( J_i \cap T(g(S^1)) \subset \{v\} \), and the crossings of the junction \( J_v \) by \( f(g(A_i)) \) are counted (see Definition 2.2.2).

Variation on the whole loop, or an allowable subarc thereof, is defined just as in Definition 2.2.5, by adding the variations on the partition elements. At this point in the development, variation is defined only relative to the given allowable partition \( F \) of \( S^1 \) and the parameterization \( g \) of \( S \) : \( \text{var}(f, F, g(S^1)) \).

Index on a Carathéodory loop \( S \) is defined exactly as in Section 2.1 with \( S = g(S^1) \) providing the parameterization of \( S \). Likewise, the definition of fractional index and Proposition 2.1.2 apply to Carathéodory loops.

Theorems 3.1.1, 3.1.2, Corollary 3.1.3, and Theorem 3.1.4 (if \( f \) is also defined on \( T(S) \)) apply to Carathéodory loops. It follows that index on a Carathéodory loop \( S \) is independent of the choice of parameterization \( g \). The Carathéodory loop \( S \) is approximated, under small homotopies, by simple closed curves \( S_i \). Allowable partitions of \( S \) can be made to correspond to allowable partitions of \( S_i \) under small homotopies. Since variation and index are invariant under suitable homotopies (see the comments after Proposition 2.2.3) we have the following theorem.

**Theorem 3.5.2.** Suppose \( S = g(S^1) \) is a parameterized Carathéodory loop in \( \subset \) and \( f : S \to \subset \) is a fixed point free map. Suppose variation of \( f \) on \( S^1 = A_0 \cup \cdots \cup A_n \) with respect to \( g \) is defined for some partition \( A_0 \cup \cdots \cup A_n \) of \( S^1 \). Then \[
\text{ind}(f, g) = \sum_{i=0}^{n} \text{var}(f, A_i, g(S^1)) + 1.
\]

### 3.6. Prime Ends

Prime ends provide a way of studying the approaches to the boundary of a simply-connected plane domain with non-degenerate boundary. See [CL66] or [MiI00] for an analytic summary of the topic and [UY51] for a more topological approach. We will be interested in the prime ends of \( U^\infty = C^\infty \setminus T(X) \). Recall that \( D^\infty = \{z \in C^\infty \mid |z| > 1\} \) is the “unit disk about \( \infty \).” The Riemann Mapping Theorem guarantees the existence of a conformal map \( \phi : D^\infty \to U^\infty \) taking \( \infty \to \infty \), unique up to the argument of the derivative at \( \infty \). Fix such a map \( \phi \). We identify \( S^1 = \partial D^\infty \) with \( \mathbb{R}/\mathbb{Z} \) and identify points \( e^{2\pi it} \) in \( \partial D^\infty \) by their argument \( t \) (mod 1). Crosscut and shadow were defined in Section 3.4.

**Definition 3.6.1 (Prime End).** A chain of crosscuts is a sequence \( \{Q_i\}_{i=1}^\infty \) of crosscuts of \( U^\infty \) such that for \( i \neq j \), \( \overline{Q_i} \cap \overline{Q_j} = \emptyset \), \( \text{diam}(Q_i) \to 0 \), and for all \( j > i \), \( Q_i \) separates \( Q_j \) from \( \infty \) in \( U^\infty \). Hence, for all \( j > i \), \( Q_j \subset \text{Sh}(Q_i) \). Two chains of crosscuts are said to be equivalent if and only if it is possible to form a sequence of crosscuts by selecting alternately a crosscut from each chain so that the resulting
sequence of crosscuts is again a chain. A prime end $\mathcal{E}$ is an equivalence class of chains of crosscuts.

If $\{Q_i\}$ and $\{Q'_i\}$ are equivalent chains of crosscuts of $U^\infty$, it can be shown that $\{\phi^{-1}(Q_i)\}$ and $\{\phi^{-1}(Q'_i)\}$ are equivalent chains of crosscuts of $\mathbb{D}^\infty$ each of which converges to the same unique point $e^{2\pi it} \in S^1 = \partial \mathbb{D}^\infty$, $t \in [0, 1)$, independent of the representative chain. Hence, we denote by $\mathcal{E}_t$ the prime end of $U^\infty$ defined by $\{Q_i\}$.

**Definition 3.6.2 (Impression and Principal Continuum).** Let $\mathcal{E}_t$ be a prime end of $U^\infty$ with defining chain of crosscuts $\{Q_i\}$. The set

$$\text{Im}(\mathcal{E}_t) = \bigcap_{i=1}^\infty \text{Sh}(Q_i)$$

is a subcontinuum of $\partial U^\infty$ called the impression of $\mathcal{E}_t$. The set

$$\text{Pr}(\mathcal{E}_t) = \{z \in \partial U^\infty \mid \text{for some chain } \{Q'_i\} \text{ defining } \mathcal{E}_t, Q'_i \to z \}$$

is a continuum called the principal continuum of $\mathcal{E}_t$.

For a prime end $\mathcal{E}_t$, $\text{Pr}(\mathcal{E}_t) \subset \text{Im}(\mathcal{E}_t)$, possibly properly. We will be interested in the existence of prime ends $\mathcal{E}_t$ for which $\text{Pr}(\mathcal{E}_t) = \text{Im}(\mathcal{E}_t) = \partial U^\infty$.

**Definition 3.6.3 (External Rays).** Let $t \in [0, 1)$ and define

$$R_t = \{z \in \mathbb{C} \mid z = \phi(r e^{2\pi it}), 1 < r < \infty\}.$$

We call $R_t$ the external ray (with argument $t$). If $x \in R_t$ then the $(X, x)$-end of $R_t$ is the bounded component $K_x$ of $R_t \setminus \{x\}$.

In this case $X$ is a continuum, $U^\infty(X)$ is simply connected, the external rays $R_t$ are all smooth and pairwise disjoint. Moreover, for each $x \in U^\infty(X)$ there exists a unique $t$ such that $x \in R_t$.

**Definition 3.6.4 (Essential crossing).** An external ray $R_t$ is said to cross a crosscut $Q$ essentially if and only if there exists $x \in R_t$ such that the $(T(X), x)$-end of $R_t$ is contained in the bounded complementary domain of $T(X) \cup Q$. In this case we will also say that $Q$ crosses $R_t$ essentially.

The results listed below are known.

**Proposition 3.6.5 ([CL66]).** Let $\mathcal{E}_t$ be a prime end of $U^\infty$. Then $\text{Pr}(\mathcal{E}_t) = \overline{R_t} \setminus R_t$. Moreover, for each $1 < r < \infty$ there is a crosscut $Q_r$ of $U^\infty$ with $\{\phi(r e^{2\pi it})\} = R_t \cap Q_r$ and $\text{diam}(Q_r) \to 0$ as $r \to 1$ and such that $R_t$ crosses $Q_r$ essentially.

**Definition 3.6.6 (Landing Points and Accessible Points).** If $\text{Pr}(\mathcal{E}_t) = \{x\}$, then we say $R_t$ lands on $x \in T(X)$ and $x$ is the landing point of $R_t$. A point $x \in \partial T(X)$ is said to be accessible (from $U^\infty$) if and only if there is an arc in $U^\infty \cup \{x\}$ with $x$ as one of its endpoints.

**Proposition 3.6.7.** A point $x \in \partial T(X)$ is accessible if and only if $x$ is the landing point of some external ray $R_t$.

**Definition 3.6.8 (Channels).** A prime end $\mathcal{E}_t$ of $U^\infty$ for which $\text{Pr}(\mathcal{E}_t)$ is non-degenerate is said to be a channel in $\partial U^\infty$ (or in $T(X)$). If moreover $\text{Pr}(\mathcal{E}_t) = \partial U^\infty = \partial T(X)$, we say $\mathcal{E}_t$ is a dense channel. A crosscut $Q$ of $U^\infty$ is said to cross the channel $\mathcal{E}_t$ if and only if $R_t$ crosses $Q$ essentially.
When $X$ is locally connected, there are no channels, as the following classical theorem proves. In this case, every prime end has degenerate principal set and degenerate impression.

**Theorem 3.6.9 (Carathéodory).** $X$ is locally connected if and only if the Riemann map $\phi : \mathbb{D}^\infty \to U^\infty = \mathbb{C}^\infty \setminus T(X)$ taking $\infty \to \infty$ extends continuously to $\mathbb{S}^1 = \partial \mathbb{D}^\infty$.

### 3.7. Oriented maps

Basic notions of (positively) oriented and confluent maps are defined in Chapter 2. In this section we study (positively) oriented maps and we will establish that it is a natural class of plane maps which are the proper generalization of an orientation preserving homeomorphism of the plane. The following lemmas are in preparation for the proof of Theorem 3.7.4.

**Lemma 3.7.1.** Suppose $f : \mathbb{C} \to \mathbb{C}$ is a perfect surjection. Then $f$ is confluent if and only if $f$ is oriented.

**Proof.** Suppose that $f$ is oriented. Let $A$ be an arc in $\mathbb{C}$ and let $C$ be a component of $f^{-1}(A)$. Suppose that $f(C) \neq A$. Let $a \in A \setminus f(C)$. Since $f(C)$ does not separate $a$ from infinity, we can choose a simple closed curve $S$ with $C \subset T(S)$, $S \cap f^{-1}(A) = \emptyset$ and $f(S)$ so close to $f(C)$ that $f(S)$ does not separate $a$ from $\infty$. Then $a \notin T(f(S))$. Since $f$ is oriented, $f(C) \subset T(f(S))$. Hence there exists a $y \in A \cap f(S)$. This contradicts the fact that $A \cap f(S) = \emptyset$. Thus $f(C) = A$.

Now suppose that $K$ is an arbitrary continuum in $\mathbb{C}$ and let $L$ be a component of $f^{-1}(K)$. Let $x \in L$ and let $A_i$ be a sequence of arcs in $\mathbb{C}$ such that $\lim A_i = K$ and $f(x) \in A_i$ for each $i$. Let $M_i$ be the component of $f^{-1}(A_i)$ containing the point $x$. By the previous paragraph $f(M_i) = A_i$. Since $f$ is perfect, $M = \limsup M_i \subset L$ is a continuum and $f(M) = K$. Hence $f$ is confluent.

Suppose next that $f : \mathbb{C} \to \mathbb{C}$ is not oriented. Then there exists a simple closed curve $S$ in $\mathbb{C}$ and $p \in T(S)$ such that $f(p) \notin T(f(S))$. Let $L$ be a half-line with endpoint $f(p)$ running to infinity in $\mathbb{C} \setminus f(S)$.

Let $L^*$ be an arc in $L$ with endpoint $f(p)$ and diameter greater than the diameter of the continuum $f(T(S))$. Let $K$ be the component of $f^{-1}(L^*)$ which contains $p$. Then $K \subset T(S)$, since $p \in T(S)$ and $L \cap f(S) = \emptyset$. Hence, $f(K) \neq L^*$, and so $f$ is not confluent.

**Lemma 3.7.2.** Let $f : \mathbb{C} \to \mathbb{C}$ be a light, open, perfect surjection. Then there exists an integer $k$ and a finite subset $B \subset \mathbb{C}$ such that $f$ is a local homeomorphism at each point of $\mathbb{C} \setminus B$, and for each point $y \in \mathbb{C} \setminus f(B)$, $|f^{-1}(y)| = k$.

**Proof.** Let $\mathbb{C}^\infty$ be the one point compactification of $\mathbb{C}$. Since $f$ is perfect, we can extend $f$ to a map of $\mathbb{C}^\infty$ onto $\mathbb{C}^\infty$ so that $f^{-1}(\infty) = \infty$. By abuse of notation we also denote the extended map by $f$. Then $f$ is a light open mapping of the compact 2-manifold $\mathbb{C}^\infty$. The result now follows from a theorem of Whyburn [Why42, X.6.3].

The following is a special case, for oriented perfect maps, of the monotone-light factorization theorem. A non-separating plane continuum is said to be *acyclic*.

**Lemma 3.7.3.** Suppose that $f : \mathbb{C} \to \mathbb{C}$ is an oriented, perfect map. It follows that $f = g \circ h$, where $h : \mathbb{C} \to \mathbb{C}$ is a monotone perfect surjection with acyclic fibers and $g : \mathbb{C} \to \mathbb{C}$ is a light, open perfect surjection.
3.7. ORIENTED MAPS

Proof. As above, \( f \) extends to a map of the sphere such that \( f(\infty) = f^{-1}(\infty) = \infty \). By the monotone-light factorization theorem [Nad92, Theorem 13.3], \( f = g \circ h \), where \( h : \mathbb{C} \to X \) is monotone, \( g : X \to \mathbb{C} \) is light, and \( X \) is the quotient space obtained from \( \mathbb{C} \) by identifying each component of \( f^{-1}(y) \) to a point for each \( y \in \mathbb{C} \). Let \( y \in \mathbb{C} \) and let \( C \) be a component of \( f^{-1}(y) \). If \( C \) were to separate \( \mathbb{C} \), then \( f(C) = y \) would be a point while \( f(T(C)) \) would be a non-degenerate continuum. Choose an arc \( A \subseteq \mathbb{C} \setminus \{y\} \) which meets both \( f(T(C)) \) and its complement and let \( x \in T(C) \setminus C \) such that \( f(x) \in A \). If \( K \) is the component of \( f^{-1}(A) \) which contains \( x \), then \( K \subseteq T(C) \). Hence \( f(K) \) cannot map onto \( A \) contradicting the fact that \( f \) is confluent. Thus, for each \( y \in \mathbb{C} \), each component of \( f^{-1}(y) \) is acyclic.

By Moore’s Plane Decomposition Theorem [Dav86], \( X \) is homeomorphic to \( \mathbb{C} \). Since \( f \) is confluent, it is easy to see that \( g \) is confluent. By a theorem of Lelek and Read [LR74] \( g \) is open since it is confluent and light (also see [Nad92, Theorem 13.26]). Since \( h \) and \( g \) factor the perfect map \( f \) through a Hausdorff space \( \mathbb{C} \), both \( h \) and \( g \) are perfect [Eng89, 3.7.5].

Below \( \partial Z \) means the boundary of the set \( Z \).

Theorem 3.7.4 (Maximum Modulus Theorem). Suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a perfect surjection. Then the following are equivalent:

1. \( f \) is either positively or negatively oriented.
2. \( f \) is oriented.
3. \( f \) is confluent.

Moreover, if \( f \) is oriented, then for any non-separating continuum \( X \), \( \partial f(X) \subset f(\partial X) \).

Proof. It is clear that (1) implies (2). By Lemma 3.7.1 every oriented map is confluent. Hence suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a perfect confluent map. By Lemma 3.7.3, \( f = g \circ h \), where \( h : \mathbb{C} \to \mathbb{C} \) is a monotone perfect surjection with acyclic fibers and \( g : \mathbb{C} \to \mathbb{C} \) is a light, open perfect surjection. By Stallings’ Theorem [Whi64] there exists a homeomorphism \( j : \mathbb{C} \to \mathbb{C} \) such that \( g \circ j \) is an analytic surjection. Then \( f = g \circ h = (g \circ j) \circ (j^{-1} \circ h) \). Since \( k = j^{-1} \circ h \) is a monotone surjection of \( \mathbb{C} \) with acyclic fibers, it is a near homeomorphism [Dav86, Theorem 25.1]. That is, there exists a sequence \( k_i \) of homeomorphisms of \( \mathbb{C} \) such that \( \lim k_i = k \). We may assume that all of the \( k_i \) have the same orientation.

Let \( f_i = (g \circ j) \circ k_i \), \( S \) a simple closed curve in the domain of \( f_i \) and \( p \in T(S) \setminus f^{-1}(f(S)) \). Note that \( \lim f_i^{-1}(f_i(S)) \subset f^{-1}(f(S)) \). Hence \( p \in T(S) \setminus f_i^{-1}(f_i(S)) \) for \( i \) sufficiently large. Moreover, since \( f_i \) converges to \( f \), \( f_i|_S \) is homotopic to \( f|_S \) in the complement of \( f(p) \) for \( i \) large. Thus for large \( i \), \( \deg(f_i)_p = \deg(f_p) \), where

\[
(f_i)_p(x) = \frac{f_i(x) - f_i(p)}{f_i(x) - f_i(p)} \quad \text{and} \quad f_p(x) = \frac{f(x) - f(p)}{|f(x) - f(p)|}.
\]

Since \( g \circ j \) is an analytic map, it is positively oriented and we conclude that \( \deg(f_i)_p = \deg(f_p) > 0 \) if \( k_i \) is orientation preserving and \( \deg(f_i)_p = \deg(f_p) < 0 \) if \( k_i \) is orientation reversing. Thus, \( f \) is positively oriented if each \( k_i \) is orientation-preserving and \( f \) is negatively oriented if each \( k_i \) is orientation-reversing.

Suppose that \( X \) is a non-separating continuum and \( f \) is oriented. Let \( y \in \partial f(X) \). Choose \( y_i \in \partial f(X) \) and rays \( R_i \), joining \( y_i \) to \( \infty \) such that \( R_i \cap f(X) = \{y_i\} \) and \( \lim y_i = y \). Choose \( x_i \in X \) such that \( f(x_i) = y_i \). Since \( f \) is confluent, there
exist closed and connected sets $C_i$, joining $x_i$ to $\infty$ such that $C_i \cap X \subset f^{-1}(y_i)$. Hence there exist $x'_i \in f^{-1}(y_i) \cap \partial X$. We may assume that $\lim x'_i = x_\infty \in \partial X$ and $f(x_\infty) = y$ as desired. \hfill \Box

We shall need the following three results in the next section.

**Lemma 3.7.5.** Let $X$ be a plane continuum and $f : \mathbb{C} \to \mathbb{C}$ a perfect, surjective map such that $f^{-1}(T(X)) = T(X)$ (i.e., $T(X)$ is fully invariant) and $f|_{\mathbb{C} \setminus T(X)}$ is confluent. Then for each $y \in \mathbb{C} \setminus T(X)$, each component of $f^{-1}(y)$ is acyclic.

**Proof.** Suppose there exists $y \in \mathbb{C} \setminus T(X)$ such that some component $C$ of $f^{-1}(y)$ is not acyclic. Then there exists $z \in T(C) \cup [f^{-1}(y) \cup T(X)]$. Then $T(X) \cup \{y\}$ does not separate $f(z)$ from infinity in $\mathbb{C}$. Let $L$ be a ray in $\mathbb{C} \setminus [T(X) \cup \{y\}]$ from $f(z)$ to infinity. Then $L = \cup L_i$, where each $L_i \subset L$ is an arc with endpoint $f(z)$. For each $i$ the component $M_i$ of $f^{-1}(L_i)$ containing $z$ maps onto $L_i$. Then $M = \cup M_i$ is a connected closed subset in $\mathbb{C} \setminus f^{-1}(y)$ from $z$ to infinity. This is a contradiction since $z$ is contained in a bounded complementary component of $f^{-1}(y)$.

**Theorem 3.7.6.** Let $X$ be a plane continuum and $f : \mathbb{C} \to \mathbb{C}$ a perfect, surjective map such that $f^{-1}(T(X)) = T(X)$ and $f|_{\mathbb{C} \setminus T(X)}$ is confluent. If $A$ and $B$ are crosscuts of $T(X)$ such that $B \cup X$ separates $A$ from $\infty$ in $\mathbb{C}$, then $f(B) \cup T(X)$ separates $f(A) \cup f(B)$ from $\infty$.

**Proof.** Suppose not. Then there exists a half-line $L$ joining $f(A)$ to infinity in $\mathbb{C} \setminus (f(B) \cup T(X))$. As in the proof of Lemma 3.7.5, there exists a closed and connected set $M \subset \mathbb{C} \setminus (B \cup X)$ joining $A$ to infinity, a contradiction. \hfill \Box

**Proposition 3.7.7.** Under the conditions of Theorem 3.7.6, if $L$ is a ray irreducible from $T(X)$ to infinity, then each component of $f^{-1}(L)$ is closed in $\mathbb{C} \setminus X$ and is a connected set from $X$ to infinity.

### 3.8. Induced maps of prime ends

Suppose that $f : \mathbb{C} \to \mathbb{C}$ is an oriented perfect surjection and $f^{-1}(Y) = X$, where $X$ and $Y$ are acyclic continua and $Y$ has no cutpoints. We will show that in this case the map $f$ induces a confluent map $F$ of the circle of prime ends of $X$ to the circle of prime ends of $Y$. This result was announced by Mayer in the early 1980’s but never appeared in print. It was also used (for homeomorphisms) by Cartwright and Littlewood in [CL51]. There are easy counterexamples that show if $f$ is not confluent then it may not induce a continuous function between the circles of prime ends. For example, if $Y = \mathbb{D}$, $X$ is the union of the unit disk and a copy of a half ray $R$ which spirals to the unit circle and $f$ is radial projection of $R$ onto the unit circle, then $f$ can be extended to a perfect map $F$ of the plane so that $F^{-1}(Y) = X$ but $F$ does not induce a continuous function from the circle of prime ends of $X$ to the circle of prime ends of $Y$.

**Theorem 3.8.1.** Let $X$ and $Y$ be non-degenerate acyclic plane continua and $f : \mathbb{C} \to \mathbb{C}$ a perfect map such that:

1. $Y$ has no cutpoint,
2. $f^{-1}(Y) = X$ and
3. $f|_{\mathbb{C} \setminus X}$ is confluent.
Let \( \varphi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus X \) and \( \psi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus Y \) be conformal mappings. Define \( \tilde{f} : \mathbb{D}^\infty \to \overline{\mathbb{D}^\infty} \) by \( \tilde{f} = \psi^{-1} \circ f \circ \varphi \).

Then \( \tilde{f} \) extends to a map \( \tilde{f} : \mathbb{C}^\infty \to \overline{\mathbb{D}^\infty} \). Moreover, \( \tilde{f}^{-1}(S^1) = S^1 \) and \( F = \tilde{f}|_{S^1} \) is a confluent map.

**Proof.** Note that \( f \) takes accessible points of \( X \) to accessible points of \( Y \). For if \( P \) is a path in \( [C \setminus X] \cup \{p\} \) with endpoint \( p \in X \), then by (2), \( f(P) \) is a path in \( [C \setminus Y] \cup \{ f(p) \} \) with endpoint \( f(p) \in Y \).

Let \( A \) be a crosscut of \( X \) such that the diameter of \( f(A) \) is less than half of the diameter of \( Y \) and let \( U \) be the bounded component of \( C \setminus \{ X \cup A \} \). Let the endpoints of \( A \) be \( x, y \in X \) and suppose that \( f(x) = f(y) \). If \( x \) and \( y \) lie in the same component of \( f^{-1}(f(x)) \) then each crosscut \( B \subseteq U \) of \( X \) is mapped to a generalized return cut of \( Y \) based at \( f(x) \) (i.e., by (1) and (3)) \( f(U) \cap Y = f(x) \) and the endpoints of \( B \) map to \( f(x) \). Note that in this case by (1), \( \partial f(U) \subseteq f(A) \cup \{ f(p) \} \).

Now suppose that \( f(x) = f(y) \) and \( x \) and \( y \) lie in distinct components of \( f^{-1}(f(x)) \). Then by unicoherence of \( C \), \( \partial U \subset A \cup X \) is a connected set and \( \partial U \not\subset A \cup f^{-1}(f(x)) \). Now \( \partial U \setminus (A \cup f^{-1}(f(x))) = \partial U \setminus f^{-1}(f(A)) \) is an open non-empty set in \( \partial U \). Thus there is a crosscut \( B \subseteq U \) such that \( B \setminus \partial U \not\subseteq f^{-1}(f(A)) \). Now \( f(B) \) is contained in a bounded component of \( C \setminus f(A) \) by Theorem 3.7.6. Since \( Y \cap f(A) = \{ f(x) \} \), \( X \) is mapped to a point \( y \). Since \( Y \setminus \{ f(x) \} \) meets \( f(B) \) and \( f(A) \), it follows by unicoherence that \( f(B) \) lies in a bounded component of \( C \setminus f(A) \). This is impossible as the diameter of \( f(A) \) is smaller than the diameter of \( Y \). It follows that there exists a \( \delta > 0 \) such that if the diameter of \( A \) is less than \( \delta \) and \( f(x) = f(y) \), then \( x \) and \( y \) must lie in the same component of \( f^{-1}(f(x)) \).

In order to define the extension \( \tilde{f} \) of \( f \) over the boundary \( S^1 \) of \( \mathbb{D}^\infty \), let \( C_i \) be a chain of crosscuts of \( \mathbb{D}^\infty \) which converge to a point \( p \in S^1 \) such that \( A_i = \varphi(C_i) \) is a null sequence of crosscuts or return cuts of \( X \) with endpoints \( a_i \) and \( b_i \) which converge to a point \( x \in X \). There are three cases to consider:

Case 1. \( f \) identifies the endpoints of \( A_i \) for some \( A_i \) with diameter less than \( \delta \). In this case the chain of crosscuts is mapped by \( f \) to a sequence of generalized return cuts based at \( f(a_i) = f(b_i) = f(x) \). Hence \( f(a_i) \) is an accessible point of \( Y \) which corresponds (under \( \psi^{-1} \)) to a unique point \( q \in S^1 \) (since \( Y \) has no cutpoints). Define \( \tilde{f}(p) = q \).

Case 2. Case 1 does not apply and there exists an infinite subsequence \( A_{i_j} \) of crosscuts such that \( f(A_{i_j}) \cap f(A_{j_i}) = \emptyset \) for \( j \neq k \). In this case \( f(A_{i_j}) \) is a chain of generalized crosscuts which converges to the point \( f(x) \in Y \). The chain \( \psi^{-1} \circ f(A_{i_j}) \) corresponds to a unique point \( q \in S^1 \). Define \( \tilde{f}(p) = q \).

Case 3. Cases 1 and 2 do not apply. Without loss of generality suppose there exists an \( i \) such that for \( j > i \), \( f(A_{i_j}) \cap f(A_{j_i}) \) contains \( f(a_i) = f(x) \). In this case \( f(A_{i_j}) \) is a sequence of generalized crosscuts based at the accessible point \( f(x) \). Define \( \tilde{f}(p) = q \).

It remains to be shown that \( \tilde{f} \) is a continuous extension of \( \tilde{f} \) and \( F \) is confluent. For continuity it suffices to show continuity at \( S^1 \). Let \( p \in S^1 \) and let \( C \) be a small crosscut of \( \mathbb{D}^\infty \) whose endpoints are on opposite sides of \( p \) in \( S^1 \) such that \( A = \varphi(C) \) has diameter less than \( \delta \) [Mil00] and such that the endpoints of \( A \) are two accessible points of \( X \). Since \( f \) is uniformly continuous near \( X \), the diameter of \( f(A) \) is small and since \( \psi^{-1} \) is uniformly continuous with respect to connected
sets in the complement of $Y$ ([UY51]), the diameter of $B = \psi^{-1} \circ f \circ \varphi(C)$ is small. Also $B$ is either a generalized crosscut or generalized return cut. Since $\hat{f}$ preserves separation of crosscuts, it follows that the image of the domain $\hat{U}$ bounded by $C$ which does not contain $\infty$ is small. This implies continuity of $\bar{f}$ at $p$.

To see that $F$ is confluent let $K \subset S^1$ be a subcontinuum and let $H$ be a component of $\bar{f}^{-1}(K)$. Choose a chain of crosscuts $C_i$ such that $\varphi(C_i) = A_i$ is a crosscut of $X$ meeting $X$ in two accessible points $a_i$ and $b_i$, $C_i \cap \bar{f}^{-1}(K) = \emptyset$ and $\lim C_i = H$. It follows from the preservation of crosscuts (see Theorem 3.7.6) that $\hat{f}(C_i)$ separates $K$ from $\infty$. Hence $\hat{f}(C_i)$ must meet $S^1$ on both sides of $K$ and $\lim \bar{f}(C_i) = K$. Hence $F(H) = \lim \bar{f}(C_i) = K$ as required.

**Corollary 3.8.2.** Suppose that $f : \mathbb{C} \to \mathbb{C}$ is a perfect, oriented map of the plane, $X \subset \mathbb{C}$ is a subcontinuum without cut points and $f(X) = X$. Let $\hat{X}$ be the component of $f^{-1}(f(X))$ containing $X$. Let $\varphi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus T(\hat{X})$ and $\psi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus T(X)$ be conformal mappings. Define $\hat{f} : \mathbb{D}^\infty \setminus \varphi^{-1}(f^{-1}(X)) \to \mathbb{D}^\infty$ by $\hat{f} = \psi^{-1} \circ f \circ \varphi$. Put $S^1 = \partial \mathbb{D}^\infty$.

Then $\hat{f}$ extends over $S^1$ to a map $f : \mathbb{D}^\infty \to \mathbb{D}^\infty$. Moreover $f^{-1}(S^1) = S^1$ and $F = f|_{S^1}$ is a confluent map.

**Proof.** By Lemma 3.7.3 $f = g \circ m$ where $m$ is a monotone perfect and onto mapping of the plane with acyclic point inverses, and $g$ is an open and perfect surjection of the plane to itself. By Lemma 3.7.2, $f^{-1}(X)$ has finitely many components. It follows that there exist a simply connected open set $V$, containing $T(X)$, such that if $U$ is the component of $f^{-1}(V)$ containing $\hat{X}$, then $U$ contains no other components of $f^{-1}(X)$. It is easy to see that $f(U) = V$ and that $U$ is simply connected. Hence $U$ and $V$ are homeomorphic to $\mathbb{C}$. Then $f|_U : U \to V$ is a confluent map. The result now follows from Theorem 3.8.1 applied to $f$ restricted to $U$. \qed
CHAPTER 4

Partitions of domains in the sphere

4.1. Kulkarni-Pinkall Partitions

Throughout this section let $K$ be a compact subset of the plane whose complement $U = \mathbb{C} \setminus K$ is connected. In the interest of completeness we define the Kulkarni-Pinkall partition of $U$ and prove the basic properties of this partition that are essential for our work in Section 4.2. Kulkarni-Pinkall [KP94] worked in closed $n$-manifolds. We will follow their approach and adapt it to our situation in the plane.

We think of $K$ as a closed subset of the Riemann sphere $\mathbb{C}^\infty$, with the spherical metric and set $U^\infty = \mathbb{C}^\infty \setminus K = U \cup \{\infty\}$. Let $\mathfrak{B}^\infty$ be the family of closed, round balls $B$ in $\mathbb{C}^\infty$ such that $\text{Int}(B) \subset U^\infty$ and $|\partial B \cap K| \geq 2$. Then $\mathfrak{B}^\infty$ is in one-to-one correspondence with the family $\mathfrak{B}$ of closed subsets $B$ of $\mathbb{C}$ which are the closure of a complementary component of a straight line or a round circle in $\mathbb{C}$ such that $\text{Int}(B) \subset U$ and $|\partial B \cap K| \geq 2$.

**Proposition 4.1.1.** If $B_1$ and $B_2$ are two closed round balls in $\mathbb{C}$ such that $B_1 \cap B_2 \neq \emptyset$ but does not contain a diameter of either $B_1$ or $B_2$, then $B_1 \cap B_2$ is contained in a ball of diameter strictly less than the diameters of both $B_1$ and $B_2$.

**Proof.** Let $\partial B_1 \cap \partial B_2 = \{s_1, s_2\}$. Then the closed ball with center $(s_1 + s_2)/2$ and radius $|s_1 - s_2|/2$ contains $B_1 \cap B_2$. \qed

If $B$ is the closed ball of minimum diameter that contains $K$, then we say that $B$ is the smallest ball containing $K$. It is unique by Proposition 4.1.1. It exists, since any sequence of balls of decreasing diameters that contain $K$ has a convergent subsequence.

We denote the *Euclidean convex hull* of $K$ by $\text{conv}_E(K)$. It is the intersection of all closed half-planes (a closed half-plane is the closure of a component of the complement of a straight line) which contain $K$. Hence $p \in \text{conv}_E(K)$ if $p$ cannot be separated from $K$ by a straight line.

Given a closed ball $B \in \mathfrak{B}^\infty$, $\text{int}(B)$ is conformally equivalent to the unit disk in $\mathbb{C}$. Hence its interior can be naturally equipped with the hyperbolic metric. Geodesics $g$ in this metric are intersections of $\text{int}(B)$ with round circles $C \subset \mathbb{C}^\infty$ which perpendicularly cross the boundary $\partial B$. For every hyperbolic geodesic $g$, $B \setminus g$ has exactly two components. We call the closure of such components *hyperbolic half-planes* of $B$. Given $B \in \mathfrak{B}^\infty$, the *hyperbolic convex hull* of $K$ in $B$ is the intersection of all (closed) hyperbolic half-planes of $B$ which contain $K \cap B$ and we denote it by $\text{conv}_H(B \cap K)$.

**Lemma 4.1.2.** Suppose that $B$ is the smallest ball containing $K \subset \mathbb{C}$ and let $c \in B$ be its center. Then $c \in \text{conv}_H(K \cap \partial B)$. 35
PARTITIONS OF DOMAINS IN THE SPHERE

Figure 4.1. Maximal balls have disjoint hulls.

Proof. By contradiction. Suppose that there exists a circle that separates the center $c$ from $K \cap \partial B$ and crosses $\partial B$ perpendicularly. Then there exists a line $\ell$ through $c$ such that a half-plane bounded by $\ell$ contains $K \cap \partial B$ in its interior. Let $B' = B + v$ be a translation of $B$ by a vector $v$ that is orthogonal to $\ell$ and directed into this halfplane. If $v$ is sufficiently small, then $B'$ contains $K$ in its interior. Hence, it can be shrunk to a strictly smaller ball that also contains $K$, contradicting that $B$ has smallest diameter.

Lemma 4.1.3. Suppose that $B_1, B_2 \in \mathfrak{B}^\infty$ with $B_1 \neq B_2$. Then

$$\text{conv}_H(B_1 \cap \partial U) \cap \text{conv}_H(B_2 \cap \partial U) \subset \partial U.$$  

In particular, $\text{conv}_H(B_1 \cap \partial U) \cap \text{conv}_H(B_2 \cap \partial U)$ contains at most two points.

Proof. A picture easily explains this, see Figure 4.1. Note that $\partial U \cap [B_1 \cup B_2] \subset \partial (B_1 \cup B_1)$. Therefore $B_1 \cap \partial U$ and $B_2 \cap \partial U$ share at most two points. The open hyperbolic chords between these points in the respective balls are disjoint.

It follows that any point in $U^\infty(X)$ can be contained in at most one hyperbolic convex hull. In the next lemma we see that each point of $U^\infty$ is indeed contained in $\text{conv}_H(B \cap K)$ for some $B \in \mathfrak{B}^\infty$. So $\{U^\infty \cap \text{conv}_H(B \cap K) \mid B \in \mathfrak{B}^\infty\}$ is a partition of $U^\infty$.

Since hyperbolic convex hulls are preserved by M"obius transformations, they are more easy to manipulate than the Euclidean convex hulls used by Bell (which are preserved only by M"obius transformations that fix $\infty$). This is illustrated by the proof of the following lemma.

Lemma 4.1.4 (Kulkarni-Pinkall inversion lemma). For any $p \in \mathfrak{C}^\infty \setminus K$ there exists $B \in \mathfrak{B}^\infty$ such that $p \in \text{conv}_H(B \cap K)$.

Proof. We prove first that there exists $B^* \in \mathfrak{B}^\infty$ such that no circle which crosses $\partial B^*$ perpendicularly separates $K \cap \partial B^*$ from $\infty$. 


Let \( B' \) be the smallest round ball which contains \( K \) and let \( B = \overline{C \setminus B'} \). Then \( B' = B \cup \{ \infty \} \in \mathcal{B}^\infty \). If \( L \) is a circle which crosses \( \partial B' \) perpendicularly and separates \( K \cap \partial B' \) from \( \infty \), then it also separates \( K \cap \partial B' \) from the center \( c' \) of \( B' \), contrary to Lemma 4.1.2. [To see this note that if \( F \) is the Möbius transformation which fixes points in the boundary of \( B' \) and interchanges the points \( \infty \) and \( c' \), then \( F(L) = L \). Hence it would follow that \( L \) separates \( c' \) from \( K \cap \partial B' \), a contradiction with Lemma 4.1.2.] Hence, \( \infty \in \text{conv}_H(B' \cap K) \).

Now let \( p \in \mathbb{C}^\infty \setminus K \). Let \( M : \mathbb{C}^\infty \to \mathbb{C}^\infty \) be a Möbius transformation such that \( M(p) = \infty \). By the above argument there exists a ball \( B^* \in \mathcal{B}^\infty \) such that \( \infty \in \text{conv}_H(B^* \cap M(K)) \). Then \( B = M^{-1}(B^*) \in \mathcal{B}^\infty \) and, since \( M \) preserves perpendicular circles, \( p \in \text{conv}_H(B \cap K) \) as desired.

From Lemmas 4.1.3 and 4.1.4, we obtain the following Theorem which is a special case of a Theorem of Kulkarni and Pinkall [KP94].

**Theorem 4.1.5.** Suppose that \( K \subset \mathbb{C} \) is a nondegenerate compact set such that its complement \( U^\infty \) in the Riemann sphere is connected. Then \( U^\infty \) is partitioned by the family

\[
\mathcal{KPP} = \{ U^\infty \cap \text{conv}_H(B \cap K) : B \in \mathcal{B}^\infty \}.
\]

Theorem 4.1.5 is the linchpin of the theory of geometric crosscuts. An analogue of it was known to Harold Bell and used by him implicitly since the early 1970’s. Bell considered non-separating plane continua \( K \) and he used the equivalent notion of Euclidean convex hull of the sets \( B \cap \partial U \) for all maximal balls \( B \in \mathcal{B} \) (see the comment following Theorem 4.2.5).

Let \( B \in \mathcal{B}^\infty \). If \( B \cap \partial U^\infty(X) \) consists of two points \( a \) and \( b \), then its (hyperbolic) hull is an open circular segment \( g \) with endpoints \( a \) and \( b \) and perpendicular to \( \partial B \). We will call the crosscut \( g \) a \( \mathcal{KPP} \) crosscut or simply a \( \mathcal{KPP} \)-chord. If \( B \cap \partial U \) contains three or more points, then we say that the hull \( \text{conv}_H(B \cap \partial U) \) is a gap. A gap has nonempty interior. Its boundary in \( \text{int}(B) \) is a union of open circular segments (with endpoints in \( K \)), which we also call \( \mathcal{KPP} \) crosscuts or \( \mathcal{KPP} \)-chords. We denote by \( \mathcal{KPP} \) the collection of all open chords obtained as above using all \( B \in \mathcal{B}^\infty \).

The following example may serve to illustrate Theorem 4.1.5.

**Example.** Let \( K \) be the unit square \( \{ x + yi : -1 \leq x, y \leq 1 \} \). There are five obvious members of \( \mathcal{B} \). These are the sets

\[
\text{Im} z \geq 1, \quad \text{Im} z \leq -1, \quad \text{Re} z \geq 1, \quad \text{Re} z \leq -1, \quad |z| \geq \sqrt{2},
\]

four of which are half-planes. These are the only members of \( \mathcal{B} \) whose hyperbolic convex hulls have non-empty interiors. However, for this example the family \( \mathcal{B} \) defined in the introduction of Section 4.1 is infinite. The hyperbolic hull of the half-plane \( \text{Im} z \geq 1 \) is the semi-disk \( \{ z : |z - i| \leq 1, \ \text{Im} z > 1 \} \). The hyperbolic hulls of the other three half-planes given above are also semi-disks. The hyperbolic hull of \( |z| \geq \sqrt{2} \) is the unbounded region whose boundary consists of parts of four circles lying (except for their endpoints) outside \( K \) and contained in the circles of radius \( \sqrt{2} \) and having centers at \(-2, 2, -2i\) and \(2i\), respectively. These hulls do not cover \( U \) as there are spaces between the hulls of the half-planes and the hull of \( |z| \geq \sqrt{2} \).

If \( C \) is a circle that circumscribes \( K \) and contains exactly two of its vertices, such as \( 1 \pm i \), then the exterior ball \( B \) bounded by \( C \) is maximal. Now \( \text{conv}_H(B \cap K) \)
is a single chord and the union of all such chords foliates the remaining spaces in $C \setminus K$.

**Lemma 4.1.6.** If $g_i$ is a sequence of $\mathcal{KP}$-chords with endpoints $a_i$ and $b_i$, and $\lim a_i = a \neq b = \lim b_i$, then $\{g_i\}$ has a convergent subsequence and $\lim g_i = C$, where $g = C \setminus \{a, b\} \in \mathcal{KP}$ is also a $\mathcal{KP}$-chord.

**Proof.** For each $i$ let $B_i \in \mathcal{B}^\infty$ such that $g_i \subset \text{conv}_\mathcal{H}(B_i \cap K)$. Then a subsequence $B_{i_j}$ converges to some $B \in \mathcal{B}^\infty$ and $g_{i_j}$ converges to a closed circular arc $C$ in $B$ with endpoints $a$ and $b$, and $C$ is perpendicular to $\partial B$. Hence $g = C \setminus \{a, b\} \subset \text{conv}_\mathcal{H}(B \cap K)$. So $g \in \mathcal{KP}$. □

By Lemma 4.1.6, the family $\mathcal{KP}$ of chords has continuity properties similar to a foliation.

**Lemma 4.1.7.** For $a, b \in \partial U^\infty$, define $C(a, b)$ as the union of all $\mathcal{KP}$-chords with endpoints $a$ and $b$. Then if $C(a, b) \neq \emptyset$, $C(a, b)$ is either a single chord, or $C(a, b) \cup \{a, b\}$ is a closed disk whose boundary consists of two $\mathcal{KP}$-chords contained in $C(a, b)$ together with $\{a, b\}$.

**Proof.** Suppose $g$ and $h$ are two distinct $\mathcal{KP}$-chords between $a$ and $b$. Then $S = g \cup h \cup \{a, b\}$ is a simple closed curve. Choose a point $z$ in the complementary domain $V$ of $S$ contained in $U^\infty$. Since the hyperbolic hulls partition $U^\infty$, there exists $B \in \mathcal{B}^\infty$ such that $z \in \text{conv}_\mathcal{H}(B \cap K)$. By Lemma 4.1.3, $\text{conv}_\mathcal{H}(B \cap K)$ can only intersect $S \cap K$ in $\{a, b\}$. So $\text{conv}_\mathcal{H}(B \cap K) \cap K = \{a, b\}$ and it follows that $V$ is contained in $C(a, b)$.

The rest of the Lemma follows from 4.1.6. □

**4.2. Hyperbolic foliation of simply connected domains**

In this section we will apply the results from Section 4.1 to the case that $K$ is a non-separating plane continuum (or, equivalently, that $U^\infty = C^\infty \setminus K$ is simply connected). The results in this section are essential to [OT07, OV09] but are not used in this paper. The reader who is only interested in the fixed point question can skip this section.

Let $\mathbb{D}$ be the open unit disk in the plane. In this section we let $\phi : \mathbb{D} \to C^\infty \setminus K = U^\infty$ be a Riemann map onto $U^\infty$. We endow $\mathbb{D}$ with the hyperbolic metric, which is carried to $U^\infty$ by the Riemann map. We use $\phi$ and the Kulkarni-Pinkall hulls to induce a closed collection $\Gamma$ of chords in $\mathbb{D}$ that is a hyperbolic geodesic lamination in $\mathbb{D}$ (see [Thu09]).

Let $g \in \mathcal{KP}$ be a chord with endpoints $a$ and $b$. Then $a$ and $b$ are accessible points in $K$ and $\phi^{-1}(g)$ is an arc in $\mathbb{D}$ with endpoints $z, w \in \partial \mathbb{D}$. Let $G$ be the hyperbolic geodesic in $\mathbb{D}$ joining $z$ and $w$. Then $G$ is an open circular arc which meets $\partial \mathbb{D}$ perpendicularly. Let $\Gamma$ be the collection of all $K$ such that $g \in \mathcal{KP}$. We will prove that $\Gamma$ inherits the properties of the family $\mathcal{KP}$ as described in Theorem 4.1.5 and Lemma 4.1.6 (see Lemma 4.2.3, Theorem 4.2.5 and the remark following 4.2.5).

Since members of $\mathcal{KP}$ do not intersect (though their closures are arcs which may have common endpoints) the same is true for distinct members of $\Gamma$. We will refer to the members of $\Gamma$ (and their images under $\phi$) as **hyperbolic chords** or **hyperbolic geodesics**. Given $g \in \mathcal{KP}$ we denote the corresponding element of $\Gamma$ by $G$ and its image $\phi(G)$ in $U^\infty$ by $g$. Note that $\Gamma$ is a lamination of $\mathbb{D}$ in the sense of
Let $D$. We may assume that the Riemann map $\phi$ of the Kulkarni-Pinkall partition. That this is indeed the case is the substance of the next lemma. The following lemma follows.

**Lemma 4.2.1** (Jørgensen [Pom92, p.91 and 93]). Let $B$ be a closed round ball such that its interior is in $U^\infty$. Let $\gamma \subset D$ be a hyperbolic geodesic. Then $\phi(\gamma) \cap B$ is connected. In particular, if $R_t$ is an external ray in $U^\infty$ and $B \in B^\infty$, then $R_t \cap B$ is connected.

If $a, b \in \partial U^\infty$, recall that $C(a, b)$ is the union of all $KP$-chords with endpoints $a$ and $b$. From the viewpoint of prime ends, all chords in $C(a, b)$ are the same. That is why all the chords in $C(a, b)$ are replaced by a single hyperbolic chord $g \in \phi(\Gamma)$. The following lemma shows this.

**Lemma 4.2.2.** Suppose $g \in KP$ and $g \subset \text{conv}_H(B \cap \partial U^\infty)$ joins the points $a, b \in \partial U^\infty$ for some $B \in B^\infty$. If $G \in \Gamma$ is the corresponding hyperbolic geodesic, then $g = \phi(G) \subset B$.

**Proof.** We may assume that the Riemann map $\phi : D \to U^\infty$ is extended over all points $x \in S^1$ so that $\phi(x)$ is an accessible point of $U^\infty$. Let $\phi^{-1}(a) = \hat{a}$, $\phi^{-1}(b) = \hat{b}$ and $\phi^{-1}(B) = \hat{B}$, and let $G$ be the hyperbolic geodesic joining the points $\hat{a}$ and $\hat{b}$ in $D$. (Note that this extended map is not necessarily continuous at points of $D$ corresponding to accessible points of $K$.) Suppose, by way of contradiction, that $x \in G \cap \hat{B}$. Let $C$ be the component of $D \setminus \phi^{-1}(g)$ which does not contain $x$. Choose $a_i \to \hat{a}$ and $b_i \to \hat{b}$ in $S^1 \cap C$ and let $H_i$ be the hyperbolic geodesic in $D$ joining the points $a_i$ and $b_i$. Then $\lim H_i = G$ and $H_i \cap \hat{B}$ is not connected for $i$ large. Hence $\phi(H_i) \cap B$ is not connected for $i$ large. This contradiction with Lemma 4.2.1 completes the proof. \hfill $\square$

**Lemma 4.2.3.** Suppose that $\{G_i\}$ is a sequence of hyperbolic chords in $\Gamma$ and suppose that $x_i \in G_i$ such that $\{x_i\}$ converges to $x \in D$. Then there is a unique hyperbolic chord $G \in \Gamma$ that contains $x$. Furthermore, $\lim G_i = \overline{G}$.

**Proof.** We may suppose that a subsequence sequence $\{G_{i_j}\}$ converges to a hyperbolic chord $G$ which contains $x$. Let $g_{i_j} \in KP$ so that $\phi^{-1}(g_{i_j})$ is an open arc which joins the endpoints of $G_i$. By Lemma 4.1.6, there is another subsequence so that $\lim g_{i_{j(i)}} = g \in KP$. It follows that $G$ is the hyperbolic chord joining the endpoints of $\phi^{-1}(g)$. Hence $g \in \Gamma$. Since the above argument applies to all subsequences, the sequence $G_i$ must converge to $G$. \hfill $\square$

So we have used the family of $KP$-chords in $U^\infty$ to stratify $D$ to the family $\Gamma$ of hyperbolic chords. In particular gaps of $\Gamma$ are no longer necessarily disjoint but they can meet at most in a common boundary chord. By Lemma 4.2.2 for each $KP$-chord $g \subset \text{conv}_H(B \cap \partial U^\infty)$ its associated hyperbolic chord $g = \phi(G) \subset B$. Hence, there is a continuous deformation of $U^\infty$ that maps $\bigcup KP$ onto $\bigcup \phi(\Gamma)$, which suggests that components of $U^\infty \setminus \bigcup \phi(\Gamma)$ naturally correspond to the interiors of the gaps of the Kulkarni-Pinkall partition. That this is indeed the case is the substance of the next lemma.

**Lemma 4.2.4.** There is a 1–1 correspondence between complementary domains $Z \subset D \setminus \bigcup \Gamma$ and the interiors of Kulkarni-Pinkall gaps $\text{conv}_H(B \cap K)$. Moreover, for each gap $Z$ of $\Gamma$ there exists a unique $B \in B^\infty$ such that $Z$ corresponds to the interior of the $KP$-gap $\text{conv}_H(B \cap K) \cap U^\infty$ in that $\partial_Z \cap D = \bigcup \{G \in \Gamma \mid g \in KP \text{ and } g \subset \partial \text{conv}_H(B \cap K)\}$ and $\phi(Z) \subset B$. 

Thurston[Thu09]. By a gap of $\Gamma$ (or of $\phi(\Gamma)$), we mean the closure of a component of $D \setminus \bigcup \Gamma$ in $D$ (or its image under $\phi$ in $U^\infty$, respectively).
Proof. Let $g$ and $h$ be two distinct $\mathcal{KP}$-chords in the boundary of the gap $\text{conv}_H(B \cap K)$ for some $B \in \mathfrak{B}^\infty$. Let $\{a, b\}$ and $\{c, d\}$ be the endpoints of $\phi^{-1}(g)$ and $\phi^{-1}(h)$, respectively. Since $g$ and $h$ are contained in the same gap, no hyperbolic chord of $\Gamma$ separates $G$ and $H$. Hence there exists a gap $Z$ of $\Gamma$ whose boundary includes the hyperbolic chords $G$ and $H$. It now follows easily that for any $g' \in \mathcal{KP}$ which is contained in the boundary of the same gap $\text{conv}_H(B \cap K)$, $G'$ is contained in the boundary of $Z$. Hence the $\mathcal{KP}$ gap $\text{conv}_H(B \cap K)$ corresponds to the gap $Z$ of $\Gamma$. Conversely, if $Z$ is a gap of $\Gamma$ in $\mathbb{D}$ then a similar argument, together with Lemmas 4.1.6 and 4.1.7, implies that $Z$ corresponds to a unique gap $\text{conv}_H(B \cap K)$ for some $B \in \mathfrak{B}^\infty$. The rest of the Lemma now follows from Lemma 4.2.2.

So if $U^\infty = \mathbb{C}^\infty \setminus K$ is endowed with the hyperbolic metric induced by $\phi$, then there exists a family of geodesic chords that share the same endpoints as elements of $\mathcal{KP}$. The complementary domains of $U^\infty \setminus \bigcup \{g \mid g \in \mathcal{KP}\}$ corresponds to the Kulkarni-Pinkall gaps. We summarize the results:

**Theorem 4.2.5.** Suppose that $K \subset \mathbb{C}$ is a non-separating continuum and let $U^\infty$ be its complementary domain in the Riemann sphere. There exists a family $\phi(\Gamma)$ of hyperbolic chords in the hyperbolic metric on $U^\infty$ such that for each $g \in \phi(\Gamma)$ there exists $B \in \mathfrak{B}^\infty$ and $g \subset \text{conv}_H(B \cap \partial U^\infty)$ so that $g$ and $g'$ have the same endpoints and $g \subset B$. Each domain $Z$ of $U^\infty \setminus \phi(\Gamma)$ naturally corresponds to a Kulkarni-Pinkall gap $\text{conv}_H(B \cap \partial U^\infty)$ The bounding hyperbolic chords of $Z$ in $U^\infty$ correspond to the $\mathcal{KP}$-chords (i.e., chords in $\mathcal{KP}$) of $\text{conv}_H(B \cap \partial U^\infty)$.

In order to obtain Bell’s Euclidean foliation [Bel76] we could have modified the $\mathcal{KP}$ family as follows. Suppose that $B \in \mathfrak{B}$. Instead of replacing a $\mathcal{KP}$-chord $g \in \text{conv}_H(B \cap K)$ by a geodesic in the hyperbolic metric on $U^\infty$, we could have replaced it by a straight line segment; i.e., the geodesic in the Euclidean metric. Then we would have obtained a family of open straight line segments. In so doing we would have replaced the gaps $\text{conv}_H(B \cap \partial U^\infty)$ by $\text{conv}_E(B \cap \partial U^\infty)$, which is the way in which Bell originally foliated $\text{conv}_E(K) \setminus K$. We hope that the above argument provides a more transparent proof of Bell’s result. Note that both in the hyperbolic and Euclidean case the elements of the foliation are not necessarily disjoint (hence we use the word “foliate” rather than “partition”). However, in both cases every point of $U^\infty$ is contained in either a unique chord or in the interior of a unique gap.

### 4.3. Schoenflies Theorem

In this short section we will show that the Schoenflies Theorem follows immediately from Theorem 4.2.5 (see [Sie05] for some recent history of this old problem and [OT07, OV09] for more details and extensions of these ideas). We want to emphasize here that no results of Chapter 3 are relied upon in Section 4.3.

**Theorem 4.3.1** (Schoenflies Theorem). Suppose that $h : S^1 \to \mathbb{C}$ is an embedding of the unit circle in the plane and $U$ is a bounded complementary domain of $h(S^1) = S$. Then there exists an embedding $H : \mathbb{D} \to \mathbb{C}$ which extends $h$.

**Proof.** Let $\mathfrak{B} = \{B_\alpha\}$ be the collection of maximal open balls in $U$ such that $|\partial B_\alpha \cap \partial U| \geq 2$. For each $\alpha$ let $F_\alpha = \text{conv}_E(\partial B_\alpha \cap \partial U)$. Let $\mathcal{L}$ be the collection of all chords in the boundaries of all the sets $F_\alpha$ and let $\mathcal{L}^*$ be the union of all
Let $\alpha$ there exists a unique $\text{var}(z)$ that $X$ elements in $U$ Kulkarni Pinkall partition of $U$ minimal with respect to these properties. We apply the Kulkarni-Pinkall partition $h$ joining the points $H$ each $w$ by Theorem 4.2.5 there exists for each $z$ the chords in $L$ which crosses $\partial B$ $\in$ contains $\subset\leq \text{var}(z)$ of diameter $KPP$ circles on the sphere containing the point at infinity. The subfamily of $\text{var}(z)$ is a plane continuum. Here we assume, as in the introduction to this paper, that $X$ be its barycenter. Then it follows easily that we can extend the map $h$ continuously by defining $H(h_{\beta}) = g_{\beta}$ for each $\beta$. Finally extend $H$ over all of $H_{\beta}$ by mapping, for each $w \in \partial H_{\beta}$ the straight line segment $w_{h_{\beta}}$ linearly onto the straight line segment joining the points $H(w)$ and $H(h_{\beta}) = g_{\beta}$. Then $H$ is the required extension of $h$.

4.4. Prime ends

We will follow the notation from Section 4.1 in the case that $K = T(X)$ where $X$ is a plane continuum. Here we assume, as in the introduction to this paper, that $f: \mathbb{C} \to \mathbb{C}$ takes the continuum $X$ into $T(X)$ with no fixed points in $T(X)$, and $X$ is minimal with respect to these properties. We apply the Kulkarni-Pinkall partition to $U^\infty = \mathbb{C}^\infty \setminus T(X)$. Recall that $KPP = \{\text{conv}_H(B \cap K) \cap U^\infty \mid B \in \mathcal{B}^\infty\}$ is the Kulkarni Pinkall partition of $U^\infty$ as given by Theorem 4.1.5.

Let $B^\infty \in \mathcal{B}^\infty$ be the maximal ball such that $\infty \in \text{conv}_H(B^\infty \cap K)$. As before we use balls on the sphere. In particular, straight lines in the plane correspond to circles on the sphere containing the point at infinity. The subfamily of $\text{var}(z)$ whose elements are of diameter $\leq \delta$ in the spherical metric is denoted by $\text{KPP}_\delta$. The subfamily of chords in $\text{KPP}$ of diameter $\leq \delta$ is denoted by $\text{KPP}_\delta$.

By Lemma 4.1.6 we know that the families $\text{KPP}$ and $\text{KPP}_\delta$ have nice continuity properties. However, $\text{KPP}$ and $\text{KPP}_\delta$ are not closed in the hyperspace of compact subsets of $\mathbb{C}^\infty$: a sequence of chords or hulls may converge to a point in the boundary of $U^\infty$ (in which case it must be a null sequence).

**Proposition 4.4.1 (Closedness).** Let $\{g_i\}$ be a convergent sequence of distinct elements in $\text{KPP}_\delta$, then either $g_i$ converges to a chord $g$ in $\text{KPP}_\delta$ or $g_i$ converges to a point of $X$. In the first case, for large $i$ and $\delta$ sufficiently small, $\text{var}(f,g,T(X)) = \text{var}(f,g,T(X))$.

**Proof.** By Lemma 4.1.6, we know that the first conclusion holds if $g = \lim g_i$ contains a point of $U^\infty$. Hence we only need to consider the case when $\lim g_i = g \subset \partial U^\infty \subset T(X)$. If the diameter of $g_i$, converged to zero, then $g$ is a point as desired. Assume that this is not the case and let $B_i$ be the maximal ball that contains $g_i$. Under our assumption, the diameters of $\{B_i\}$ do not decay to zero. Let $B \in \mathcal{B}^\infty$ be the limit of a subsequence $B_i$. Then $\lim g_i$ is a piece of a round circle which crosses $\partial B$ perpendicularly. Hence $\lim g_i \cap \text{int}(B) \neq \emptyset$, contradicting the fact
that \( g \subset \partial U_\infty \subset T(X) \). Note that for \( \delta \) sufficiently small, \( g \cap f(\overline{g}) = \emptyset \). Hence, \( \text{var}(f, g, T(X)) \) and \( \text{var}(f, g, T(X)) \) are defined for all \( i \) sufficiently large. Then last statement in the Lemma follows from stability of variation (see Section 2.2).

**Corollary 4.4.2.** For each \( \varepsilon > 0 \), there exist \( \delta > 0 \) such that for all \( g \in \mathcal{KP} \) with \( g \subset B(T(X), \delta) \), \( \text{diam}(g) < \varepsilon \).

**Proof.** Suppose not, then there exist \( \varepsilon > 0 \) and a sequence \( g_i \) in \( \mathcal{KP} \) such that \( \lim \text{diam}(g_i) \geq \varepsilon \) a contradiction to Proposition 4.4.1. \( \square \)

The proof of the following well-known proposition is omitted.

**Proposition 4.4.3.** For each \( \varepsilon > 0 \) there exists \( \delta > 0 \) such that for each open arc \( A \) with distinct endpoints \( a,b \) such that \( \overline{A} \cap T(X) = \{a,b\} \) and \( \text{diam}(A) < \delta \), \( T(T(X) \cup A) \subset B(T(X), \varepsilon) \).

**Proposition 4.4.4.** Let \( \varepsilon, \delta \) be as in Proposition 4.4.3 above with \( \delta < \varepsilon/2 \) and let \( B \in \mathcal{B}_\infty \). Let \( A \) be a crosscut of \( T(X) \) such that \( \text{diam}(A) < \delta \). If \( x \in T(A \cup T(X)) \cap \text{conv}_H (B \cap T(X)) \setminus T(X) \) and \( d(x,A) \geq \varepsilon \), then the radius of \( B \) is less than \( \varepsilon \). Hence, \( \text{diam}(\text{conv}_H (B \cap T(X))) < 2\varepsilon \).

**Proof.** Let \( z \) be the center of \( B \). If \( d(z,T(X)) < \varepsilon \) then \( \text{diam}(B) < 2\varepsilon \) and we are done. Hence, we may assume that \( d(z,T(X)) \geq \varepsilon \). We will show that this leads to a contradiction. By Proposition 4.4.3 and our choice of \( \delta, z \in C^\infty \setminus T(A \cup X) \). The straight line segment \( \ell \) from \( x \) to \( z \) must cross \( T(X) \cup A \) at some point \( w \). Since the segment \( \ell \) is in the interior of the maximal ball \( B \), it is disjoint from \( T(X) \), so \( w \in A \). Hence \( d(x,w) \geq \varepsilon \) and, since \( x \in B \), \( B(w, \varepsilon) \subset B \). This is a contradiction since \( A \subset B(w, \varepsilon) \) and \( \delta < \varepsilon/2 \) so \( \overline{A} \) would be contained in the interior of \( B \) which is impossible since \( A \) is a crosscut of \( T(X) \). \( \square \)

**Proposition 4.4.5.** Let \( C \) be a crosscut of \( T(X) \) and let \( A \) and \( B \) be disjoint closed sets in \( T(X) \) such that \( C \cap A \neq \emptyset \neq C \cap B \). For each \( x \in C \), let \( F_x \in \mathcal{KPP} \) so that \( x \in F_x \). If each \( \overline{F_x} \) intersects \( A \cup B \), then there exists an \( F_\infty \in \mathcal{KPP} \) such that \( \overline{F_\infty} \) intersects \( A, B \) and \( C \).

**Proof.** Let \( a \in A,b \in B \) be the endpoints of \( C \). Let \( C_a,C_b \subset C \) be the set of points \( x \in C \) such that \( \overline{F_x} \) intersects \( A \) or \( B \), respectively. Then \( C_a \) and \( C_b \) are closed subsets by Proposition 4.4.1. Note that \( d(A,B) > 0 \). If \( C_a = \emptyset \), choose \( x_i \in C \) converging to \( a \in A \cap C \). Let \( F_{x_i} = \text{conv}_H (B_i \cap T(X)) \), where \( B_i \in \mathcal{B} \) and assume that \( B_\infty = \lim B_i \). Then by Lemma 4.1.6, \( \overline{F_{x_i}} \cap B \neq \emptyset \) and \( \text{lim} F_{x_i} = \overline{F_\infty} \subset \text{conv}_H (B_\infty \cap K) \). Then \( \overline{F_\infty} \cap B \neq \emptyset \) and \( a \in A \cap C \cap \overline{F_\infty} \). Suppose now \( C_a \neq \emptyset \neq C_b \). Then \( C_a \) and \( C_b \) are closed and, since \( C \) is connected, \( C_a \cap C_b \neq \emptyset \). Let \( y \in C_a \cap C_b \). Then \( \overline{F_y} \cap A \neq \emptyset \neq \overline{F_y} \cap B \) and \( y \in F_y \cap C \). \( \square \)

Proposition 4.4.5 allows us to replace small crosscuts which essentially cross the external ray \( R_t \) with non-trivial principal continuum by small nearby \( \mathcal{KPP} \)-chords which also essentially cross \( R_t \). For if \( C \) is a small crosscut of \( T(X) \) with endpoints \( a \) and \( b \) which crosses the external ray \( R_t \) essentially, let \( A \) and \( B \) be the closures of the sets in \( T(X) \) accessible from \( a \) and \( b \), respectively by small arcs missing \( R_t \). If the \( F_\infty \) of proposition 4.4.5 is a gap \( \text{conv}_H (B \cap T(X)) \), then a \( \mathcal{KPP} \)-chord in its boundary crosses \( R_t \) essentially.

Fix a Riemann map \( \varphi : D^\infty \to U^\infty = C^\infty \setminus T(X) \) with \( \varphi(\infty) = \infty \). Recall that an external ray \( R_t \) is the image of the radial line segment with argument \( 2\pi t \) under the map \( \varphi \).
Proposition 4.4.6. Suppose the external ray $R_t$ lands on $x \in T(X)$, and \( \{g_i\}_{i=1}^\infty \) is a sequence of crosscuts of $T(X)$ converging to $x$ such that there exists a null sequence of arcs $A_i \subset \mathbb{C} \setminus T(X)$ joining $g_i$ to $R_t$. Then for sufficiently large $i$, $\text{var}(f, g_i, T(X)) = 0$.

Proof. Since $f$ is fixed point free on $T(X)$ and $f(x) \in T(X)$, we may choose a small ball $W$ with center $x$ in $\mathbb{C}$ such that $f(W) \cap (W \cup R_t) = \emptyset$. For sufficiently large $i$, $A_i \cup g_i \subset W$. Then for each such $i$ there exists a junction $J_i$ starting from a point in $g_i$, with all of its legs staying in $W$ close to $A_i$ until it reaches $R_t$, and then staying close to $R_t$ to $\infty$. By our choice of $W$, $\text{var}(f, g_i, T(X)) = 0$. \(\square\)

Proposition 4.4.7. Suppose that for an external ray $R_t$ we have $R_t \cap \text{int(conv}_T(T(X))) \neq \emptyset$. Then there exists $x \in R_t$ such that $(T(X), x)$-end of $R_t$ is contained in $\text{conv}_T(T(X))$. In particular there exists a chord $g \in KP$ such that $R_t$ crosses $g$ essentially.

Proof. External rays in $U^\infty$ correspond to geodesic half-lines starting at infinity in the hyperbolic metric on $\mathbb{C}^\infty \setminus T(X)$. Half-planes are conformally equivalent to disks. Therefore, Jørgensen’s lemma applies: the intersection of $R_t$ with a half-plane is connected, so it is a half-line. Since the Euclidean convex hull of $T(X)$ is the intersection of all half-planes containing $T(X)$, the intersection of all half-planes containing $T(X)$, $R_t \cap \text{conv}_T(T(X))$ is connected. \(\square\)

Lemma 4.4.8. Let $\mathcal{E}_t$ be a channel (that is, a prime end such that $\text{Pr}(\mathcal{E}_t)$ is non-degenerate) in $T(X)$. Then for each $x \in \text{Pr}(\mathcal{E}_t)$, for every $\delta > 0$, there is a chain $\{g_i\}_{i=1}^\infty$ of chords defining $\mathcal{E}_t$ selected from $KP_\delta$ such that $d(g_i, C_i) \to 0$ and $R_t$ crosses each $g_i$ essentially. By Proposition 4.4.4, the sequence $g_i$ converges to $\{x\}$. \(\square\)

Lemma 4.4.9. Suppose an external ray $R_t$ lands on $a \in T(X)$ with $\{a\} = \text{Pr}(\mathcal{E}_t) \neq \text{Im}(\mathcal{E}_t)$. Suppose $\{x_i\}_{i=1}^\infty$ is a collection of points in $U^\infty$ with $x_i \to x \in \text{Im}(\mathcal{E}_t) \setminus \{a\}$ and $\phi^{-1}(x_i) \to t$. Then there is a sequence of $KP$-chords $\{g_i\}_{i=1}^\infty$ such that for sufficiently large $i$, $g_i$ separates $x_i$ from $\infty$, $g_i \to a$ and $\phi^{-1}(g_i) \to t$.

Proof. The existence of the chords $g_i$ again follows from the remark following Proposition 4.4.5. It is easy to see that $\lim \phi^{-1}(g_i) \to t$. \(\square\)

4.4.1. Auxiliary Continua. We use $KP$-chords to form Carathéodory loops around the continuum $T(X)$.

Definition 4.4.10. Fix $\delta > 0$. Define the following collections of chords:

\[
KP_\delta^+ = \{g \in KP_\delta \mid \text{var}(f, g, T(X)) \geq 0\}
\]

\[
KP_\delta^- = \{g \in KP_\delta \mid \text{var}(f, g, T(X)) \leq 0\}
\]

\[
KP_\delta = KP_\delta^+ \cup KP_\delta^-
\]

To each collection of chords above, there corresponds an auxiliary continuum defined as follows:

\[
T(X)_\delta = T(X) \cup (\cup KP_\delta)
\]

\[
T(X)_\delta^+ = T(X) \cup (\cup KP_\delta^+)
\]
4. PARTITIONS OF DOMAINS IN THE SPHERE

\[ T(X)_s = T(X \cup (\cup \mathcal{KP}_s^+)) \]

Proposition 4.4.11. Let \( Z \in \{ T(X)_s, T(X)_g^+, T(X)_g^- \} \), and correspondingly \( W \in \{ \mathcal{KP}_s^+, \mathcal{KP}_s^+, \mathcal{KP}_s^- \} \). Then the following hold:

1. \( Z \) is a non-separating plane continuum.
2. \( \partial Z \subset T(X) \cup (\cup W) \).
3. Every accessible point \( y \) in \( \partial Z \) is either a point of \( T(X) \) or a point interior to a chord \( g \in W \).
4. If \( y \in \partial Z \cap g \) with \( g \in W \), then \( y \) is accessible, \( g \subset \partial Z \) and \( \partial Z \) is locally connected at each point of \( g \). Hence, if \( \phi : D^\infty \to C^\infty \setminus Z \) is the Riemann map and \( R_t \) is an external ray landing at \( y \), then \( \phi \) extends continuously to an open interval in \( S^1 \) containing \( t \). Moreover, if \( y \in \partial Z \cap [g \setminus g] \), then \( \phi \) extends continuously over a half open \( J \subset S^1 \) with endpoint \( t \) so that \( \phi(J) \subset g \).

Proof. By Proposition 4.4.1, \( T(X) \cup (\cup W) \) is compact. Moreover, \( T(X) \cup (\cup W) \) is connected since each crosscut \( A \in W \) has endpoints in \( T(X) \). Hence, the topological hull \( T(T(X) \cup (\cup W)) \) is a non-separating plane continuum, establishing (1).

Since \( Z \) is the topological hull of \( T(X) \cup (\cup W) \), no boundary points can be in complementary domains of \( T(X) \cup (\cup W) \). Hence, \( \partial Z \subset T(X) \cup (\cup W) \), establishing (2). Conclusion (3) follows immediately.

Suppose \( y \in \partial Z \cap g \) with \( g \in W \). Then \( \text{Sh}(g) \subset Z \) and there exists \( y_i \in C \setminus Z \) such that \( \lim y_i = y \). We may assume that all the points \( y_i \) are on the “same side” of the arc \( g \) (i.e., \( y_i \in C \setminus \text{Sh}(g) \)). This side of \( g \) is either (1) a limit of \( \mathcal{KP} \)-chords \( g_j \), or (2) there exists a gap \( \text{conv}_H(B \cap X) \) on this side with \( g \) in its boundary. In case (1), \( g \subset \text{Sh}(g_j) \) and, since \( y_i \in C \setminus Z \) for all \( i \), \( g_j \notin W \). Hence each \( g_j \subset C \setminus Z \) for all \( j \). It follows that every point of \( g \) is accessible, \( g \subset \partial Z \) and \( \partial Z \) is locally connected at each point of \( g \). In case (2) there exists a chord \( g' \neq g \) in the boundary of \( \text{conv}_H(B \cap X) \) which separates \( g \) from infinity. Then \( g' \notin W \) and the interior of \( \text{conv}_H(B \cap X) \subset C \setminus Z \). Hence the same conclusion follows.

The last part of (4) follows from the proof of Carathéodory’s theorem (see [Pom92]).

Proposition 4.4.12. \( T(X)_s \) is locally connected; hence, \( \partial T(X)_s \) is a Carathéodory loop.

Proof. Suppose that \( T(X)_s \) is not locally connected. Then \( T(X)_s \) has a non-trivial impression and there exist \( 0 < \varepsilon < \delta/2 \) and a chain \( A_i \) of crosscuts of \( T(X)_s \) such that \( \text{diam}(\text{Sh}(A_i)) > \varepsilon \) for all \( i \). We may assume that \( \lim A_i = y \in T(X)_s \).

By Proposition 4.4.11 (4) we may assume \( y \in X \). Choose \( z_i \in \text{Sh}(A_i) \) such that \( d(z_i, y) > \varepsilon \). We can enlarge the crosscut \( A_i \) of \( T(X)_s \) to a crosscut \( C_i \) of \( T(X) \) as follows. Suppose that \( A_i \) joins the points \( a_i^+ \) and \( a_i^- \) in \( T(X)_s \). If \( a_i^+ \in T(X)_s \), put \( y_i^+ = a_i^+ \). Otherwise \( a_i^+ \) is contained in a chord \( g_i^+ \in \mathcal{KP}_s^+ \), with endpoints in \( T(X)_s \), which is contained in \( T(X)_s \). Since \( \lim A_i = y \), we can select one of these endpoints and call it \( y_i^+ \) such that \( d(y_i^+, a_i^+) \to 0 \). Define \( g_i^- \) and \( y_i^- \) similarly. Then \( g_i^+ \cup A_i \cup g_i^- \) contains a crosscut \( C_i \) of \( T(X) \) joining the points \( y_i^+ \) and \( y_i^- \) such that \( \lim C_i = y \). We claim that \( z_i \in \text{Sh}(C_i) \). To see this note that, since \( z_i \in \text{Sh}(A_i) \), there exists a half-ray \( R_i \subset C \setminus \{ T(X)_s \} \) joining \( z_i \) to infinity such that \( |R_i \cap A_i| \) is an odd number and each intersection is transverse. Since \( R_i \cap C_i = R_i \cap A_i \), it follows...
that $z_i \in \text{Sh}(C_i)$. Let $\text{conv}_H(B_i \cap X)$ be the unique hull of the Kulkarni-Pinkall partition $\mathcal{KPP}$ which contains $z_i$. Since $\text{diam}(C_i) \to 0$ and $d(z_i, y) > \varepsilon$, it follows from Proposition 4.4.4 that $\text{diam}(\text{conv}_H(B_i \cap X)) < 2\varepsilon < \delta$. This contradicts the fact that $z_i \in \mathbb{C} \setminus T(X)_\delta$ and completes the proof. \qed
Part 2

Applications of basic theory
CHAPTER 5

Description of main results of Part 2

We begin by describing the results obtained in Part 2. These results are applications of the tools developed in Part 1. We will say that a continuum $X$ is decomposable if there exist two proper subcontinua $A, B$ of $X$ such that $X = A \cup B$. A continuum which is not decomposable, is called indecomposable.

5.1. Outchannels

In Chapter 6 we will study outchannels. Outchannels were introduced by Bell to establish that a minimal counterexample to the Plane Fixed Point Problem must be an indecomposable continuum. In Chapter 6 we will recover this result and strengthen it by showing that the outchannel in a minimal counterexample to the Plane Fixed Point Problem is unique: there exists exactly one prime end $E_t$ which corresponds to a dense channel with non zero variation. It will follow that the variation of this channel must be $-1$ while all other small crosscuts, which do not cross this channel essentially, must have variation zero. Let us assume that $f : \mathbb{C} \to \mathbb{C}$ with a forward invariant non-separating continuum $X$ presents a (possibly existing) minimal counterexample to the Plane Fixed Point Problem.

We construct a specific locally connected (but not invariant) continuum $X' \supset X$ by adding small crosscuts to $X$. This will be done in a careful way; we will only add Kulkarni-Pinkall crosscuts from $K\mathbb{P}$. This construction is used to show that if there is a minimal counterexample $(X, f)$ to the Plane Fixed Point Problem, then there exists a continuum $Z$ such that the following facts hold.

1. $Z \supset X$;
2. there exists a one-to-one map $\varphi : \mathbb{R} \to Z$,
3. $\varphi(\mathbb{R})$ is the set of accessible points of $Z$,
4. as $t \to \infty$, $\varphi(t)$ and $\varphi(-t)$ run along opposite sides of the outchannel.

Moreover, the same construction is important in the proof of the uniqueness of the outchannel.

These ideas are also applied in [BO09]. There it was shown that in certain cases a minimal subcontinuum $X$ without a fixed point must be fully invariant. As an important tool it was shown in that paper that the map $f$ can be modified on $\mathbb{C} \setminus X$ to a map $g$ such that $g(R_t) = R_t$, $g$ maps points on $R_t$ closer to infinity and $g$ locally interchanges the two sides of $R_t$. Here $R_t$ is the conformal external ray which represents the prime end corresponding to the outchannel. Note that if $X$ is fully invariant then a prime end which corresponds to the outchannel has the property that, in a defining sequence $\{C_i\}$ of crosscuts of the prime end $f(C_i)$ separates $C_i$ from infinity in $U^\infty(X)$ (thus justifying the name “outchannel”).

Suppose that $f : \mathbb{C} \to \mathbb{C}$ is a perfect map, $X$ is a continuum, $f$ has no fixed point in $T(X)$ and $X$ is minimal with respect to $f(X) \subset T(X)$. Fix $\eta > 0$ such that
for each $\mathcal{KP}$-chord $\mathbf{g} \subset T(X)_\eta$, $\mathbf{g} \cap f(\mathbf{g}) = \emptyset$ and $f$ is fixed point free on $T(X)_\eta$. In this case we will say that $\eta$ defines variation near $X$ and that the triple $(f, X, \eta)$ satisfies the standing hypothesis. As usual, for a continuum $X$ let $U^\infty = C^\infty \setminus T(X)$.

**Definition 5.1.1 (Outchannel).** Suppose that the triple $(X, f, \eta)$ satisfies the standing hypothesis. An outchannel of the non-separating plane continuum $T(X)$ is a prime end $E_i$ of $U^\infty$ such that for some chain $\{\mathbf{g}_i\}$ of crosscuts defining $E_i$, $\text{var}(f, \mathbf{g}_i, T(X)) \neq 0$ for every $i$. We call an outchannel $E_i$ of $T(X)$ a geometric outchannel if and only if for sufficiently small $\delta$, every chord in $\mathcal{KP}_\delta$, which crosses $E_i$ essentially, has nonzero variation. We call a geometric outchannel negative (respectively, positive) (starting at $\mathbf{g} \in \mathcal{KP}$) if and only if every $\mathcal{KP}$-chord $\mathbf{h} \subset T(X)_\eta \cap \text{Sh}(\mathbf{g})$, which crosses $E_i$ essentially, has negative (respectively, positive) variation.

### 5.2. Fixed points in invariant continua

In this Section we describe the results obtained in Section 7.1 of Chapter 7. The main result of Section 7.1 solves the Plane Fixed Point Problem in the affirmative for positively oriented maps of the plane. Namely, the following theorem is proven.

**Theorem 7.1.3.** Suppose $f : \mathbb{C} \to \mathbb{C}$ is a positively oriented map and $X$ is a continuum such that $f(X) \subset T(X)$. Then there exists a point $x_0 \in T(X)$ such that $f(x_0) = x_0$.

### 5.3. Fixed points in non-invariant continua – the case of dendrites

As described in Chapter 1, in the rest of Chapter 7 we want to extend Theorem 7.1.3 to at least some non-invariant continua. We are motivated by the interval case in which to conclude that there exists a fixed point in an interval it is enough to know that the endpoints of the interval map in opposite directions, and the invariantness of the interval itself is not crucial.

In Section 7.2 we extend, in the spirit of the interval case, a well-known result according to which a map of a dendrite into itself has a fixed point (Theorem 1.0.2, see [Nad92]). We show the existence of fixed points in non-invariant dendrites and, with some additional conditions, obtain also results related to the number of periodic points of $f$. To state the precise results we need some definitions.

**Definition 5.3.1 (Boundary scrambling for dendrites).** Suppose that $f$ maps a dendrite $D_1$ to a dendrite $D_2 \supset D_1$. Put $E = \overline{D_2 \setminus D_1 \cap D_1}$ (observe that $E$ may be infinite). If for each non-fixed point $e \in E$, $f(e)$ is contained in a component of $D_2 \setminus \{e\}$ which intersects $D_1$, then we say that $f$ has the boundary scrambling property or that it scrambles the boundary. Observe that if $D_1$ is invariant then $f$ automatically scrambles the boundary.

The following theorem is the first result obtained in Section 7.2.

**Theorem 7.2.2.** Let $f : D_1 \to D_2$ be a map, where $D_1$ and $D_2$ are dendrites and $D_1 \subset D_2$. The following claims hold.

(1) If $a, b \in D_1$ are such that $a$ separates $f(a)$ from $b$ and $b$ separates $f(b)$ from $a$, then there exists a fixed point $c \in (a, b)$. Thus, if $e_1 \neq e_2 \in E$ are such that each $f(e_i)$ belongs to a component of $D_2 \setminus \{e_i\}$ disjoint from $D_1$ then there is a fixed point $c \in (e_1, e_2)$. 

(2) If $f$ scrambles the boundary, then $f$ has a fixed point.

To give the next definition we recall that if $x \in Y$ then the valence of $Y$ at $x$, $\text{val}_Y(x)$, is defined as the number of connected components of $Y \setminus \{x\}$, and $x$ is said to be a cutpoint (of $Y$) if $\text{val}_Y(x) > 1$.

**Definition 5.3.2 (Weakly repelling periodic points).** In the situation of Definition 5.3.1, let $a \in D_1$ be a fixed point and suppose that there exists a component $B$ of $D_1 \setminus \{a\}$ such that arbitrarily close to $a$ in $B$ there exist fixed cutpoints of $D_1$ or points $x$ separating $a$ from $f(x)$. Then we say that $a$ is a weakly repelling fixed point (of $f$ in $B$). A periodic point $a \in D_1$ is said to be simply weakly repelling if there exists $n$ and a component $B$ of $D_1 \setminus \{a\}$ such that $a$ is a weakly repelling fixed point of $f^n$ in $B$.

We use the notions introduced in Definition 5.3.2 to prove Theorem 7.2.6.

**Theorem 7.2.6.** Suppose that $f : D \to D$ is continuous where $D$ is a dendrite and all its periodic points are weakly repelling. Then $f$ has infinitely many periodic cutpoints.

This theorem is applied in Theorem 7.2.7 where it is shown that if $g : J \to J$ is a topological polynomial on its dendritic Julia set (e.g., if $f$ is a complex polynomial with a dendritic Julia set) then it has infinitely many periodic cutpoints.

### 5.4. Fixed points in non-invariant continua – the planar case

In parallel with the dendrite case, we want to extend Theorem 7.1.3 to a larger class of maps of the plane and non-invariant continua such that certain “boundary” conditions are satisfied. This is accomplished in Section 7.3.

**Definition 5.4.1.** Suppose that $f : \mathbb{C} \to \mathbb{C}$ is a positively oriented map and $X \subset \mathbb{C}$ is a non-separating continuum. Suppose that there exist $n \geq 0$ disjoint non-separating continua $Z_i$ such that the following properties hold:

1. $f(X) \setminus X \subset \bigcup_i Z_i$;
2. for all $i$, $Z_i \cap X = K_i$ is a non-separating continuum;
3. for all $i$, $f(K_i) \cap [Z_i \setminus K_i] = \emptyset$.

Then the map $f$ is said to scramble the boundary (of $X$). If instead of (3) we have

(3a) for all $i$, either $f(K_i) \subset K_i$, or $f(K_i) \cap Z_i = \emptyset$

then we say that $f$ strongly scrambles the boundary (of $X$); clearly, if $f$ strongly scrambles the boundary of $X$, then it scrambles the boundary of $X$. In either case, the continua $K_i$ are called exit continua (of $X$).

Observe that if in Definition 5.4.1 $n = 0$, then $X$ must be invariant (i.e., $f(X) \subset X$).

**Remark 5.4.2.** Since $Z_i$ and $Z_i \cap X = K_i \neq \emptyset$ are non-separating continua and sets $Z_i$ are pairwise disjoint, then $X \cup \bigcup Z_i$ is a non-separating continuum. Loosely, scrambling the boundary means that $f(X)$ can only “grow” off $X$ within the sets $Z_i$ and through the sets $K_i \subset X$ while any set $K_i$ itself cannot be mapped outside $X$ within $Z_i$, with more specific restrictions upon the dynamics of $K_i$'s in the case of strong scrambling.

The following theorem extends Theorem 7.1.3 onto some non-invariant continua.
5. DESCRIPTION OF MAIN RESULTS OF PART 2

Theorem 7.3.3. In the situation of 5.4.1, if \( f \) is a positively oriented map which strongly scrambles the boundary of \( X \), then \( f \) has a fixed point in \( X \).

We specify the above theorem for positively oriented maps with isolated fixed points as follows. Given a non-separating continuum \( X \subset \mathbb{C} \), a positively oriented map \( f \) and a fixed point \( p \in X \), we define what it means that \( f \) repels outside \( X \) at \( p \) (see Definition 7.4.5; basically, it means that there exists an invariant external ray to \( X \) which lands at \( p \) and along which the points are repelled away from \( p \) by \( f \)). We also need the next definition which is closely related to that of the index of the map on a simple closed curve.

Definition 5.4.3. Suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a positively oriented map with isolated fixed points and \( x \) is a fixed point of \( f \). Then the local index of \( f \) at \( x \), denoted by \( \text{ind}(f,x) \), is defined as \( \text{ind}(f,S) \) where \( S \) is a small simple closed curve around \( x \).

Then we prove the following theorem.

Theorem 7.4.8. Suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a positively oriented map with isolated fixed points, and \( X \subset \mathbb{C} \) is a non-separating continuum or a point. Suppose that the conditions (1)-(3) in 5.4.1 are satisfied. Moreover, suppose that the following conditions hold.

1. For each fixed point \( p \in X \) we have that \( \text{ind}(f,p) = 1 \) and \( f \) repels outside \( X \) at \( p \).
2. The map \( f \) scrambles the boundary of \( X \). Moreover, for each \( i \) either \( f(K_i) \cap Z_i = \emptyset \), or there exists a neighborhood \( U_i \) of \( K_i \) with \( f(U_i \cap X) \subset X \).

Then \( X \) is a point.

5.5. The polynomial case

Theorems 7.2.7 and 7.4.7 apply to polynomials acting on the complex plane. These theorems allow us to obtain corollaries dealing with the existence of periodic points in certain parts of the Julia set of a polynomial and with the degeneracy of certain continua (e.g., impressions). To discuss this we need the following standard notation.

Suppose that \( P : \mathbb{C} \to \mathbb{C} \) is a complex polynomial of degree \( d \). A \( P \)-periodic point \( a \) of period \( n \) is called repelling if \( |(P^n)'(a)| > 1 \), parabolic if \( (P^n)'(a) \) is a root of unity (i.e., for an appropriate \( k \) we will have \( |(P^n)'(a)|^k = 1 \)) and irrational neutral if \( (p^n)'(a) = e^{2\pi i\alpha} \) with \( \alpha \) irrational. The closure of the union of all repelling periodic points of \( P \) is called the Julia set of \( P \) and is denoted by \( J_P \). Then the set \( U^\infty(J_P) = U^\infty \) (i.e., the unbounded component of \( \mathbb{C} \setminus J_P \) is called the basin of attraction of infinity and the set \( K_P = \mathbb{C} \setminus U^\infty = T(J_P) \) is called the “filled-in” Julia set.

Components of \( \mathbb{C} \setminus J_P \) are called Fatou domains. A Fatou domain is said to be attracting (Siegel, respectively) if it contains a periodic point which is attracting (irrational neutral, respectively); an irrational neutral periodic point like that is said to be a Siegel (periodic) point. A bounded periodic Fatou domain is said to be parabolic if it contains no periodic points (in this case all its points converge to the same parabolic periodic orbit which meets the boundary of the domain). Finally,
an irrational neutral periodic point which belongs to \( J_P \) is said to be a Cremer (periodic) point.

The set \( U^\infty \) is foliated by so-called (conformal) external rays \( R_\alpha \) of arguments \( \alpha \in \mathbb{S}^1 \). By [DH85a], if the degree of \( P \) is \( d \) and \( \sigma_d : \mathbb{C} \to \mathbb{C} \) is defined by \( \sigma_d(z) = z^d \), then \( P(R_\alpha) = R_{\sigma_d(\alpha)} \). Denote by \( C_\ast \) the set of all preimages of critical points in \( U^\infty(J_P) \) \( (C_\ast = \emptyset \) if \( J_P \) is connected). If \( J_P \) is connected, each \( \alpha \in \mathbb{S}^1 \) corresponds to a unique external ray and all external rays are smooth and pairwise disjoint. In general, \( R_\alpha \) is smooth and unique if and only if \( R_\alpha \cap C_\ast = \emptyset \). Other external rays are one-sided limits of smooth rays; it follows that they are non-smooth and there are at most countably many of them (in fact, for each \( \alpha \in \mathbb{S}^1 \) there exist at most two external rays \( R_\alpha^\pm \) with argument \( \alpha \) and each is a one sided limit of smooth external rays, see [LP96] for further details).

It is known that two distinct external rays are not homotopic in the complement of \( K_P \) (with the landing point fixed under the homotopies). Given an external ray \( R_\alpha \) of \( K_P \), we denote by \( \Pi(R_\alpha) = \overline{K_\alpha} \setminus R_\alpha \) the principal continuum of \( R_\alpha \). Given a set \( \mathfrak{R} \) of external rays, we extend the above notation by setting \( \Pi(\mathfrak{R}) = \bigcup_{R_\alpha \in \mathfrak{R}} \Pi(R) \).

Now we are ready to give the following technical definition (see Figure 7.4 for an illustration).

**Definition 5.5.1** (General puzzle-piece). Let \( P : \mathbb{C} \to \mathbb{C} \) be a polynomial. Let \( X \subset K_P \) be a non-separating subcontinuum or a point such that the following holds:

1. There exists \( m \geq 0 \) and \( m \) pairwise disjoint non-separating continua/points \( E_1 \subset X, \ldots, E_m \subset X \).
2. There exist \( m \) finite sets of external rays \( A_1 = \{ R_{a_1}, \ldots, R_{a_1} \}, \ldots, A_m = \{ R_{a_m}, \ldots, R_{a_m} \} \) with \( i_k \geq 2, 1 \leq k \leq m \).
3. We have \( \Pi(A_j) \subset E_j \) (so the set \( E_j \cup (\bigcup_{k=1}^{j} R_{a_k}) = E_j^1 \) is closed and connected).
4. \( X \) intersects a unique component \( C_X \) of \( \mathbb{C} \setminus \bigcup E_j \).
5. For each Fatou domain \( U \) either \( U \cap X = \emptyset \) or \( U \subset X \).

We call such \( X \) a general puzzle-piece and call the continua \( E_i \) the exit continua of \( X \). For each \( k \), the set \( E_k^1 \) divides the plane into \( i_k \) open sets which we will call wedges (at \( E_k \)); denote by \( W_k \) the wedge which contains \( X \setminus E_k \) (it is well-defined by (4) above).

Note that if \( m = 0 \), \( \bigcup E_j = \emptyset \) and \( C_X = \mathbb{C} \); so, any non-separating continuum in \( K_P \) with the empty set of exit continua satisfying (5) is a general puzzle-piece. Observe also, that there is a natural situation in which general puzzle-pieces can occur. Suppose that \( J_P \) is connected, conditions (1) - (3) are satisfied, and all continua \( E_j \) are contained in \( J_P \) while the continuum \( X \) is not yet defined. Suppose that there exists a component \( C \) of \( \mathbb{C} \setminus \bigcup E_j \) such that the boundary of \( C \) meets every \( E_j, 1 \leq j \leq m \). Let \( X = (C \cap K_P) \cup (\bigcup E_j) \). Then it is easy to see that \( X \) is a general puzzle-piece. However, our definition allows for a wider variety of general puzzle-pieces (like, e.g., non-separating invariant subcontinua of \( J_P \)).

For convenience call a fixed point \( x \) of a polynomial \( P \) non-rotational if there is a fixed external ray landing at \( x \) (it follows that each such point is either repelling or parabolic). We are ready to state the main result of Section 7.5.
Theorem 7.5.2. Let $P$ be a polynomial with filled-in Julia set $K_P$ and let $Y$ be a non-degenerate periodic component of $K_P$ such that $P^p(Y) = Y$. Suppose that $X \subseteq Y$ is a non-degenerate general puzzle-piece with $m \geq 0$ exit continua $E_1, \ldots, E_m$ such that $P^p(X) \cap C_X \subset X$ and either $P^p(E_i) \subset W_i$, or $E_i$ is a $P^p$-fixed point. Then at least one of the following claims holds:

1. $X$ contains a $P^p$-invariant parabolic domain,
2. $X$ contains a $P^p$-fixed point which is neither repelling nor parabolic, or
3. $X$ has an external ray $R$ landing at a repelling or parabolic $P^p$-fixed point such that $P^p(R) \cap R = \emptyset$ (i.e., $P^p$ locally rotates at some parabolic or repelling $P^p$-fixed point).

Equivalently, suppose that $Y$ is a non-degenerate periodic component of $K_P$ such that $P^p(Y) = Y$, $X \subset Y$ is a general puzzle-piece with $m \geq 0$ exit continua $E_1, \ldots, E_m$ such that $P^p(X) \cap C_X \subset X$ and either $P^p(E_i) \subset W_i$, or $E_i$ is a $P^p$-fixed point; if, moreover, $X$ contains only non-rotational $P^p$-fixed points and does not contain $P^p$-invariant parabolic domains, then it is degenerate.

We also prove in Corollary 7.5.4 that an impression of an invariant external ray, to the filled in Julia set, which contains only repelling or parabolic periodic points is degenerate.
CHAPTER 6

Outchannels and their properties

6.1. Outchannels

In this section we will always let \( f : \mathbb{C} \to \mathbb{C} \) be a continuous function. Suppose that \( X \) is a minimal continuum such that \( f(X) \subset T(X) \) and \( f \) has no fixed point in \( T(X) \). We show that \( X \) has at least one negative outchannel. We will always assume that \( (f, X, \eta) \) satisfies the standing hypothesis (see Definition 5.1.1 and the paragraph preceding 5.1.1) and see Section 4.4.1 for the notation \( T(X)_\delta^- \). In particular, \( f \) is fixed point free on \( T(X)_\eta \). Note that for each \( KP \)-chord \( g \) in \( T(X)_\eta \), \( \text{var}(f, g, T(X)) = \text{var}(f, g) \) is defined.

**Lemma 6.1.1.** Suppose that \( (f, X, \eta) \) satisfy the standing hypothesis and \( \delta \leq \eta \). Let \( Z \in \{ T(X)_\delta^+, T(X)_\delta^- \} \). Fix a Riemann map \( \varphi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus Z \) such that \( \varphi(\infty) = \infty \). Suppose \( R_t \) lands at \( x \in \partial Z \). Then there is an open interval \( M \subset \partial \mathbb{D}^\infty \) containing \( t \) such that \( \varphi \) can be extended continuously over \( M \).

**Proof.** Suppose that \( Z = T(X)_\delta^+ \) and \( R_t \) lands on \( x \in \partial Z \). By proposition 4.4.11 we may assume that \( x \in X \). Note first that the family of chords in \( KP_\delta \) form a closed subset of the hyperspace of \( \mathbb{C} \setminus X \), by Proposition 4.4.1. By symmetry, it suffices to show that we can extend \( \psi \) over an interval \( [t', t] \subset S^1 \) for \( t' < t \).

Let \( \phi : \mathbb{D}^\infty \to \mathbb{C} \setminus T(X) \) be the Riemann map for \( T(X) \). Then there exists \( s \in S^1 \) so that the external ray \( R_s \) of \( \mathbb{C} \setminus T(X) \) lands at \( x \). Suppose first that there exists a chord \( g \in KP_\delta \) such that \( G = \varphi^{-1}(g) \) has endpoints \( s' \) and \( s \) with \( s' < s \). Since \( KP_\delta \) is closed, there exists a minimal \( s'' \leq s' < s \) such that there exists a chord \( h \in KP_\delta \) so that \( H = \varphi^{-1}(h) \) has endpoints \( s'' \) and \( s \). Then \( h \subset \partial Z \) and \( \phi \) can be extended over an interval \( [t', t] \) for some \( t' < t \), by Proposition 4.4.11 (4).

Suppose next that no such chord \( g \) exists. Choose a junction \( J_x \) for \( T(X)_\delta^- \) and a neighborhood \( W \) of \( x \) such that \( f(W) \cap [W \cup J_x] = \emptyset \). We will first show that there exists \( \nu \leq \delta \) such that \( x \in \partial T(X)_\nu \). For suppose that this is not the case. Then there exists a sequence \( g_i \in KP \) of chords such that \( x \in \text{Sh}(g_{i+1}) \subset \text{Sh}(g_i) \), \( \lim g_i = x \) and \( \text{var}(f, g_i) > 0 \) for all \( i \). This contradicts Proposition 4.4.6. Hence \( x \in \partial T(X)_\nu \) for some \( \nu > 0 \). We may assume that \( \nu \) is so small that any chord of \( KP_\nu \) with endpoint \( x \) is contained in \( W \).

By Proposition 4.4.12, the boundary of \( T(X)_\nu \) is a simple closed curve \( S \) which must contain \( x \). If there exists a chord \( h \in KP_\nu \) with endpoint \( x \) such that \( H \) has endpoints \( s' \) and \( s \) with \( s' < s \), then, since \( h \subset W \), \( f(h) \cap J_x = \emptyset \), \( \text{var}(f, h) = 0 \) and \( h \in KP_\delta^- \), a contradiction. Similarly, all chords \( h \) close to \( x \) in \( S \) so that \( H \) has endpoints less than \( s \) and which are contained in \( W \) have \( \text{var}(f, h) = 0 \) by Proposition 4.4.6. Hence a small interval \( [x', x] \subset S \), in the counterclockwise order on \( S \) is contained in \( T(X)^-_\delta \). It now follows easily that a similar arc exists in the boundary of \( T(X)^-_\delta \) and the desired result follows. \( \square \)
6. OUTCHANNELS AND THEIR PROPERTIES

By a narrow strip we mean the image of an embedding \( h : \{(x, y) \in \mathbb{C} \mid x \geq 0 \text{ and } -1 < y < 1\} \to \mathbb{C} \) such that \( h \) has a continuous extension over the closure of its domain and \( \lim_{x \to \infty} \text{diam}(h(\{x\} \times [-1, 1])) = 0. \)

**Lemma 6.1.2.** Suppose that \((f, X, \eta)\) satisfy the standing hypothesis. If there is a chord \( g \subset T(X) \) of \( T(X) \) of negative (respectively, positive) variation, such that there is no fixed point in \( T(T(X) \cup g) \), then there is a negative (respectively, positive) geometric outchannel \( E_t \) of \( T(X) \) starting at \( g \).

Moreover, if \( E_t \) is a positive (negative) geometric outchannel starting at the KP-chord \( g \), and \( \mathcal{G} = \bigcup\{\text{conv}_H(B \cap T(X)) \mid \text{conv}_H(B \cap T(X)) \subset T(X)_\eta \cap \text{Sh}(g) \) and a chord in \( \text{conv}_H(B \cap T(X)) \) crosses \( R_t \) essentially\{. Then \( \mathcal{G} \) is an infinite narrow strip in the plane whose remainder is contained in \( T(X) \) and which is bordered by a KP-chord and two halflines \( H_1 \) and \( H_2 \) (see figure 6.1).

**Proof.** Without loss of generality, assume \( \text{var}(f, g, T(X)) = \text{var}(f, g) < 0 \). If \( g \) is such that for any chord \( h \subset T(X \cup g) \), \( h \subset T(X)_\eta \), put \( g' = g \). Otherwise consider the boundary of \( T(X)_{\delta} \) \( (\delta < \eta) \) which is locally connected by Proposition 4.4.12 and, hence, a Carathéodory loop. Then a continuous extension \( g : S^1 \to \partial T(X)_{\delta} \) of the Riemann map \( \phi : \mathbb{D}^\infty \to \mathbb{C}^\infty \setminus T(X)_{\delta} \) exists. Whence the boundary of \( T(X)_{\delta} \) contains a sub-path \( A = g([a, b]) \), which is contained in \( \text{Sh}(g) \), whose endpoints coincide with the endpoints of \( g \). Note that for each component \( C \) of \( A \setminus X \), \( \text{var}(f, C) \) is defined. Then it follows from Proposition 3.4.4, applied to a Carathéodory path, that there exists a component \( C = g' \) such that \( \text{var}(f, g') < 0 \).
Note that $g'$ is a $\mathcal{KP}$-chord contained in the boundary of $T(X)_\delta$.
By taking $\delta$ sufficiently small we can assume that for any chord $h \subset \overline{\text{Sh}(g')}$, $h \subset T(X)_{\eta}$.

To see that a geometric outchannel, starting with $g'$, exists, note that for any chord $g'' \subset \overline{\text{Sh}(g')}$ with $\text{var}(f, g'', X) < 0$, if $g'' = \lim g_i$, then there exists $i$ such that for any chord $h$ which separates $g_i$ and $g''$ in $U^\infty(X)$, $\text{var}(f, h, X) = \text{var}(f, g'', X) < 0$. This follows since $f'(g)$ is close to $f(h)$ and, hence, crosses a junction $J_\nu$ in the same way (we can slightly change the junction $J_\nu$ with $v \in g''$ to a junction with vertex in $h$ without changing the crossings of the images of the crosscuts with the junction). If $g''$ is isolated on the side closest to $X$, then $g'' \subset \text{conv}_H(B \cap T(X))$, where $\text{conv}_H(B \cap T(X))$ is a gap, such that $g''$ separates $\text{conv}_H(B \cap T(X)) \setminus g''$ from infinity in $U^\infty(X)$. Again by Proposition 3.4.4, there exists $h \neq g''$ in $\text{conv}_H(B \cap T(X))$ such that $g''$ separates $h$ from infinity in $U^\infty(X)$ and $\text{var}(f, h, X) < 0$. It follows from these two facts that there exists a maximal family of $\mathcal{KP}$ crosscuts, all of which have negative variation and are such that for any three members of the family, one separates the other two in $U^\infty(X)$. Hence this maximal family determines a geometric outchannel. Each chord $h$ in this family corresponds to a unique maximal ball $B_h$. It is now not difficult to see that the union of all the sets $\text{conv}_H(B_h \cap T(X))$ is a narrow strip.

6.1.1. Invariant Channel in $X$. We are now in a position to prove Bell’s principal result on any possible counter-example to the fixed point property, under our standing hypothesis.

**Lemma 6.1.3.** Suppose $E_t$ is a geometric outchannel of $T(X)$ under $f$. Then the principal continuum $\Pr(E_t)$ of $E_t$ is invariant under $f$. So $\Pr(E_t) = X$.

**Proof.** Let $x \in \Pr(E_t)$. Then for some chain $\{g_i\}_{i=1}^\infty$ of crosscuts defining $E_t$ selected from $\mathcal{KP}_3$, we may suppose $g_i \to x \in \partial T(X)$ (by Lemma 4.4.8) and $\text{var}(f, g_i, X) \neq 0$ for each $i$. The external ray $R_t$ meets all $g_i$ and there is, for each $i$, a junction from $g_i$ which “parallels” $R_t$. Since $\text{var}(f, g_i, X) \neq 0$, each $f(g_i)$ intersects $R_t$. Since $\text{diam}(f(g_i)) \to 0$, we have $f(g_i) \to f(x)$ and $f(x) \in \Pr(E_t)$. We conclude that $\Pr(E_t)$ is invariant.\hfill $\Box$

**Theorem 6.1.4 (Dense channel, Bell).** If $(X, f, \eta)$ satisfy our standing hypothesis then $T(X)$ contains a negative geometric outchannel; hence, $\partial U^\infty = \partial T(X) = X = f(X)$ is an indecomposable continuum.

**Proof.** By Lemma 4.4.12 $\partial T(X)_\eta$ is a Carathéodory loop. Since $f$ is fixed point free on $T(X)_\eta$, $\text{ind}(f, \partial T(X)_\eta) = 0$. Consequently, by Theorem 3.2.2 for Carathéodory loops, $\text{var}(f, \partial T(X)_\eta) = -1$. By the summability of variation on $\partial T(X)_\eta$, it follows that on some chord $g \subset \partial T(X)_\eta$, $\text{var}(f, g, T(X)) < 0$. By Lemma 6.1.2, there is a negative geometric outchannel $E_t$ starting at $g$. Since $\Pr(E_t)$ is invariant under $f$ by Lemma 6.1.3, it follows that $\Pr(E_t)$ is an invariant subcontinuum of $\partial U^\infty \subset \partial T(X) \subset X$. So by the minimality condition in our Standing Hypothesis, $\Pr(E_t)$ is dense in $X$. It then follows from a theorem of Rutt [Rut35] that $X$ is an indecomposable continuum.\hfill $\Box$

**Theorem 6.1.5.** Assume that $(X, f, \eta)$ satisfy our standing hypothesis and $\delta \leq \eta$. Then the boundary of $T(X)_\delta$ is a simple closed curve. The set of accessible points in the boundary of each of $T(X)^+\delta$ and $T(X)^-\delta$ is an at most countable union of pairwise disjoint continuous one-to-one images of $\mathbb{R}$. 

6.1. OUTCHANNELS
6. Outhannels and Their Properties

Proof. By Theorem 6.1.4, \( X \) is indecomposable, so it has no cut points. By Proposition 4.4.12, \( \partial T(X)_\delta \) is a Carathéodory loop. Since \( X \) has no cut points, neither does \( T(X)_\delta \). A Carathéodory loop without cut points is a simple closed curve.

Let \( Z \in \{ T(X)_\delta^+, T(X)_\delta^- \} \) with \( \delta \leq \eta \). Fix a Riemann map \( \phi : \mathbb{D}^\infty \rightarrow \mathbb{C}^\infty \setminus Z \) such that \( \phi(\infty) = \infty \). Corresponding to the choice of \( Z \), let \( \mathcal{W} \in \{ \mathcal{KP}^+\delta, \mathcal{KP}^-\delta \} \). Apply Lemma 6.1.1 and find the maximal collection \( \mathcal{J} \) of disjoint open subarcs of \( \partial \mathbb{D}^\infty \) over which \( \phi \) can be extended continuously. The collection \( \mathcal{J} \) is countable. Since \( X \) has no cutpoints the extension is one-to-one over \( \cup \mathcal{J} \). Since angles that correspond to accessible points are dense in \( \partial \mathbb{D}^\infty \), so is \( \cup \mathcal{J} \). If \( Z = T(X)_\delta^+ \), then it is possible that \( \cup \mathcal{J} \) is all of \( \partial \mathbb{D}^\infty \) except one point, but it cannot be all of \( \partial \mathbb{D}^\infty \) since there is at least one negative geometric outhannel by Theorem 6.1.4.

Theorem 6.1.5 still leaves open the possibility that \( Z \in \{ T(X)_\delta^+, T(X)_\delta^- \} \) has a very complicated boundary. The set \( C = \partial \mathbb{D}^\infty \setminus \cup \mathcal{J} \) is compact and zero-dimensional. Note that \( \phi \) is discontinuous at points in \( C \). We may call \( C \) the set of outhannels of \( Z \). In principle, there could be an uncountable set of outhannels, each dense in \( X \). The one-to-one continuous images of half lines in \( \mathbb{R} \) lying in \( \partial Z \) are the “sides” of the outhannels. If two elements \( J_1 \) and \( J_2 \) of the collection \( \mathcal{J} \) happen to share a common endpoint \( t \), then the prime end \( \mathcal{E}_t \) is an outhannel in \( Z \), dense in \( X \), with images of half lines \( \phi(J_1) \) and \( \phi(J_2) \) as its sides. It seems possible that an endpoint \( t \) of \( J \in \mathcal{J} \) might have a sequence of elements \( J_t \) from \( \mathcal{J} \) converging to it. Then the outhannel \( \mathcal{E}_t \) would have only one (continuous) “side.” Such exotic possibilities are eliminated in the next section.

In the proposition below we summarize several of the results in this section and show that an arc component \( K \) of the set of accessible points of the boundary of \( T(X)_\delta^- \) is efficient in connecting close points in \( K \). Note that it will follow later from Theorem 6.2.1 that there are no chords of positive variation. Hence \( T(X)_\delta^- = T(X)_\delta \) which is always a simple closed curve.

Proposition 6.1.6. Suppose that \( (X, \mathcal{J}, \eta) \) satisfy our standing hypothesis, that the boundary of \( T(X)_\delta^- \) is not a simple closed curve, \( \delta \leq \eta \) and that \( K \) is an arc component of the boundary of \( T(X)_\delta^- \) so that \( K \) contains an accessible point. Let \( \varphi : \mathbb{D}^\infty \rightarrow \mathbb{C}^\infty \setminus T(X)_\delta^- \) be a conformal map such that \( \varphi(\infty) = \infty \). Then:

1. \( \varphi \) extends continuously and injectively to a map \( \tilde{\varphi} : \mathbb{D}^\infty \rightarrow \tilde{\mathbb{U}}^\infty \), where \( \mathbb{D}^\infty \setminus \tilde{\mathbb{U}}^\infty \) is a dense and open subset of \( \mathbb{S}^1 \) which contains \( K \) in its image. Let \( \tilde{\varphi}^{-1}(K) = (t', t) \subset \mathbb{S}^1 \) with \( t' < t \) is the counterclockwise order on \( \mathbb{S}^1 \).
   Hence \( \tilde{\varphi} \) induces an order \( < \) on \( K \). If \( x < y \in K \), we denote by \( (x, y) \) the subarc of \( K \) from \( x \) to \( y \) and by \( (x, \infty) = \cup_{y > x} (x, y) \).
2. \( \mathcal{E}_t \) and \( \mathcal{E}_\mathcal{J} \) are positive geometric outhannels of \( T(X) \).
3. Let \( R_t \) be the external ray of \( T(X)_\delta^- \) with argument \( t \). There exists \( s \in R_t \), \( B \in \mathcal{B}^\infty \) and \( \mathfrak{g} \in \mathcal{KP} \) such that \( s \in \mathfrak{g} \subset \text{conv}_\mathcal{H}(B \cap X) \) and \( s \) is the last point of \( R_t \) in \( \text{conv}_\mathcal{H}(B \cap X) \) (from \( \infty \)), \( \mathfrak{g} \) crosses \( R_t \) essentially and for each \( B' \in \mathcal{B}^\infty \) with \( \text{conv}_\mathcal{H}(B' \cap X) \setminus X \subset \text{Sh}(\mathfrak{g}) \), \( \text{diam}(B') < \delta \).
4. There exists \( x' \in K \) such that if \( B' \in \mathcal{B}^\infty \) with \( \text{int}(B') \subset \text{Sh}(\mathfrak{g}) \), then \( \text{conv}_\mathcal{H}(B' \cap X) \cap (x', \infty) \) is a compact ordered subset of \( K \) so that if \( C \) is \( \mathcal{KP} \)-crosscut in the boundary of \( \text{conv}_\mathcal{H}(B' \cap X) \) with both endpoint in \( K \), then \( C \subset K \).
(5) Let $\mathfrak{B}^\infty \subset \mathfrak{B}^\infty$ be the collection of all $B \in \mathfrak{B}^\infty$ such that $R_t$ crosses a chord in the boundary of $\text{conv}_H(B \cap X)$ essentially and $\text{int}(B) \subset \text{Sh}(g)$.

Then $\mathfrak{S} = \bigcup_{B \in \mathfrak{B}^\infty} \text{conv}_H(B \cap X)$ is a narrow strip in the plane, bordered by two halflines $H_1$ and $H_2$, which compactify on $X$ and one of $H_1$ or $H_2$ contains the set $\langle \hat{x}', \infty \rangle$ for some $\hat{x}' \in K$.

In particular, if $\max(\hat{x}, \hat{x}') < p < q$ and $\text{diam}(\langle p, q \rangle) > 2\delta$, then there exists a chord $g \in \mathcal{KP}$ such that one endpoint of $g$ is in $\langle p, q \rangle$ and $g$ crosses $R_t$ essentially.

An analogous conclusion holds for $T(X)_J$ since its boundary cannot be a simple closed curve (clearly $T(X)_J$ must contain a crosscut of negative variation).

**Proof.** By Proposition 4.4.11 and Theorem 6.1.5, and its proof, $\varphi$ extends continuously and injectively to a map $\hat{\varphi} : \hat{D}^\infty \to U^\infty$ and (1) holds.

By Lemma 6.1.1, the external ray $R_t$ does not land. Hence there exist a chain $g_i$ of $\mathcal{KP}_\delta$ chords which define the prime end $\mathcal{E}_t$. If for any $i \text{ var}(f, g_i) \leq 0$, then $g_i \subset T(X)_J$ a contradiction with the definition of $t$. Hence $\text{var}(f, g_i) > 0$ for all $i$ sufficiently small and $\mathcal{E}_t$ is a positive geometric outchannel by the proof of Lemma 6.1.2. Hence (2) holds.

The proof of (3) is straightforward and is left to the reader.

Suppose that the endpoints of $g$ are $e$ and $f$ with $f \in K$. Choose $\hat{x} > f$ in $K$ so that $\hat{x}$ is the endpoint of a $\mathcal{KP}$ crosscut which is contained in $\text{conv}_H(B \cap X)$ with $B \subset \text{Sh}(g)$. Let $B' \in \mathfrak{B}^\infty$ with $\text{int}(B') \subset \text{Sh}(g)$, $\hat{x} \not\in B'$ and $\langle \hat{x}, \infty \rangle \setminus \text{conv}_H(B' \cap X)$ not connected. Suppose $(a, b)$ is a bounded component of $\langle \hat{x}, \infty \rangle \setminus \text{conv}_H(B' \cap X)$ with endpoints in $B'$. Note that there must exist a chord $h \in \mathcal{KP}$ with endpoints $a$ and $b$. If $\text{var}(f, h) \leq 0$ we are done. Hence $\text{var}(f, h) > 0$. By Lemma 6.1.2, there is a geometric outchannel $\mathcal{E}_x$ starting at $h$. This outchannel disconnects the arc $\langle a, b \rangle$ between $a$ and $b$, a contradiction. Hence (4) holds.

Next choose $\hat{x}' \in K$ such that each point of $\langle \hat{x}', \infty \rangle$ is accessible from $\text{Sh}(g)$. Then each subarc $\langle p, q \rangle$ of $\langle \hat{x}', \infty \rangle$ of diameter bigger than $2\delta$ cannot be contained in a single element of the $\mathcal{KP}$ partition. Hence there exists a $\mathcal{KP}$-chord $g$ which crosses $R_t$ essentially and has one endpoint in $\langle p, q \rangle$.

Note that for each chord $h \subset \text{Sh}(g)$ which crosses $R_t$ essentially, $\text{var}(f, h) > 0$. By Lemma 6.1.2, $\bigcup_{B \in \mathfrak{B}^\infty} \text{conv}_H(B \cap X)$ is a strip in the plane, bordered by two halflines $H_1, H_2$, which compactify on $X$. These two halflines, consist of chords in $\mathcal{KP}_\delta$ and points in $X$, one of which, say $H_1$ meets $\langle \hat{x}', \infty \rangle$. If $\langle \hat{x}', \infty \rangle$ is not contained in $H_1$ then, as in the proof of (4), there exists a chord $h \subset H_1$ with $\text{var}(f, h) > 0$ joining two points of $x, y \in \langle \hat{x}, \infty \rangle$. As above this leads to a contradiction and the proof is complete. \qed

### 6.2. Uniqueness of the Outchannel

Theorem 6.1.4 asserts the existence of at least one negative geometric outchannel which is dense in $X$. We show below that there is exactly one geometric outchannel, and that its variation is $-1$. Of course, $X$ could have other dense channels, but they are “neutral” as far as variation is concerned.

**Theorem 6.2.1 (Unique Outchannel).** If $(X, f, \eta)$ satisfy the standing hypothesis then there exists a unique geometric outchannel $\mathcal{E}_t$ for $X$, which is dense in $X = \partial T(X)$. Moreover, for any sufficiently small chord $g$ in any chain defining $\mathcal{E}_t$,
60 6. OUTCHANNELS AND THEIR PROPERTIES

\[ \text{Figure 6.2. Uniqueness of the negative outchannel.} \]

\[ \text{var}(f, g, X) = -1, \text{ and for any sufficiently small chord } g' \text{ not crossing } R_t \text{ essentially, } \text{var}(f, g', X) = 0. \]

**Proof.** Suppose by way of contradiction that \( X \) has a positive outchannel. Let \( 0 < \delta \leq \eta \) such that if \( M \subset T(B(T(X), 2\delta)) \) with \( \text{diam}(M) < 2\delta \), then \( f(M) \cap M = \emptyset \). Since \( X \) has a positive outchannel, \( \partial T(X)^{-}_{\delta} \) is not a simple closed curve. By Theorem 6.1.5 \( \partial T(X)^{-}_{\delta} \) contains an arc component \( K \) which is the one-to-one continuous image of \( \mathbb{R} \). Note that each point of \( K \) is accessible.

Let \( \varphi : D^\infty \to V^\infty = \mathbb{C} \setminus T(X)^{-}_{\delta} \) a conformal map. By Proposition 6.1.6, \( \varphi \) extends continuously and injectively to a map \( \tilde{\varphi} : \tilde{D}^\infty \to \tilde{V}^\infty \), where \( \tilde{D}^\infty \setminus D^\infty \) is a dense and open subset of \( S^1 \) which contains \( K \) in its image. Then \( \tilde{\varphi}^{-1}(K) = \)
There exists \( \langle g \rangle \) of \( K \). A variation is defined on each component of \( a \) with diameter \( d \), such that there is a chord \( x \in \partial W \) of \( \langle X \rangle \). Let \( p,d \) be the endpoint of \( t \) essentially at \( \langle X \rangle \). If not, then there exists a \( \mathcal{K}P \)-chord \( g \) which contains a point of \( \langle a \rangle \) and an arc \( h \subset K \) such that \( a \in h \). Let \( p \) be the endpoint of \( h \) such that \( p < a \). See Figure 6.2.

Since \( X \in \langle x,\infty \rangle \) there are components of \( \langle b,\infty \rangle \cap W \) which are arbitrarily close to \( a_0 \). Choose \( b < c < d \) in \( K \) so that the \( \langle c,d \rangle \) is the closure of a component of \( W \cap \langle b,\infty \rangle \) such that:

1. \( a \) and \( d \) lie in the same component of \( \partial W \setminus \{b,c\} \).
2. There exists \( z \in \langle c,d \rangle \cap \partial W \cap W \) and an arc \( I \subset \{a_0,z\} \cup [W \setminus \{p,d\}] \) joining \( a_0 \) to \( z \).
3. There is a \( \mathcal{K}P \)-chord \( g \subset W \) with \( z \) and \( y \) as endpoints which crosses \( R_t \) essentially. Hence, \( \text{var}(f,g) > 0 \).
4. \( \text{diam}(f(g)) < d(J_{a_0}^+ \setminus W,J_{a_0}^+ \setminus W) \).

Conditions (1) and (2) follow because \( J_{a_0} \) is a connected and closed set from \( a_0 \) to \( \infty \) in \( \{a_0\} \cup [\mathbb{C} \setminus T(X)] \) and the ray \( \langle b,\infty \rangle \) approaches both \( a_0 \) and \( p \). Conditions (3) and (4) follow from Proposition 6.1.6. If \( d \in X \), put \( q = d \). Otherwise, let \( q \in \langle a \rangle \). Since \( \langle a \rangle \) approaches \( b \) and \( \partial W \), the \( \langle a \rangle \) approaches \( b \) and \( \partial W \). The arc \( A' \) must enter \( W \) through \( C \). Since we want to apply the Lollipop lemma, we will modify the arc \( A' \) to a new arc \( A \) which is disjoint from \( I \).

Let \( A \) be the set of points in \( A' \cup C \) accessible from \( \infty \) in \( \mathbb{C} \setminus \{S' \cup C\} \). Then \( A \) is a bumping arc from \( p \) to \( q \). \( A \cap I = \emptyset \), \( \text{var}(f,A) \) is defined, \( S = A \cup \{p,q\} \) is a simple closed curve with \( T(S) \subset T(S) \) and \( f \) is fixed point free on \( T(S) \). Note
that \( y \in A \). Then the Lollipop lemma applies to \( S \) with \( R = T((a_0, z) \cup I) \) and \( L = T(I \cup (z, q) \cup A \cup (p, a_0)) \).

Claim: \( f(z) \in R \). Hence by Corollary 3.3.2, \( \langle a_0, z \rangle \) contains a chord \( g_1 \) with \( \text{var}(f, g_1) < 0 \).

Proof of Claim. Note that the positive direction along \( g \) is from \( z \) to \( y \). Since \( z, y \in X \), \( \{f(z), f(y)\} \subset X \subset T(S) = R \cup L \). Choose a junction \( J_z \) such that \( J_{a_0} \setminus W \subset J_z \) and \( J_z \) runs close to \( \langle a_0, z \rangle \) on its way to \( g \). In particular we may assume that \( J_z \cap R = \{z\} \). Since \( g \) crosses \( R_t \) essentially, \( \text{var}(f, g) > 0 \). For \( * \in \{-, i, +\} \), let \( C_z^* \) be the union of components of \( J_z^* \setminus W \) which are disjoint from \( J_{a_0}^* \). Then \( C_z^* \) separates \( R \cup C_z^+ \) from \( L \cup C_z^- \) in \( C \setminus W \) (see figure 6.2). Since \( f(g) \cap J_{a_0} = \emptyset \), if \( f(z) \notin R \), \( \text{var}(f, g) \leq 0 \), a contradiction. Hence \( f(z) \in R \) (and, in fact, \( f(y) \in L \)) as desired.

Since \( f(z) \in R \), \( \langle a_0, z \rangle \) contains a chord \( g_1 \) with \( \text{var}(f, g_1) < 0 \). Repeating the same argument, replacing \( a_0 \) by \( z \) we obtain a second chord \( g_2 \) contained in \( \langle z, \infty \rangle \) such that \( \text{var}(f, g_2) < 0 \).

We will now show that the existence of two distinct chords \( g_1 \) and \( g_2 \) in \( K \) with variation \( < 0 \) on each leads to a contradiction. Recall that \( a_0 \in \langle b, \infty \rangle \). Hence we can find \( y' \in \langle b, \infty \rangle \) with \( y' \in X \) such that \( g_1 \cup g_2 \subset \langle a_0, y' \rangle \) and there exists a small arc \( I' \subset W \) such that \( I' \setminus \langle a_0, y' \rangle = \{a_0, y'\} \). Since \( f(I') \cap J_{a_0} = \emptyset \), \( \text{var}(f, I') = 0 \). We may also assume that \( f \) is fixed point free on \( T(S'') \), where \( S'' = I' \cup \langle a_0, y' \rangle \). Since \( \langle a_0, y' \rangle \) contains both \( g_1, g_2 \) and no chords of positive variation, \( \text{var}(f, \langle a_0, y' \rangle) \leq -2 \) and \( \text{var}(f, S'') \leq -2 \). Then \( \text{ind}(f, S'') = \text{var}(f, S'') + 1 \leq -1 \) a contradiction with Theorem 3.1.4. Hence \( X \) has no positive geometric outchannel.

By Theorems 6.1.4 and 3.2.2, \( X \) has exactly one negative outchannel and its variation is \(-1 \). \( \square \)

Note that the following Theorem follows from Lemma 6.1.6 and Theorem 6.2.1.

**Theorem 6.2.2.** Suppose that \( X \) is a minimal counterexample to the Plane Fixed Point Problem. Then there exists \( \delta > 0 \) such that the continuum \( Y = T(X)^2 \) is a non-separating continuum, \( f \) is fixed point free on \( Y \) and all accessible points of \( Y \) are contained in one arc component \( K \) of the boundary of \( Y \). In other words, \( Y \) is homeomorphic to a disk with exactly one channel removed which corresponds to the unique geometric outchannel of variation \(-1 \) of \( X \). This channel compactifies on \( X \). The sides of this channel are halflines consisting entirely of chords of zero variation and points in \( X \). There exist arbitrarily small homeomorphisms of tails of these halflines to a tail of \( R_t \) which is the external ray corresponding to this channel.
CHAPTER 7

Fixed points

In this chapter we study fixed points in invariant and non-invariant continua under positively oriented maps. We also obtain corollaries dealing with complex polynomials (the applications of these corollaries to complex dynamics are described in Chapter 1.)

7.1. Fixed points in invariant continua

In this section we will consider a positively oriented map of the plane. As we shall see below, a straightforward application of the tools developed above will give us the desired fixed point result. We will often assume, by way of contradiction, that \( f : \mathbb{C} \to \mathbb{C} \) is a positively oriented map, \( X \) is a plane continuum such that \( f(X) \subset T(X) \) and \( T(X) \) contains no fixed points of \( f \).

**Lemma 7.1.1.** Let \( f : \mathbb{C} \to \mathbb{C} \) be a map and \( X \subset \mathbb{C} \) a continuum such that \( f(X) \subset T(X) \). Suppose \( C = (a, b) \) is a crosscut of the continuum \( T(X) \). Let \( v \in (a, b) \) be a point and \( J_v \) be a junction such that \( J_v \cap (X \cup C) = \{v\} \). Then there exists an arc \( I \) such that \( S = I \cup C \) is a simple closed curve, \( T(X) \subset T(S) \) and \( f(I) \cap J_v = \emptyset \).

**Proof.** Since \( f(X) \subset T(X) \) and \( J_v \cap X = 0 \), it is clear that there exists an arc \( I \) with endpoints \( a \) and \( b \) sufficiently close to \( T(X) \) such that \( I \cup C \) is a simple closed curve, \( T(X) \subset T(I \cup C) \) and \( f(I) \cap J_v = \emptyset \). This completes the proof.

**Corollary 7.1.2.** Suppose \( X \subset \mathbb{C} \) is a continuum, \( f : \mathbb{C} \to \mathbb{C} \) a positively oriented map such that \( f(X) \subset T(X) \). Then for each crosscut \( C \) of \( T(X) \) such that \( f(\overline{C}) \cap \overline{C} = \emptyset \), \( \text{var}(f, C) \geq 0 \)

**Proof.** Suppose that \( C = (a, b) \) is a crosscut of \( T(X) \) such that \( f(\overline{C}) \cap \overline{C} = \emptyset \) and \( \text{var}(f, C) \neq 0 \). Choose a junction \( J_v \) such that \( J_v \cap (X \cup C) = \{v\} \) and \( v \in C \setminus X \). By Lemma 7.1.1, there exists an arc \( I \) such that \( S = I \cup C \) is a simple closed curve and \( f(I) \cap J_v = \emptyset \). Moreover, by choosing \( I \) sufficiently close to \( X \), we may assume that \( v \in \mathbb{C} \setminus f(S) \). Hence \( \text{var}(f, C) = \text{Win}(f, S, v) \geq 0 \) by the remark following Definition 2.2.2.

**Theorem 7.1.3.** Suppose \( f : \mathbb{C} \to \mathbb{C} \) is a positively oriented map and \( X \) is a continuum such that \( f(X) \subset T(X) \). Then there exists a point \( x_0 \in T(X) \) such that \( f(x_0) = x_0 \).

**Proof.** Suppose we are given a continuum \( X \) and \( f : \mathbb{C} \to \mathbb{C} \) a positively oriented map such that \( f(X) \subset T(X) \). Assume that \( f|_{T(X)} \) is fixed point free. Choose a simple closed curve \( S \) such that \( X \subset T(S) \) and points \( a_0 < a_1 < \ldots < a_n \) in \( S \cap X \) such that for each \( i \), \( C_i = (a_i, a_{i+1}) \) is a sufficiently small crosscut of \( X \), \( f(\overline{C_i}) \cap \overline{C_i} = \emptyset \) and \( f|_{T(S)} \) is fixed point free. By Corollary 7.1.2, \( \text{var}(f, C_i) \geq 0 \).
for each $i$. Hence by Theorem 3.2.2, $\text{ind}(f, S) = \sum \text{var}(f, C_i) + 1 \geq 1$. This contradiction with Theorem 3.1.4 completes the proof.

Corollary 7.1.4. Suppose $f : \mathbb{C} \to \mathbb{C}$ is a perfect, oriented map and $X$ is a continuum such that $f(X) \subset T(X)$. Then there exists a point $x_0 \in T(X)$ of period at most 2.

Proof. By Theorem 3.7.4, $f$ is either positively or negatively oriented. In either case, the second iterate $f^2$ is positively oriented and must have a fixed point in $T(X)$ by Theorem 7.1.3.

7.2. Dendrites

Here we generalize Theorem 1.0.2 on the existence of fixed points in invariant dendrites to non-invariant dendrites. We also show that in certain cases the dendrite must contain infinitely many periodic cutpoints. Given two points $a, b$ of a dendrite we denote by $[a, b], \{a, b\}, (a, b)$ the unique closed, semi-open and open arcs connecting $a$ and $b$ in the dendrite. Unless specified otherwise, the situation considered in this subsection is as follows: $D_1 \subset D_2$ are dendrites and $f : D_1 \to D_2$ is a continuous map. Set $E = D_2 \setminus D_1 \cap D_1$. In other words, $E$ consists of points at which $D_2$ “grows” out of $D_1$. Observe that more than one component of $D_2 \setminus D_1$ may “grow” out of a point $e \in E$. We assume that $D_1$ is non-degenerate.

As an important tool we will need the following retraction closely related to the described above situation.

Definition 7.2.1. For each $x \in D_2$ there exists a unique arc (possibly a point) $[x, d_x]$ such that $[x, d_x] \cap D_1 = \{d_x\}$. Hence there exists a natural monotone retraction $r : D_2 \to D_1$ defined by $r(x) = d_x$, and the map $g = g_f = r \circ f : D_1 \to D_1$ is a continuous map of $D_1$ into itself. We call the map $r$ the natural retraction (of $D_2$ onto $D_1$) and the map $g$ the retracted (version of) $f$.

The map $g$ is designed to make $D_1$ invariant so that Theorem 1.0.2 applies to $g$ and allows us to conclude that there are $g$-fixed points. Theorem 7.2.2 extends the result for $\mathbb{R}$ claiming that if there are points $a < b$ in $\mathbb{R}$ mapped by $f$ in different directions, then there exists a fixed point $c \in (a, b)$ (see Introduction, Subsection 1.1). Let us recall the notion of the boundary scrambling property which is first introduced in Definition 5.3.1.

Definition 5.3.1 (Boundary scrambling for dendrites). Suppose that $f$ maps a dendrite $D_1$ to a dendrite $D_2 \supset D_1$. Put $E = D_2 \setminus D_1 \cap D_1$ (observe that $E$ may be infinite). If for each non-fixed point $e \in E$, $f(e)$ is contained in a component of $D_2 \setminus \{e\}$ which intersects $D_1$, then we say that $f$ has the boundary scrambling property or that it scrambles the boundary. Observe that if $D_1$ is invariant then $f$ automatically scrambles the boundary.

We are ready to prove the following theorem.

Theorem 7.2.2. The following claims hold.

1. If $a, b \in D_1$ are such that $a$ separates $f(a)$ from $b$ and $b$ separates $f(b)$ from $a$, then there exists a fixed point $c \in (a, b)$. Thus, if $e_1 \neq e_2 \in E$ are such that $f(e_i)$ belongs to a component of $D_2 \setminus \{e_i\}$ disjoint from $D_1$ then there is a fixed point $c \in (e_1, e_2)$.

2. If $f$ scrambles the boundary, then $f$ has a fixed point in $D_1$. 
Observe that the fixed points found in (1) are cutpoints of $D_1$ (and hence of $D_2$).

Proof. (1) Set $a_0 = a$. Then we find a sequence of points $a_{-1}, a_{-2}, \ldots$ in $(a, b)$ such that $f(a_{-n-1}) = a_{-n}$ and $a_{-n-1}$ separates $a_{-n}$ from $b$. Clearly, $\lim_{n \to \infty} a_{-n} = c \in (a, b)$ is a fixed point as desired (by the assumptions $c$ cannot be equal to $b$). If there are two points $e_1 \neq e_2 \in E$ such that $f(e_i)$ belongs to a component of $D_2 \setminus \{e_i\}$ disjoint from $D_1$ then the above applies to them.

(2) Assume that there are no $f$-fixed points $e \in E$. By Theorem 1.0.2 $g_f = g$ has a fixed point $p \in D_1$. It follows from the fact that $f$ scrambles the boundary that points of $E$ are not $g$-fixed. Hence $p \notin E$.

In general, a $g$-fixed point is not necessarily an $f$-fixed point. In fact, it follows from the construction that if $f(x) \neq g(x)$, then $f$ maps $x$ to a point belonging to a component of $D_2 \setminus D_1$ which “grows” out of $D_1$ at $r \circ f(x) = g(x) \in E$. Thus, since $g(p) = p$ but $p \notin E$, then $g(p) = f(p) = p$.

Remark 7.2.3. It follows from Theorem 7.2.2 that the only behavior of points in $E$ which does not force the existence of a fixed point in $D_1$ is when one point $e \in E$ maps into a component of $D_2 \setminus \{e\}$ disjoint from $D_1$ whereas any other point $e' \in E$ maps into the component of $D_2 \setminus \{e'\}$ which is not disjoint from $D_1$.

Now we suggest conditions under which a map of a dendrite has infinitely many periodic cutpoints; the result will then apply in cases related to complex dynamics. Let us recall the notion of a weakly repelling periodic point which is first introduced in Definition 5.3.2.

Definition 5.3.2 (Weakly repelling periodic points). In the situation of Definition 5.3.1, let $a \in D_1$ be a fixed point and suppose that there exists a component $B$ of $D_1 \setminus \{a\}$ such that arbitrarily close to $a$ in $B$ there exist fixed cutpoints of $D_1$ or points $x$ separating $a$ from $f(x)$. Then we say that $a$ is a weakly repelling fixed point (of $f$ in $B$). A periodic point $a \in D_1$ is said to be simply weakly repelling if there exists $n$ and a component $B$ of $D_1 \setminus \{a\}$ such that $a$ is a weakly repelling fixed point of $f^n$ in $B$.

Now we can prove Lemma 7.2.4.

Lemma 7.2.4. Let $a$ be a fixed point of $f$ and $B$ be a component of $D_1 \setminus \{a\}$. Then the following two claims are equivalent:

1. $a$ is a weakly repelling fixed point for $f$ in $B$;
2. either there exists a sequence of fixed cutpoints of $f|_B$, converging to $a$,
or otherwise, there exists a point $y \in B$ which separates $a$ from $f(y)$ such that there are no fixed cutpoints in the component of $B \setminus \{y\}$ containing $a$ in its closure (in the latter case for any $z \in (a, y]$ the point $z$ separates $f(z)$ from $a$ and each backward orbit of $y$ in $(a, y]$ converges to $a$).

In particular, if $a$ is a weakly repelling fixed point for $f$ in $B$ then $a$ is a weakly repelling fixed point for $f^n$ in $B$ for any $n \geq 1$.

Proof. Let us show that (2) implies (1). We may assume that there exists a point $y \in B$ which separates $a$ from $f(y)$ such that there are no fixed cutpoints in the component $W$ of $B \setminus \{y\}$ containing $a$ in its closure. Choose a point $z \in (a, y)$. Since there are no fixed cutpoints of $f$ in $W$, Theorem 7.2.2(1) implies that $f(z)$ cannot be separated from $y$ by $z$. Hence $f(z)$ is separated from $a$ by $z$, and $a$ is
weakly repelling for $f$ in $B$. Moreover, we can take preimages of $y$ in $(a, y)$, then take their preimages even closer to $a$, inductively. Any so constructed backward orbit of $y$ in $(a, y)$ converges to $a$ because it converges to a fixed point of $f$ in $[a, y]$ and $a$ is the only such fixed point.

Now, suppose that (1) holds. We may assume that there exists a neighborhood $U$ of $a$ in $B$ such that there are no fixed cutpoints of $f$ in $U$. If $a$ is weakly repelling in $B$ for $f$, we can choose a point $y \in U$ so that $y$ separates $a$ from $f(y)$ as desired.

It remains to prove the last claim of the lemma. Indeed, we may assume that there is no sequence of $f^n$-fixed cutpoints in $B$ converging to $a$. Choose a neighborhood $U$ of $a$ which contains no $f^n$-fixed cutpoints in $U \cap B$. By (2) we can choose a point $y \in U \cap B$ such that $y$ separates $a$ from $f(y)$ so that there is a sequence of preimages of $y$ under $f$ which converges to $a$ monotonically. Choosing the $n$-th preimage $z$ we will see that $z$ separates $a$ from $f^n(z)$ with other parts of the second set of conditions of (2) also fulfilled. By the above $a$ is weakly repelling for $f^n$ in $B$ as desired.

Let $B$ be a component of $D_1 \setminus \{a\}$ where $a$ is fixed. Suppose that $a$ is a weakly repelling fixed point for $f$ in $B$ which is not a limit of fixed cutpoints of $f$ in $B$. Since the set of all vertices of $D_2$ together with their images under $f$ and powers of $f$ is countable (see Theorem 10.23 [Nad92]), we can choose $y$ from Lemma 7.2.4 so that $y$ and all cutpoints $x$ in its backward orbit have $\text{val}_{D_2}(x) = 2$. From now on to each fixed point $a$ which is weakly repelling for $f$ in a component $B$ of $D_1 \setminus \{a\}$, but is not a limit point of fixed cutpoints in $B$, we associate a point $x_a \in B$ of valence 2 in $D_2$ separating $a$ from $f(x_a)$ and such that all cutpoints in the backward orbit of $x_a$ are of valence 2 in $D_2$. We also associate to $a$ a semi-neighborhood $U_a$ of $a$ in $\overline{B}$ which is the component of $\overline{B} \setminus \{x_a\}$ containing $a$. We choose $x_a$ so close to $a$ that the diameter of $U_a$ is less than one third of the diameter of $B$.

The next lemma shows that in some cases a fixed point $p$ from Theorem 7.2.2(2) can be chosen to be a cutpoint of $D_1$. Recall that an endpoint of a continuum $X$ is a point $a$ such that the number $\text{val}_X(a)$ of components of $X \setminus \{a\}$ equals 1.

**Lemma 7.2.5.** Suppose that $f$ scrambles the boundary. Then either there is a fixed point of $f$ which is a cutpoint of $D_1$, or, otherwise, there exists a fixed endpoint $a$ of $D_1$ such that if $C_a$ is the component of $D_2 \setminus \{a\}$ containing $D_1 \setminus \{a\}$, then $a$ is not weakly repelling for $f$ in $C_a$.

**Proof.** Suppose that $f$ has no fixed cutpoints. By Theorem 7.2.2(2), the set of fixed points of $f$ is not empty. Hence we may assume that all fixed points of $f$ are endpoints of $D_1$ and, by way of contradiction, $f$ is weakly repelling at any such fixed point $a$ in the component of $D_2 \setminus \{a\}$ containing $D_1 \setminus \{a\}$. Suppose $a$ and $b$ are distinct fixed points of $f$. Let us show that either $U_a \subset U_b$ or $U_b \subset U_a$, or $U_a \cap U_b = \emptyset$. Set $\text{diam}(D_1) = \varepsilon$.

Recall that $x_a, x_b$ are cutpoints of $D_2$ of valence 2. Now, first we assume that $b \in U_a$. If $x_a \notin U_a$, then $U_a \subset U_b$ as desired. Suppose that $x_b \in U_a$. We will show that $x_b \in [b, x_a]$. Indeed, otherwise $U_b$ would contain the component $Q$ of $D_1 \setminus \{x_a\}$, not containing $a$. However, by the choice of the size of $U_a$ we see that $\text{diam}(Q) \geq 2\varepsilon/3$ and therefore $\text{diam}(U_b) \geq 2\varepsilon/3$, a contradiction with the choice of the size of $U_b$. Hence $x_b \in [b, x_a]$ which implies that $U_b \subset U_a$. Now assume that $b \notin U_a$ and $a \notin U_b$. Then it follows that $x_a, x_b \in [a, b]$ and that the order of points in $[b, a]$ is $b, x_b, x_a, a$ which implies that $U_b \cap U_a = \emptyset$. 


Let $B$ be a dendrite. Suppose that $f : D \to D$ is continuous and all its periodic points are weakly repelling. Then $f$ has infinitely many periodic cutpoints.

**Proof.** By way of contradiction we assume that there are finitely many periodic cutpoints of $f$. Let us show that each endpoint $b$ of $D$ with $f(b) = b$ is a weakly repelling fixed point for $f$. Since the only component of $D \setminus \{b\}$ is $D \setminus \{b\}$, we will not be mentioning this component anymore. By the assumptions of the Theorem $b$ is weakly repelling for some power $f^m$ with $m \geq 1$. Then by Lemma 7.2.4 and by the assumption we can choose a point $x \neq b$ such that (1) $x$ is not a vertex or endpoint of $D$, (2) for each point $z \in \{b, x\}$ we have that $z$ separates $b$ from $f^m(z)$, and (3) the component $U$ of $D \setminus \{x\}$ containing $b$, contains no periodic cutpoints.

On the other hand, by way of contradiction we assume that $b$ is not weakly repelling for $f$. Then, again by Lemma 7.2.4, no point $z \in \{b, x\}$ is such that $z$ separates $b$ from $f^m(y)$, a contradiction. To find $y$ we apply the following construction. First, observe that there exists a point $d_1 \in \{b, x\}$ such that $f([b, x]) \supset [b, f(x)] \supset [b, d_1]$. Let $X_1 = \{z \in [b, x] \mid f(z) \in [b, x]\}$ be the set of points mapped into $[b, x]$ by $f$. Then $f(X_1) \supset [b, d_1]$ and all points of $X_1$ map towards $b$ on $[b, x]$. We can apply the same observation to $[b, d_1]$ instead of $[b, x]$. In this way we obtain a point $d_2 \in [b, d_1]$ and a set $X_2 = \{z \in [b, d_1] \mid f(z) \in X_1\}$ such that $[b, d_2] \subset f^2(X_2) \subset [b, d_1]$ and all points of $f(X_2)$ are mapped towards $b$ by $f$. Repeating this argument, we will find points of $[b, x]$ mapped towards $b$ and staying on $[b, x]$ for $m$ steps in a row. This contradicts the previous paragraph and proves that if $b$ is weakly repelling for $f^m$, then it is weakly repelling for $f$. Now by Lemma 7.2.4 $b$ is weakly repelling for $f^n$ for all $n \geq 1$.

Let $f$ have infinitely many periodic cutpoints $a^1, \ldots, a^k$ of $f$. For each $a^i$ there exists $N_i$, such that $a^i$ is fixed for $f^{N_i}$ and there exists a component $B^i$ of $D \setminus \{a^i\}$ such that $a^i$ is weakly repelling for $f^{N_i}$ in $B^i$. Set $N = N_1 \cdots N_k$. Then it follows from Lemma 7.2.4 that each $a^i$ is fixed for $f^N$ and weakly repelling for $f^N$ in $B^i$. Observe that, as we showed above, the endpoints of $D$ which are fixed under $f^N$ are in fact weakly repelling for $f^N$. Without loss of generality we may use $f$ for $f^N$ in the rest of the proof.

Let $A = \bigcup_{i=1}^k a^i$ and let $B$ be a component of $D \setminus A$. Then $\overline{B}$ is a subdendrite of $D$ to which the above tools apply: $D$ plays the role of $D_1$, $\overline{B}$ plays the role of $D_2$, $\partial B$ plays the role of $\partial D_1$, and $E$ is exactly the boundary $\partial B$ of $B$ (by the construction $\partial B \subset A$). Suppose that each point $a \in \partial B$ is weakly repelling in $B$. Then all fixed points of $f$ in $B$ are endpoints of $B$, and all of them are weakly repelling for $f$. Thus, by Lemma 7.2.5 there exists a fixed cutpoint in $B$, a contradiction. Hence for some $a \in \partial B$ we have that $a$ is *not* weakly repelling in $\overline{B} \setminus \{a\}$. By the assumption there exists a

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**7.2. Dendrites**

Consider an open covering of the set of all fixed points $a \in D_1$ by their neighborhoods $U_a$ and choose a finite subcover. By the above we may assume that it consists of pairwise disjoint sets $U_{a_1}, \ldots, U_{a_k}$. Consider the component $Q$ of $D_1 \setminus \{x_{a_1}, \ldots, x_{a_k}\}$ whose endpoints are the points $x_{a_1}, \ldots, x_{a_k}$ and perhaps some endpoints of $D_1$. Then $f|_Q$ is fixed point free. On the other hand, $f|_D$, scrambles the boundary, and hence it is easy to see that $f|_Q$, with $Q$ considered as a subdendrite of $D_2$, scrambles the boundary too, a contradiction with Theorem 7.2.2. □

Lemma 7.2.5 is helpful in the next theorem.

**Theorem 7.2.6.** Let $D$ be a dendrite. Suppose that $f : D \to D$ is continuous and all its periodic points are weakly repelling. Then $f$ has infinitely many periodic cutpoints.

**Proof.** By way of contradiction we assume that there are finitely many periodic cutpoints of $f$. Let us show that each endpoint $b$ of $D$ with $f(b) = b$ is a weakly repelling fixed point for $f$. Since the only component of $D \setminus \{b\}$ is $D \setminus \{b\}$, we will not be mentioning this component anymore. By the assumptions of the Theorem $b$ is weakly repelling for some power $f^m$ with $m \geq 1$. Then by Lemma 7.2.4 and by the assumption we can choose a point $x \neq b$ such that (1) $x$ is not a vertex or endpoint of $D$, (2) for each point $z \in \{b, x\}$ we have that $z$ separates $b$ from $f^m(z)$, and (3) the component $U$ of $D \setminus \{x\}$ containing $b$, contains no periodic cutpoints.

On the other hand, by way of contradiction we assume that $b$ is not weakly repelling for $f$. Then, again by Lemma 7.2.4, no point $z \in \{b, x\}$ is such that $z$ separates $b$ from $f^m(y)$, a contradiction. To find $y$ we apply the following construction. First, observe that there exists a point $d_1 \in \{b, x\}$ such that $f([b, x]) \supset [b, f(x)] \supset [b, d_1]$. Let $X_1 = \{z \in [b, x] \mid f(z) \in [b, x]\}$ be the set of points mapped into $[b, x]$ by $f$. Then $f(X_1) \supset [b, d_1]$ and all points of $X_1$ map towards $b$ on $[b, x]$. We can apply the same observation to $[b, d_1]$ instead of $[b, x]$. In this way we obtain a point $d_2 \in [b, d_1]$ and a set $X_2 = \{z \in [b, d_1] \mid f(z) \in X_1\}$ such that $[b, d_2] \subset f^2(X_2) \subset [b, d_1]$ and all points of $f(X_2)$ are mapped towards $b$ by $f$. Repeating this argument, we will find points of $[b, x]$ mapped towards $b$ and staying on $[b, x]$ for $m$ steps in a row. This contradicts the previous paragraph and proves that if $b$ is weakly repelling for $f^m$, then it is weakly repelling for $f$. Now by Lemma 7.2.4 $b$ is weakly repelling for $f^n$ for all $n \geq 1$.

Let $f$ have infinitely many periodic cutpoints $a^1, \ldots, a^k$ of $f$. For each $a^i$ there exists $N_i$ such that $a^i$ is fixed for $f^{N_i}$ and there exists a component $B^i$ of $D \setminus \{a^i\}$ such that $a^i$ is weakly repelling for $f^{N_i}$ in $B^i$. Set $N = N_1 \cdots N_k$. Then it follows from Lemma 7.2.4 that each $a^i$ is fixed for $f^N$ and weakly repelling for $f^N$ in $B^i$. Observe that, as we showed above, the endpoints of $D$ which are fixed under $f^N$ are in fact weakly repelling for $f^N$. Without loss of generality we may use $f$ for $f^N$ in the rest of the proof.

Let $A = \bigcup_{i=1}^k a^i$ and let $B$ be a component of $D \setminus A$. Then $\overline{B}$ is a subdendrite of $D$ to which the above tools apply: $D$ plays the role of $D_2$, $\overline{B}$ plays the role of $D_1$, and $E$ is exactly the boundary $\partial B$ of $B$ (by the construction $\partial B \subset A$). Suppose that each point $a \in \partial B$ is weakly repelling in $B$. Then all fixed points of $f$ in $B$ are endpoints of $B$, and all of them are weakly repelling for $f$. Thus, by Lemma 7.2.5 there exists a fixed cutpoint in $B$, a contradiction. Hence for some $a \in \partial B$ we have that $a$ is *not* weakly repelling in $\overline{B} \setminus \{a\}$. By the assumption there exists a
component, say, $C'$, of $D \setminus \{a\}$ disjoint from $B$ such that $a$ is weakly repelling in $C'$. Let $C$ be the component of $D \setminus A$ non-disjoint from $C'$ with $a \in \partial C$.

We can now apply the same argument to $C$. If all boundary points of $C$ are weakly repelling for $f$ in $C$, then by Lemma 7.2.5 $C$ will contain a fixed cutpoint, a contradiction. Hence there exists a point $d \in A$ such that $d$ is not weakly repelling for $f$ in $C$ and a component $F$ of $D \setminus A$ whose closure meets $C$ at $d$, and $d$ is weakly repelling in $F$. Note that $B \cap F = \emptyset$. Clearly, after finitely many steps this process will have to end (recall, that $D$ is a dendrite), ultimately leading to a component $Z$ of $D \setminus A$ such that all fixed points of $f$ in $Z$ are endpoints of $Z$ at which $f$ is weakly repelling. Again, Lemma 7.2.5 applies to $Z$ and there exists a fixed cutpoint in $Z$, a contradiction. □

An important application of Theorem 7.2.6 is to dendritic topological Julia sets. They can be defined as follows. Consider an equivalence relation $\sim$ on the unit circle $S^1 \subseteq \mathbb{C}$. Equivalence classes of $\sim$ will be called ($\sim$-)classes and will be denoted by boldface letters. A $\sim$-class consisting of two points is called a leaf; a class consisting of at least three points is called a gap. Fix an integer $d > 1$ and define the map $\sigma_d : S^1 \rightarrow S^1$ by $\sigma_d(z) = z^d$, where $z$ is a complex number with $|z| = 1$. Then the equivalence $\sim$ is said to be a $(d\text{-})$invariant lamination (this is more restrictive than Thurston’s definition in [Thu09]) if:

1. $\sim$ is closed: the graph of $\sim$ is a closed subset of $S^1 \times S^1$;
2. $\sim$ defines a lamination, i.e., it is unlinked: if $g_1$ and $g_2$ are distinct $\sim$-classes, then their convex hulls $\text{Ch}(g_1), \text{Ch}(g_2)$ in the unit disk $\mathbb{D}$ are disjoint,
3. $\sim$ is forward invariant: for a $\sim$-class $g$, the set $\sigma_d(g)$ is a $\sim$-class too which implies that
4. $\sim$ is backward invariant: for a $\sim$-class $g$, its preimage $\sigma_d^{-1}(g) = \{x \in S^1 : \sigma_d(x) \in g\}$ is a union of $\sim$-classes;
5. for any gap $g$, the map $\sigma_d|g : g \rightarrow \sigma_d(g)$ is a map with positive orientation, i.e., for every connected component $(s,t)$ of $S^1 \setminus g$ the arc $(\sigma_d(s), \sigma_d(t))$ is a connected component of $S^1 \setminus \sigma_d(g)$.

The lamination in which all points of $S^1$ are equivalent is said to be degenerate. It is easy to see that if a forward invariant lamination $\sim$ has a $\sim$-class with non-empty interior then $\sim$ is degenerate. Hence equivalence classes of any non-degenerate forward invariant lamination are totally disconnected.

Let $\sim$ define an invariant lamination. A $\sim$-class $g$ is periodic if $\sigma^n(g) = g$ for some $n \geq 1$. Let $p : S^1 \rightarrow J_\sim = S^1 / \sim$ be the quotient map of $S^1$ onto its quotient space $J_\sim$. We can extend the equivalence relation $\sim$ to an equivalence relation $\equiv$ of the entire plane by defining $x \equiv y$ if either $x$ and $y$ are contained in the convex hull of one equivalence class of $\sim$, or $x = y$. Then the quotient map $m : \mathbb{C} \rightarrow \mathbb{C}/ \equiv$ is a monotone map whose point inverses are convex continua or points. Note that $p(S^1) = S^1 / \sim = m(\mathbb{D}) = \mathbb{D}/ \equiv$. Let $f_\sim : J_\sim \rightarrow J_\sim$ be the map induced by $\sigma_d$. We call $J_\sim$ a topological Julia set and the induced map $f_\sim$ a topological polynomial. Recall that a branched covering map $f : X \rightarrow Y$ is a finite-to-one and open map for which there exists a finite set $F \subset Y$ such that $f|_{X \setminus f^{-1}(F)}$ is a covering map. Note that $f_\sim$ is a branched covering map, and in particular, $f_\sim$ has finitely many critical points (i.e., points where $f$ is not locally one-to-one). It is easy to see that if $g$ is a $\sim$-class then $\text{val}_{f_\sim}(p(g)) = |g|$ where by $|A|$ we denote the cardinality of a set $A$. 


Suppose that the topological Julia set $J_\infty$ is a dendrite and $f_\sim : J_\infty \to J_\infty$ is a topological polynomial. Then all periodic points of $f_\sim$ are weakly repelling and $f_\sim$ has infinitely many periodic cutpoints.

**Proof.** Suppose that $x$ is an $f_\sim$-fixed point and set $g = p^{-1}(x)$. Then $\sigma_g(g) = g$. Suppose first, that $x$ is an endpoint of $J_\infty$. Then $g$ is a singleton. Choose $y \neq x \in J_\infty$. Then the unique arc $[x,y] \subset J_\infty$ contains points $y_k \to x$ of valence 2 because there are no more than countably many vertices of $J_\infty$ (see Theorem 10.23 in [Nad92]). It follows that $\sim$-classes $p^{-1}(y_k)$ are leaves separating $g$ from the rest of the circle and repelled from $g$ under the action of $\sigma_d$ which is expanding. Hence $f_\sim(y_k)$ is separated from $x$ by $y_k$ and so $x$ is weakly repelling.

Suppose that $x$ is not an endpoint. Choose a very small connected neighborhood $U$ of $x$. It is easy to see that each component $A$ of $U \setminus \{x\}$ corresponds to a single non-degenerate chord $\ell_A$ in the boundary of the Euclidean convex hull, $\text{Ch}(g) = G$, of $g$. Recall that $S^1 = \mathbb{R} \setminus \mathbb{Z}$ and that the endpoints $a_A$ and $b_A$ of $\ell_A$ are points in $S^1$. Denote by $\sigma_d(\ell_A)$ the chord with endpoints $\sigma_d(a_A)$ and $\sigma_d(b_A)$ and by $|\ell_A| = \min\{|a_A - b_A|, 1 - |a_A - b_A|\}$, the length of $\ell_A$. Since $f_\sim$ is a branched covering map, for each component $A$ of $U \setminus \{x\}$ there exists a unique component $B = h(A)$ of $U \setminus \{x\}$ such that $f_\sim(A) \cap B \neq \emptyset$. This defines a map $h$ from the set $\mathcal{A}$ of all components of $U \setminus \{x\}$ to itself. It follows that for each chord $\ell_A \subset \partial G$, $\sigma_d(\ell_A)$ is a non-degenerate chord in $\partial G$.

Suppose that there exist $\ell_A \subset \partial G$ and $n > 0$ such that $\sigma_d^n(\ell_A) = \ell_A$. Then it follows that the endpoints of $\ell_A$ are fixed under $\sigma_d^{2n}$. Connect $x$ to a point $y \in A$ with the arc $[x,y]$, and choose, as in the first paragraph, a sequence of points $y_k \in [x,y], y_k \to x$ of valence 2. Then again by the expanding properties of $\sigma_d^n$ it follows that $f_\sim(y_k)$ is separated from $x$ by $y_k$ and so $x$ is weakly repelling (for $f_\sim^{2n}$ in $A$).

It remains to show that there must exist a component $A$ of $U \setminus \{x\}$ with $\sigma_d^n(\ell_A) = \ell_A$ for some $n > 0$. Clearly $\partial G$ can contain at most finitely many chords $\ell_A$ such that its length $L(\ell_A) \geq 1/(2(d+1))$. If $L(\ell) < 1/(2(d+1))$, then $L(\sigma_d(\ell)) = d \cdot L(\ell)$ (i.e. $\sigma_d$ expands the length of small leaves by the factor $d$).

Since the family of chords in the boundary of $G$ is forward invariant and for each chord $\ell_A$ with $L(\ell_A) < 1/(2(d+1))$, $L(\sigma_d(\ell_A)) = d \cdot L(\ell_A)$, such a periodic chord must exist (since if this is not the case there must exist an infinite number of distinct leaves in the boundary of $G$ of length bigger than $1/(2(d+1))$, a contradiction.

Hence all periodic points of $f_\sim$ are weakly repelling and by Theorem 7.2.6 $f_\sim$ has infinitely many periodic cutpoints. 

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**7.3. Non-invariant continua and positively oriented maps of the plane**

In this subsection we will extend Theorem 7.1.3 and obtain a general fixed point theorem which shows that if a non-separating plane continuum, not necessarily invariant, maps in an appropriate way, then it contains a fixed point. However we begin with Lemma 7.3.1 which gives a sufficient condition for the non-negativity of the variation of an arc.

**Lemma 7.3.1.** Let $f : C \to C$ be positively oriented, $X$ a continuum and $C = [a,b]$ a bumping arc of $X$ such that $f(a), f(b) \in X$ and $f(C) \cap C = \emptyset$. Let $v \in C \setminus X$, and let $J_v$ be a continuum at $v \in C$ defined as in Definition 2.2.2. If there exists a continuum $K$ disjoint from $J_v$ such that $C$ is a bumping arc of $K$. 

Consider two cases. First, let $f$ for all $v \in D$. Suppose that $D$ whose closures contain $z$. Previous section. For each $v$ in both cases.

Then the map $f$ is said to strongly scramble the boundary (of $X$). If instead of (3) we have (3a) for all $i$, either $f(K_i) \subset K_i$, or $f(K_i) \cap Z_i = \emptyset$ then we say that $f$ strongly scrambles the boundary (of $X$); clearly, if $f$ strongly scrambles the boundary of $X$, then it scrambles the boundary of $X$. In either case, the continua $K_i$ are called exit continua (of $X$).

We will always use the same notation (for $X$, $Z_i$ and $K_i$) introduced in Definition 5.4.1 unless explicitly stated otherwise. Let us make a few remarks. First, even though we use the notion only for positively oriented maps $f$, the definitions can be given for all continuous functions. Also, observe, that in the situation of Definition 5.4.1 if $X$ is invariant then $f$ automatically strongly scrambles the boundary because the set of exit continua can be taken to be empty. We will also agree that the choice of the sets $Z_i$ is optimal in the sense that if $(f(X) \setminus X) \cap Z_i = \emptyset$ for some $i$, then the set $Z_i$ will be removed from the list. In particular, all continua $Z_i$ contain points from $f(X) \setminus X$ and hence all continua $K_i$ have points from $f(X) \setminus X \cap X$.

By Remark 5.4.2 $X \cup (\bigcup Z_i)$ is a non-separating continuum. Suppose that we are in the situation of the previous section, $D_1 \subset D_2$ are dendrites, $E = D_2 \setminus D_1 \cap D_1 = \{z_1, \ldots, z_l\}$ is finite, and $f : D_1 \to D_2$ scrambles the boundary in the sense of the previous section. For each $z_i$ consider the union of all components of $D_2 \setminus D_1$ whose closures contain $z_i$, and denote by $Z_i$ the closure of their union. In other words, $Z_i$ is the closed connected piece of $D_2$ which "grows" out of $D_1$ at $z_i$. Then $Z_i \cap D_1 = \{z_i\}$, each $Z_i$ is a dendrite itself, and $f$ strongly scrambles the boundary in the sense of the new definition too. This explains why we use similar terminology in both cases.

From now on we fix a positively oriented map $f$. Even though some of the main applications of the results are to polynomial maps, this generality is well justified because in some arguments (e.g., when dealing with parabolic points) we have to locally perturb our map to make sure that the local index at a parabolic
fixed point equals 1, and this leads to the loss of analytic properties of the map (see Lemma 7.5.1). Let us now prove the following technical lemma.

**Lemma 7.3.2.** Suppose that \( f \) is positively oriented, scrambles the boundary of \( X \), \( Q \) is a bumping arc of \( X \) such that its endpoints map back into \( X \) and \( f(Q) \cap Q = \emptyset \). Then \( \text{var}(f, Q) \geq 0 \).

**Proof.** Suppose first that \( Q \setminus \bigcup Z_i \neq \emptyset \) and choose \( v \in Q \setminus \bigcup Z_i \). Since \( v \in Q \setminus \bigcup Z_i \) and \( X \cup (\bigcup Z_i) \) is non-separating, there exists a junction \( J_v \), with \( v \in Q \), such that \( J_v \setminus [X \cup Q \cup (\bigcup Z_i)] = \{v\} \) and, hence, \( J_v \cap f(X) \subset \{v\} \). Now the desired result follows from Lemma 7.3.1.

Suppose now that \( Q \setminus \bigcup Z_i = \emptyset \). Then \( Q \subset Z_i \) for some \( i \) and so \( Q \cap X \subset K_i \). In particular, both endpoints of \( Q \) belong to \( K_i \). Choose a point \( v \in Q \). Then again there is a junction connecting \( v \) and infinity outside \( X \) (except possibly for \( v \)). Since all sets \( Z_j, j \neq i \) are positively distant from \( v \) and \( X \cup (\bigcup_{j \neq i} Z_i) \) is non-separating, the junction \( J_v \) can be chosen to avoid all sets \( Z_j, j \neq i \). Now, by part (3) of Definition 5.4.1, \( f(K_i) \cap J_v \subset \{v\} \), hence by Lemma 7.3.1 \( \text{var}(f, Q) \geq 0 \). \( \square \)

Lemma 7.3.2 is applied in Theorem 7.3.3 in which we show that a map which strongly scrambles the boundary has fixed points. In fact, Lemma 7.3.2 is a major technical tool in our other results too. Indeed, suppose that a positively oriented map \( f \) scrambles the boundary of \( X \). If we can construct a bumping simple closed curve \( S \) around \( X \) which has a partition into bumping arcs (links of \( S \)) whose endpoints map into \( X \) (or at least into \( T(S) \)) and whose images are disjoint from themselves, then Lemma 7.3.2 would imply that the variation of \( S \) is non-negative. By Theorem 3.2.2 this would imply that the index of \( S \) is positive. Hence by Theorem 3.1.4 there are fixed points in \( T(S) \). Choosing \( S \) to be sufficiently tight around \( X \) we see that there are fixed points in \( X \). Thus, the construction of a tight bumping simple closed curve \( S \) with a partition satisfying the above listed properties becomes a major task.

For the sake of convenience we now sketch the proof of Theorem 7.3.3 which allows us to emphasize the main ideas rather than details. The main steps in constructing \( S \) are as follows. First we assume by way of contradiction that \( f \) has no fixed points in \( X \). By Theorem 7.1.3 then \( f(X) \not\subset X \) and \( f(K_i) \not\subset K_i \) for any \( i \). By the definition of strong scrambling then \( f(K_i) \) is “far away” from \( Z_i \) for any \( i \). Choose a tight bumping simple closed curve \( S \) around \( X \) with very small links. We need to construct a partition of \( S \) into bumping arcs whose endpoints map into \( X \) (or at least into \( T(S) \)) and whose images are disjoint from themselves. Since there are no fixed points in \( X \), we may assume that all links of \( S \) move off themselves. However some of them may have endpoint(s) mapping outside \( X \) which prevents the corresponding partition from being the one we are looking for. So, we enlarge these links by consecutive concatenating them to each other until the images of the endpoints of these concatenations are inside \( X \) and these concatenations still map off themselves (the latter needs to be proven which is a big part of the proof of Theorem 7.3.3).

The bumping simple closed curve \( S \) then remains as before, but the partition changes because we enlarge some links. Still, the construction shows that the new partition is satisfactory, and since \( S \) can be chosen arbitrarily tight, this implies the existence of a fixed point in \( X \) as explained before. Thus, a new development
is that we are able to construct a partition of \( S \) which has all the above listed necessary properties having possibly very long links.

To achieve the goal of replacing some links in \( S \) by their concatenations we consider the links with at least one endpoint mapped outside \( X \) in detail (indeed, Lemma 7.3.2 already applies to all other links) and use the fact that \( f \) strongly scrambles the boundary. The idea is to consider consecutive links of \( S \) with endpoints mapped into \( Z_i \setminus X \). Their concatenation is a connected piece of \( S \) with endpoints (and a lot of other points) belonging to \( X \) and mapping into one \( Z_i \). If we begin the concatenation right before the images of links enter \( Z_i \setminus X \) and stop it right after the images of the links exit \( Z_i \setminus X \) we will have one condition of Lemma 7.3.2 satisfied because the endpoints of the thus constructed new “big” concatenation link \( T \) of \( S \) map into \( X \).

We need to verify that \( T \) moves off itself under \( f \). This is easy to see for the end-links of \( T \): each end-link has the image “crossing” into \( X \) from \( Z_i \setminus X \), hence the images of end-links are close to \( K_i \). However the set \( K_i \) is mapped “far away” from \( Z_i \) by the definition of strong scrambling and because none of the \( K_j \)'s is invariant by the assumption. This implies that the end-links themselves must be far away from \( K_j \) and hence from \( Z_j \). If now we move from link to link inside \( T \) we see that those links cannot approach \( Z_i \) too closely because if they do, they will have to “be close to \( K_i \)”, and their images will have to be close to the image of \( K_i \) which is far away from \( Z_i \), a contradiction with the fact that all links in \( T \) have endpoints which map into \( Z_i \setminus X \). In other words, the dynamics of \( K_i \) prevents the new bigger links from getting even close to \( Z_i \) under \( f \) which shows that \( T \) moves off itself as desired (after all, the images of new bigger links are close to the set \( Z_i \setminus X \)).

Given a compact set \( K \) denote by \( B(K, \varepsilon) \) the open set of all points whose distance to \( K \) is less than \( \varepsilon \). By \( d(\cdot, \cdot) \) we denote the distance between two points or sets.

**Theorem 7.3.3.** Suppose \( f : \mathbb{C} \to \mathbb{C} \) is positively oriented, \( X \) is a non-separating continuum and \( f \) strongly scrambles the boundary of \( X \). Then \( f \) has a fixed point in \( X \).

**Proof.** If \( f(X) \subset X \) then the result follows from Theorem 7.1.3. Similarly, if there exists \( i \) such that \( f(K_i) \subset K_{i+1} \), then \( f \) has a fixed point in \( K_i \subset X \) and we are also done. Hence we may assume that \( f(X) \setminus X \neq \emptyset \), there are \( m > 0 \) sets \( Z_i, i = 1, \ldots, m \), \( (f(X) \setminus X) \cap Z_i \neq \emptyset \) for any \( i \), and \( f(K_i) \cap Z_i = \emptyset \) for all \( i \) (making these claims we rely upon the fact that \( f \) strongly scrambles the boundary). Suppose that \( f_X \) is fixed point free. Then there exists \( \varepsilon > 0 \) such that for all \( x \in X \), \( d(x, f(x)) > \varepsilon \). We may assume that \( 2 \varepsilon < \min \{ d(Z_i, Z_j) \mid i \neq j \} \). We now choose constants \( \eta', \eta, \delta \) and a bumping simple closed curve \( S \) (whose initial links are crosscuts) of \( X \) so that the following holds.

1. \( 0 < \eta' < \eta < \delta < \varepsilon/3 \).
2. For each \( x \in X \setminus B(K_i, 3\delta) \) we have \( d(f(x), Z_i) > 3\delta \).
3. For each \( x \in X \setminus B(K_i, 3\delta) \) we have \( d(x, Z_i) > 3\eta \).
4. For each \( i \) there is a point \( x_i \in X \) with \( f(x_i) = z_i \in Z_i \) and \( d(z_i, X) > 3\eta \).

Since by Theorem 3.7.4 \( \partial f(X) \subset f(\partial X) \) and \( X \) is non-separating, we may assume that \( x_i \in \partial X \).

5. \( X \subset T(S) \) and \( A = X \cap S = \{ a_0 < \cdots < a_n < a_{n+1} = a_0 \} \) with points of \( A \) numbered in the positive circular order around \( S \).
(6) \( f|_{T(S)} \) is fixed point free.

(7) For any \( Q_i = (a_i, a_{i+1}) \subset S \), \( \text{diam}(Q_i) + \text{diam}(f(Q_i)) < \eta \).

(8) For any \( x, y \in X \) with \( d(x, y) < \eta' \) we have \( d(f(x), f(y)) < \eta \).

(9) \( A \) is an \( \eta' \)-net in \( \partial X \) (i.e., the Hausdorff distance between \( A \) and \( \partial X \) is less than \( \eta' \)).

Observe that \( Q_i \) is contained in the closed ball centered at \( a_i \) of radius \( \text{diam}(Q_i) \) and \( f(Q_i) \) is contained in the closed ball centered at \( f(a_i) \) of radius \( \text{diam}(f(Q_i)) \); hence by (7) and since \( d(x, f(x)) > \varepsilon \) for all \( x \in X \) we see that \( Q_i \cap f(Q_i) = \emptyset \) for every \( i \) (we rely on the triangle inequality here too).

Claim 1. There exists a point \( a_j \in A \) such that \( f(a_j) \in X \setminus \bigcup_i B(Z_i, \eta) \).

Proof of Claim 1. Set \( B(Z_i, 3\eta) = T_i \). We will show that there exists a point \( x \in \partial X \) with \( f(x) \in X \setminus \bigcup T_i \). Indeed, suppose first that \( m = 1 \). Then by (2) and the assumption that \( f(K_1) \cap Z_1 = \emptyset \) for each \( i \) we have \( f(K_1) \subset X \setminus T_1 \), and we can choose any point of \( K_1 \cap \partial X \) as \( x \). Now, suppose that \( m \geq 2 \). Observe that by the choice of \( \varepsilon \) and by (1) the compacta \( T_i \) are pairwise disjoint. By (4) for each \( i \) there are points \( x_i \) in \( \partial X \) such that \( f(x_i) \in Z_i \subset T_i \). Since the sets \( f^{-1}(T_i) \cap X \) are pairwise disjoint non-empty compacta we see that the set \( V = \partial X \setminus \bigcup f^{-1}(T_i) \) is non-empty (because \( \partial X \) is a continuum). Now we can choose any point of \( V \) as \( x \).

Notice now that by the choice of \( A \) (see (9)) we can find a point \( a_j \) such that \( d(a_j, x) \leq \eta' \) which by (8) implies that \( d(f(a_j), f(x)) < \eta \) and hence \( f(a_j) \in X \setminus (\bigcup B(Z_i, \eta)) \) as desired. \( \square \)

By Claim 1, we assume without loss of generality, that \( f(a_0) \in X \setminus \bigcup B(Z_i, \eta) \). Now, by (4) there exists a point \( x_1 \) such that \( f(x_1) = z_1 \) is more than \( 3\eta \)-distant
Given a map $f$, since by (7) for each link $\ell$ we have diam$(f(Q_\ell)) < \varepsilon/3$ for any $u$ and $d(Z_n Z_l) > 2\varepsilon$ for all $s \neq t$, there exists $i(1) = 1$ such that for each $i \in \{1, \ldots, n\}$ with $f(a_i) \in Z_i(1)$ for all $n(1) \leq r \leq m(1) - 1$.

(3) $f(a_m(1)) \in X$.

Consider the arc $Q' = \{a_{n(1) - 1}, a_{m(1)}\} \subset S$ and show that $f(Q') \cap Q' = \emptyset$. As we walk along $Q'$ and mark the $f$-images of points $a_{n(1) - 1}, a_{n(1)}, \ldots, a_{m(1)}$, we begin in $X$ at $f(a_{n(1) - 1})$, then enter $Z_i(1) \setminus X$ and walk inside it, and then exit $Z_i(1) \setminus X$ at $f(a_{m(1)}(1)) \in X$. Since every step in this walk is rather short (by (7) $\text{diam}(Q_i) + \text{diam}(f(Q_i)) < \eta$), we see that $d(f(a_{n(1) - 1}), Z_i(1) \setminus X) < \eta$. On the other hand, for each $r, n(1) \leq r \leq m(1) - 1$, we have $f(a_r) \in Z_i(1) \setminus X$. Thus, $d(f(a_r), Z_i(1)) < \eta$ for each $n(1) - 1 \leq r \leq m(1)$. Since by (7) for each link $Q$ of $S$ we have $\text{diam}(Q) + \text{diam}(f(Q)) < \eta$, we now see by the triangle inequality that $d(f(Q'), Z_i(1)) < 2\eta$.

This implies that for $n(1) - 1 \leq r \leq m(1)$, $d(a_{r}, K_i(1)) > 3\delta$ (because otherwise by (2) $f(a_r)$ would be farther away from $Z_i(1)$ than $3\delta > \eta$, a contradiction) and so $d(a_r, Z_i(1)) > 3\eta$ (because $a_r \in X \setminus B(K_i(1), 3\delta)$ and by (3)). Since by (7) for each link of $S$ we have $\text{diam}(Q) + \text{diam}(f(Q)) < \eta$, then $d(Q', Z_i(1)) > 2\eta$.

Therefore $f(Q') \cap Q' = \emptyset$. This allows us to replace the original division of $S$ into links $Q_0, \ldots, Q_{m(1) - 1}$ by a new one in which $Q'$ plays the role of a new link; in other words, we simply delete the points $\{a_{n(1) - 1}, \ldots, a_{m(1) - 1}\}$ from $A$. Thus, $Q'$ is a bumping arc whose endpoints map back into the continuum $X$ and such that $f(Q') \cap Q' = \emptyset$. Therefore $Q'$ satisfies the conditions of Lemma 7.3.2, and so $\text{var}(f, Q') \geq 0$. Observe also that for $Q'$ the associated continuum $Z_i(1)$ is well-defined because the distance between distinct continua $Z_i$ is greater than $2\varepsilon$. Replace the string of links $\{Q_0, \ldots, Q_{m(1) - 1}\}$ in $S$ by the single link $Q' = Q'_0$ which has as endpoints $a_{n(1) - 1}$ and $a_{m(1)}$. Continuing in the same manner and moving along $S$, in the end we obtain a finite set $A = \{a_0 = a'_0 < a'_1 < \cdots < a'_k\} \subset A$ such that for each $i$ we have $f(a'_i) \in X \subset T(S)$ and for each arc $Q'_i = [a'_i, a'_{i+1}]$ we have $f(Q'_i) \cap Q'_i = \emptyset$. In other words, we will construct a partition of $S$ satisfying all the required properties: its links are bumping arcs whose endpoints map back into $X$ and whose images are disjoint from themselves. As outlined after Lemma 7.3.2, this yields a contradiction. More precisely, by Theorem 3.2.2, $\text{ind}(f, S) = \sum \text{var}(f, Q'_i) + 1$, and since by Lemma 7.3.2, $\text{var}(f, Q'_i) \geq 0$ for all $i$, $\text{ind}(f, S) \geq 1$ contradicting the fact that $f$ is fixed point free in $T(S)$ (see Theorem 3.1.4).

### 7.4. Maps with isolated fixed points

In this section we assume that all maps $f : C \to C$ are positively oriented maps with isolated fixed points.

**Definition 7.4.1.** Given a map $f : X \to Y$ we say that $c \in X$ is a critical point of $f$ if for each neighborhood $U$ of $c$, there exist $x_1 \neq x_2 \in U$ such that $f(x_1) = f(x_2)$. Hence, if $x$ is not a critical point of $f$, then $f$ is locally one-to-one near $x$. 


If a point \( x \) belongs to a non-degenerate continuum collapsed to a point under \( f \) then \( x \) is critical; also any point which is an accumulation point of collapsing continua is critical. However in these cases the map near \( x \) may be monotone. A more interesting case is when the map near \( x \) is not monotone; then \( x \) is a branchpoint of \( f \) and it is critical even if there are no collapsing continua close by. One can define the local degree \( \deg_f(a) \) as the number of components of \( f^{-1}(y) \) non-disjoint from a small neighborhood of \( a \) (\( y \) then should be chosen close to \( f(a) \)). It is well-known that for a positively oriented map \( f \) and a point \( a \) which is a component of \( f^{-1}(f(a)) \) the local degree \( \deg_f(a) \) equals the winding number \( \text{win}(f, S, f(a)) \) for any small simple closed curve \( S \) around \( a \). Then branchpoints are exactly the points at which the local degree is greater than 1. Notice that since we do not assume any smoothness, a critical point may well be both fixed (periodic) and topologically repelling in the sense that some small neighborhoods of \( c = f(c) \) map over themselves by \( f \).

Let us recall the notion of a local index of a map at a point which is first introduced in Definition 5.4.3.

**Definition 5.4.3.** Suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a positively oriented map with isolated fixed points and \( x \) is a fixed point of \( f \). Then the local index of \( f \) at \( x \), denoted by \( \text{ind}(f, x) \), is defined as \( \text{ind}(f, S) \) where \( S \) is a small simple closed curve around \( x \).

It is easy to see that, since \( f \) is positively oriented and has isolated fixed points, the local index is well-defined, i.e. does not depend on the choice of \( S \). By modifying a translation map one can give an example of a homeomorphism of the plane which has exactly one fixed point \( x \) with local index 0. Still in some cases the local index at a fixed point must be positive.

**Definition 7.4.2.** Let \( f : \mathbb{C} \to \mathbb{C} \) be a map. A fixed point \( x \) is said to be **topologically repelling** if there exists a sequence of simple closed curves \( S_j \to \{x\} \) such that \( x \in \text{int}(T(S_j)) \subset T(S_j) \subset \text{int}(T(f(S_j))) \). A fixed point \( x \) is said to be **topologically attracting** if there exists a sequence of simple closed curves \( S_j \to \{x\} \) not containing \( x \) and such that \( x \in \text{int}(T(f(S_j))) \subset T(f(S_j)) \subset \text{int}(T(S_j)) \).

**Lemma 7.4.3.** Let \( f : \mathbb{C} \to \mathbb{C} \) be a positively oriented map with isolated fixed points. If \( a \) is a topologically repelling fixed point then we have that \( \text{ind}(f, a) = \deg_f(a) \geq 1 \). If however \( a \) is a topologically attracting fixed point then \( \text{ind}(f, a) = 1 \).

**Proof.** Consider the case of the repelling fixed point \( a \). Then it follows that, as \( x \) runs along a small simple closed curve \( S \) with \( a \in T(S) \), the vector from \( x \) to \( f(x) \) produces the same winding number as the vector from \( a \) to \( f(x) \). As we remarked before, it is well-known that this winding number equals \( \deg_f(a) \); on the other hand, \( \text{ind}(f, S) > 0 \) since \( f \) is positively oriented and has isolated fixed points. The argument for an attracting fixed point is similar.

If however a fixed point \( x \) is neither topologically repelling nor topologically attracting, then \( \text{ind}(f, x) \) could be greater than 1 even in the non-critical case. Indeed, by definition \( \text{ind}(f, x) \) coincides with the winding number of \( f(z) - z \) on a small simple closed curve \( S \) around \( x \) with respect to the origin. If, e.g., \( f \) is rational and \( f'(x) \neq 1 \) then this implies that \( \text{ind}(f, x) = 1 \). However if \( f'(x) = 1 \) then \( \text{ind}(f, x) \) is the multiplicity at \( x \) (i.e., the local degree of the map \( f(z) - z \)
Suppose that \( f \) is positively oriented and has isolated fixed points. Then for any simple closed curve \( S \subset \mathbb{C} \), which contains no fixed points of \( f \), its index equals the sum of local indices taken over all fixed points in \( T(S) \). In particular if for each fixed point \( p \in T(S) \) we have that \( \text{ind}(f, x) = 1 \) then \( \text{ind}(f, S) \) equals the number \( n(f, S) \) of fixed points inside \( T(S) \).

Theorem 7.4.4 implies Theorem 3.1.4 for positively oriented maps with isolated fixed points (indeed, if \( \text{ind}(f, S) \neq 0 \) then by Theorem 7.4.4 there must exist fixed points in \( T(S) \)), and actually provides more information. By the above analysis, Lemma 7.4.3 and Theorem 7.4.4, \( \text{ind}(f, S) \) equals the number \( n(f, S) \) of fixed points inside \( T(S) \) if all \( f \)-fixed points in \( T(S) \) are either topologically attracting, or such that if \( f \) has a complex derivative \( f' \) at \( x \), and \( f'(x) \neq 1 \); if \( f \)-fixed points can also be topologically repelling, then \( \text{ind}(f, S) \geq n(f, S) \).

In the spirit of the previous parts of the paper, we are still concerned with finding \( f \)-fixed points inside non-invariant continua of which \( f \) (strongly) scrambles the boundary. However we now specify the types of fixed points we are looking for. Thus, the main result of this subsection proves the existence of specific fixed points in non-degenerate continua satisfying the appropriate boundary conditions and shows that in some cases such continua must be degenerate. It is in this form that we apply the result later in this subsection.

Recall that an essential crossing of an external ray \( R \) and a crosscut \( Q \) was defined in Definition 3.6.4; there an external-ray \( R_t \) is said to cross a crosscut \( Q \) essentially if and only if there exists \( x \in R_t \) such that the bounded component of \( R_t \setminus \{x\} \) is contained in the bounded complementary domain of \( T(X) \cup Q \). The fact that a crosscut crosses a ray essentially can be similarly restated in the language of the uniformization plane (i.e., if the ray and the crosscut are replaced by their counterparts in the uniformization plane while \( X \) is replaced by the unit disk in the uniformization plane).

For the next definition we need to make an observation. Suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a map and \( D \) is a closed Jordan disk with interior non-disjoint from a continuum \( X \) such that \( f(D \cap X) \subset X \) and \( f(\partial D \setminus X) \cap D = \emptyset \). Suppose in addition that \( |\partial D \cap X| \geq 2 \). Then the closure of any component \( Q \) of \( \partial D \setminus X \) is a bumping arc whose endpoints map back into \( X \) and such that \( f(Q) \cap Q = \emptyset \) (indeed, \( f(\partial D \setminus X) \cap D = \emptyset \) implies that \( f(Q) \cap Q = \emptyset \)). Thus, \( \text{var}(f, Q) \) is well-defined.

Definition 7.4.5. Let \( f \) be a positively oriented map and \( X \) a continuum. If \( f(p) = p \) and \( p \in \partial X \) then we say that \( f \) repels outside \( X \) at \( p \) provided there exists a ray \( R \subset \mathbb{C} \setminus X \) from \( \infty \) which lands on \( p \) and a sequence of simple closed curves \( S^j \) bounding closed disks \( D^j \) such that \( D^1 \supset D^2 \supset \ldots \), \( p \in \text{int}(D^j) \), \( \cap D^j = \{p\} \), \( f(D^j \cap X) \subset X \), \( f(S^j \setminus X) \cap D^j = \emptyset \) and for each \( j \) there exists a component \( Q^j \) of \( S^j \setminus X \) such that \( Q^j \cap R \neq \emptyset \) and \( \text{var}(f, Q^j) \neq 0 \).

Definition 7.4.5 gives some information about dynamics around \( p \).

Lemma 7.4.6. Suppose that \( f : \mathbb{C} \to \mathbb{C} \) is a positively oriented map, \( X \) a non-separating continuum, \( p \in \partial X \) such that \( f(p) = p \) and \( f \) repels outside \( X \) at \( p \). If \( R \) is the ray from the Definition 7.4.5 then \( f(Q^j) \cap R \neq 0 \). Moreover, if \( f \) scrambles the boundary of \( X \), then \( \text{var}(f, Q^j) > 0 \).
Thus, even though $R$ above may be non-invariant, there are crosscuts approaching $p$ which are mapped by $f$ "along $R$ farther away from $p$".

**Proof.** Take $z \in R \cap Q^j$ so that $(z, \infty)_R \cap Q^j = \emptyset$. Choose a junction with $[z, \infty)_R$ (the subray of $R$ running from $z$ to infinity) playing the role of $R_i$ and two other rays close to $[z, \infty)_R$ on both sides. Then $f(Q^j) \cap R = \emptyset$ implies var($f, Q^j$) = 0, a contradiction. The second claim follows by Lemma 7.3.2.

We will use the following version of uniformization. Let $X$ be a non-separating continuum and $\varphi : \mathbb{D}_\infty \to C \setminus X$ an onto conformal map such that $\varphi(\infty) = \infty$ (here $\mathbb{D}_\infty = C \setminus \mathbb{D}$ is the complement of the closed unit disk). Thus, we choose the uniformization, under which the complement $C \setminus \overline{D}$ of the closed unit disk corresponds to the complement $C \setminus X$ of $X$. Of course, the same can be considered on the two-dimensional sphere $C^\infty$ which is sometimes more convenient. Notice, that since $\overline{D}$ is a non-separating continuum in $C$, we can use for it the usual terminology (crosscuts, shadows, etc). Also recall that the shadow of a crosscut $C$ of a nonseparating continuum $X$ is the bounded component of $C \setminus [X \cup C]$ (and not its closure).

Then, given a crosscut $C$ of $X$ with endpoints $x, y$, we can associate to its endpoints external angles as follows. It is well known [Mil00] that $\varphi^{-1}(C)$ is a crosscut of the closed unit disk with endpoints $\alpha, \beta$. It follows that we can extend $\varphi$ by defining $\varphi(\alpha) = x$ and $\varphi(\beta) = y$. Note that this extension is not necessarily continuous. In this case we say that $\alpha$ corresponds to $x$ and $\beta$ corresponds to $y$. There is a unique arc $I \subset S^1$ with endpoints $\alpha, \beta$, contained in the shadow of $\varphi^{-1}(C)$. Assuming that the positive orientation on $I$ is from $\alpha$ to $\beta$, we choose the appropriate orientation of $C$ (i.e., in this case from $x$ to $y$) and call such an oriented $C$ positively oriented.

Observe that in this situation if $D$ is a disk around a point $x \in X$ then components of $\partial D \setminus X$ are crosscuts of $X$ whose $\varphi$-preimages are crosscuts of $\overline{D}$ in the uniformization plane. However, these preimage-crosscuts in $\mathbb{D}_\infty$ may be located all over $S^1$.

The next theorem is the main result of this subsection.

**Theorem 7.4.7.** Suppose that $f$ is a positively oriented map of the plane with only isolated fixed points, $X \subset C$ is a non-separating continuum or a point, and the following conditions hold.

1. For each fixed point $p \in X$ we have that $x \in \partial X$, ind($f, p$) = 1 and $f$ repels outside $X$ at $p$.

2. The map $f$ scrambles the boundary of $X$. Moreover, using the notation from Definition 5.4.1 it can be said that for each $i$ either $f(K_i) \cap Z_i = \emptyset$, or there exists a neighborhood $U_i$ of $K_i$ with $f(U_i \cap X) \subset X$.

Then $X$ is a point.

**Proof.** Suppose that $X$ is not a point. Since $f$ has isolated fixed points, there exists a simply connected neighborhood $V$ of $X$ such that all fixed points $\{p_1, \ldots, p_m\}$ of $f|_{\overline{V}}$ belong to $X$. The idea of the proof is to construct a tight bumping simple closed curve $S$ such that $X \subset T(S) \subset V$ and var($f, S$) $\geq m$. Hence ind($f, S$) = var($f, S$) + 1 $\geq m + 1$ while by Theorem 7.4.4 our assumptions imply that ind($f, S$) = $m$, a contradiction.
First we need to make a few choices of neighborhoods and constants; we assume that there are \( n \) exit continua \( K_1, \ldots, K_n \). We also assume that they are numbered so that for \( 1 \leq i \leq n_1 \) we have \( f(K_i) \cap Z_i = \emptyset \) and for \( n_1 < i \leq n \) we have \( f(K_i) \cap Z_i \neq \emptyset \). Choose for each \( i \) a small neighborhood \( U_i \) of \( K_i \) as follows.

**CHOOSING NEIGHBORHOODS \( U_i \) OF EXIT CONTINUA \( K_i \)**

1. By assumption (2) of the theorem we may assume that \( f(U_i \cap X) \subset X \) for each \( i \) with \( n_1 < i \leq n \).
2. By continuity we may assume that \( d(U_i \cup Z_i, f(U_i)) > 0 \) for \( 1 \leq i \leq n_1 \).
3. We may assume that \( T(X \cup \bigcup U_i) \subset V \) and \( U_i \cap U_k = \emptyset \) for all \( i \neq k \).
4. We may assume that every fixed point of \( f \) contained in \( U_i \) is contained in \( K_i \).

Let \( \{p_1, \ldots, p_t\} \) be all fixed points of \( f \) in \( X \setminus \bigcup_i K_i \) and let \( \{p_{t+1}, \ldots, p_m\} \) be all the fixed points contained in \( \bigcup K_i \). Observe that then by part (4) of the choice of neighborhoods \( U_i \) we have \( p_i \in X \setminus \bigcup U_i \) if \( 1 \leq i \leq t \). Also, it follows that for each \( j, t+1 \leq j \leq m \), there exists a unique \( r_j \), \( 1 \leq r_j \leq n \), such that \( p_j \in K_{r_j} \). For each fixed point \( p_j \in X \) we rely upon Definition 7.4.5 and, as specified in that definition, choose a ray \( R_j \subset \mathbb{C} \setminus X \) landing on \( p_j \). Now we choose closed disks \( D_j \) around each \( p_j \) from Definition 7.4.5 so that, in addition to properties from Definition 7.4.5 (listed below as (6), (9) and (10)), they satisfy the following conditions.

**THE CHOICE OF CLOSED DISKS \( D_j \)**

5. \( D_j \cap R_i = \emptyset \) for all \( i \neq j \).
6. \( f(S_j \setminus X) \cap D_j = \emptyset \).
7. \( T(X \cup \bigcup_j D_j) \subset V \).
Denote by $Q(j, s)$ all components of $S_j \setminus X$; then there exists a component, $Q(j, s(j))$, of $S_j \setminus X$, with $\text{var}(f, Q(j, s(j))) > 0$ and $Q(j, s(j)) \cap R_j \neq \emptyset$ (this is possible by Definition 7.4.5 and Lemma 7.4.6).

Let us use our standard uniformization $\varphi : \mathbb{D}_\infty \to \mathbb{C} \setminus X$ described before the statement of the theorem. It serves as an important tool; in particular it allows us to pull crosscuts from the plane containing $X$ to $\mathbb{D}_\infty$ and introduce the appropriate orientation on all these crosscuts.

Claim A. Suppose that $W \subset D_j$ is a Jordan disk around $p_j$ (e.g., $W$ may coincide with $D_j$) and $C$ is a crosscut which is a component of $\partial W \setminus X$. Then the shadow of $Q(i, s(i)), i \neq j$ is not contained in the shadow of $C$ (thus, the shadows of $Q(j, s(j))$ and $Q(i, s(i))$ are disjoint). Moreover, if $W$ is sufficiently small, then the shadow of $Q(j, s(j))$ is not contained in the shadow of $C$ either.

Proof of Claim A. By condition (8) from the choice of the disks $D_j$ those disks are pairwise disjoint. Hence all the crosscuts $Q(r, t)$ are pairwise disjoint, and $C \cap Q(i, s(i)) = \emptyset$. If the shadow of $Q(i, s(i))$ is contained in the shadow of $C$, then the ray $R_i$ intersects $C$, contradicting condition (5) from the choice of the disks $D_j$. Hence $\text{Sh}(Q(j, s(j))) \cap \text{Sh}(Q(i, s(i))) = \emptyset$ for $i \neq j$ (otherwise, because the crosscuts are pairwise disjoint, one of the shadows would contain the other one, impossible by the just proven). Now, if $\text{Sh}(Q(j, s(j)) \subset \text{Sh}(C)$, then, in $\mathbb{D}_\infty$, $\varphi^{-1}(C)$ shields $\varphi(Q(j, s(j))$ from infinity and must be, together with $C$, of a bounded away from zero size. Hence, if $W$ is very small, this cannot happen.

Now we define another collection of disks around the points $p_j$. By Claim A for each $j$ we choose a small Jordan disk $D'_j$ from Definition 7.4.5 around $p_j$ so that no shadow $\text{Sh}(Q(i, s(i))$ is contained in the shadow of any crosscut $C$ which is a component of $(\partial D'_j) \setminus X$. In particular, for each such $C$, $f(C) \cap C = \emptyset$. Let us now choose a few constants.

**THE CHOICE OF CONSTANTS $\eta < \delta < \varepsilon$**

Choose $\varepsilon > 0$ such that for all $x \in X \setminus \bigcup D'_j$, $d(x, f(x)) > 3\varepsilon$ and for each crosscut $C$ of $X$ of diameter less than $\varepsilon$ with at least one endpoint outside of $\bigcup D_j$ we have that $f(C)$ is disjoint from $C$ (observe that outside any given neighborhood of $\{p_1, \ldots, p_m\}$ all points of $X$ move under $f$ by a bounded away from zero distance).

Choose $\delta > 0$ so that the following several inequalities hold:

(a) $3\delta < \varepsilon$,
(b) $3\delta < d(Z_i, Z_j)$ for all $i \neq j$,
(c) $3\delta < d(Z_i, [X \cup f(X)] \setminus [Z_i \cup U_i])$ for each $i$,
(d) $3\delta < d(K_i, C \setminus U_i)$ for each $i$,
(e) if $f(K_i) \cap Z_i = \emptyset$, then $3\delta < d(f(U_i), Z_i \cup U_i)$.

By continuity choose $\eta > 0$ such that for each set $H \subset V$ of diameter less than $\eta$ we have $\text{diam}(H) + \text{diam}(f(H)) < \delta$ and that $d(D'_i, D'_j) > \eta, i \neq j$. 

(8) $[D_j \cup f(D_j)] \cap [D_k \cup f(D_k)] = \emptyset$ for all $j \neq k$.
(9) $f(D_j \cap X) \subset X$ (this is possible because $f$ repels outside $X$ at each fixed point of $f$ and by Definition 7.4.5).
(10) Denote by $Q(j, s)$ all components of $S_j \setminus X$; then there exists a component, $Q(j, s(j))$, of $S_j \setminus X$, with $\text{var}(f, Q(j, s(j))) > 0$ and $Q(j, s(j)) \cap R_j \neq \emptyset$ (this is possible by Definition 7.4.5 and Lemma 7.4.6).
(11) $[D_j \cup f(D_j)] \cap \bigcup U_t = \emptyset$ for all $1 \leq j \leq t$.
(12) If $i < j \leq m$ then $[D_j \cup f(D_j)] \subset U_{r_i}$.
By (11) and (13) above, if a set \( H \subset V \) is of diameter at most \( \eta \) and \( H \not\subset \bigcup D'_j \), then \( f(H) \cap H = \emptyset \). Indeed, otherwise let \( x \in H \setminus \bigcup D'_j \) and \( y \in H \) be such that \( f(y) \in H \). Then by (13) \( d(x, f(x)) > 3\varepsilon \) while by (15) the triangle inequality \( d(x, f(y)) + d(f(y), f(x)) < \delta < \varepsilon /3 \), a contradiction.

Consider the family \( \mathcal{E}_X \) of all components of the sets \( (\partial D'_j)^c \setminus X \), and the crosscuts \( Q(i, s(i)) \). By condition (8) from the choice of the closed disks \( D_j \), the disks \( D_j \) are pairwise disjoint; hence, the crosscuts in \( \mathcal{E}_X \) are pairwise disjoint.

Let \( T \) be the topological hull \( T = T(X \cup (\bigcup D'_j) \cup \bigcup Q(i, s(i))) \). Then \( T \) is a non-separating continuum. Call \( C \in \mathcal{E}_X \) an unshielded (crosscut of \( X \)) if it is a part of \( \partial T \) and denote the family of all such crosscuts by \( \mathcal{E}^u_X \). By Claim A all crosscuts \( Q(i, s(i)) \) are unshielded. Call \( \varphi \)-preimages of unshielded crosscuts unshielded (crosscuts of \( \partial \mathcal{D} \)) and denote their family by \( \mathcal{E}^u_{\mathcal{D}} \). Clearly, any two unshielded crosscuts have disjoint shadows.

For each \( C \in \mathcal{E}^u_X \), let \( C_\mathcal{D} = \varphi^{-1}(C) \). Note that there are at most finitely many crosscuts \( C \in \mathcal{E}^u_X \) with \( \text{diam}(C) \geq \eta /30 \). Let \( C^1, \ldots, C^q \) be the collection of all crosscuts \( Q(i, s(i)) \) and all crosscuts in \( \mathcal{E}^u_X \) with diameters at least \( \eta /30 \). By definition 7.4.5, \( f(C^j) \cap \overline{C^j} = \emptyset \) for each \( j \). Then the crosscuts \( C^j = \varphi^{-1}(C^j) \) are all pairwise disjoint and have disjoint shadows. Hence we may assume that, if the endpoints of \( C^j \) are \( \alpha_j, \beta_j \), then \( \alpha_1 < \beta_1 < \cdots < \alpha_q < \beta_q < \alpha_{q+1} = \alpha_1 \) in the positive circular order around \( S^1 \).

For each \( i, 1 \leq i \leq q \), choose a finite chain of crosscuts \( F^i_1 \) in \( \mathbb{D}_\infty \) with endpoints \( \gamma_j, \gamma_{j+1} \) where \( \beta_i = \gamma_1 < \gamma_2 < \cdots < \gamma_k = \alpha_i+1 \) so that all closures of crosscuts from the collection \( \{ C^i_1, \ldots, C^i_q \} \bigcup \{ F^i_1 \}_{i,j} \) (except for the adjacent crosscuts which share endpoints) and their shadows are pairwise disjoint (this can be easily done, e.g. because accessible points on the boundary of \( X \) are dense), \( \varphi(F^i_1) = G^i_j \) is a crosscut of \( X \) and \( \text{diam}(G^i_j) < \eta /30 \) for all \( i, j \). In addition we may assume that non-adjacent \( G^i_j \) have disjoint sets of endpoints. Let \( Y = T(\bigcup C^i_1 \bigcup F^i_1) \). Then \( Y \) is a Jordan disk whose boundary is a simple closed curve \( \overline{S''} \subset \mathbb{D}_\infty \). Let \( S'' = \varphi(\overline{S''}) \). The set \( S'' \cap X \) is finite. It partitions \( S'' \) into links which include all \( C^i \)'s. However, some links of the form \( G^i_j \) may be very close to a fixed point of \( f \) and may not move off themselves under \( f \). Hence we modify \( S'' \) as follows.

Claim B. There exists a bumping simple closed curve \( S \) such that:

1. \( \bigcup_{i \leq 1} C^j \subset S \).
2. \( \bigcup_{i \geq 1} C^j \subset S \).
3. all components of \( S \setminus (X \cup \bigcup D'_j \bigcup \bigcup C^i) \) are of diameter less than \( \eta \).
4. for each \( i \) components of \( S \cap \text{int}(D'_j) \) are so small that they stay far away from the fixed points and are moved off themselves by \( f \).

Let \( Z = T(S'' \bigcup \bigcup D'_j) \). Then \( Z \) is a Jordan disk whose boundary is a simple closed curve \( S' \), and all crosscuts \( C^i \) are still contained in \( S' \). We modify \( S' \), keeping all \( C^i \)'s but changing \( S'' \setminus \bigcup C^i \) so that the resulting bumping simple closed curve \( S \) can be partitioned into finitely many links each of which does not go deep into the interior of \( \bigcup D'_j \) and, hence, moves off itself under \( f \).

Consider the crosscuts \( \mathcal{E}^u_{\mathcal{D}} \) in \( \mathbb{D}_\infty \). If the chain \( \{ F^i_1, \ldots, F^i_{s(i)} \} \) intersects a crosscut \( Q \in \mathcal{E}^u_{\mathcal{D}} \) let \( p_Q \) and \( r_Q \) be the first and last point of intersection of the arc \( \cup_i F^i_j \) and \( Q \). Then \( p_Q \neq r_Q \). If \( \varphi([p_Q, r_Q]) \) is small then move forward along \( S'' \). Otherwise suppose that the endpoints of \( \varphi(Q) \) are \( a_Q \) and \( b_Q \) and assume that \( a_Q < b_Q \) in the positive order around \( X \). Suppose that \( p_Q \in F_j^i \) which has endpoints...
γ_j and γ_{j+1} and ρ_Q ∈ F^n which has endpoints γ_k and γ_{k+1} and γ_{j+1} ≤ γ_k. Replace the subarc from γ_j up to γ_{k+1} in S′′ by an arc joining the same endpoints whose ϕ-image is very close to ϕ(Q). Moving forward along S′′ in the positive direction and making finitely often similar modifications, we obtain the desired simple closed curve S. This completes the proof of Claim B.

We want to compute the variation of S. Each link Q(j, s(j)) contributes at least 1 towards var(f, S), and we want to show that all other links have non-negative variation. To do so we want to apply Lemma 7.3.2. Hence we need to verify that all links of S are bumping arcs whose endpoints map back into X such that their images are disjoint from themselves. By Claim B, all links of S move off themselves. However some links of S may have endpoints mapped off X. To ensure that for our bumping simple closed curve endpoints e of its links map back into X we have to replace some of the finite chains of links of S by one link which is their concatenation (this is similar to what was done in Theorem 7.3.3). Then we will have to check if the new “bigger” links still have images disjoint from themselves.

Denote by A the initial partition of S into links which are called A-links.

**Claim C.** There exists a partition ′ of S whose links are bumping arcs with endpoints mapped back into X and whose images are disjoint from themselves. Moreover, ′-links are concatenations of A-links of S such that all A-links of S of type Q(i, t) remain ′-links of S.

**Proof of Claim C.** Suppose that X ∩ S = A = {a_0 < a_1 < · · · < a_n} and a_0 ∈ A is such that f(a_0) ∈ X (by arguments similar to those in Theorem 7.3.3 one can show that we can make this assumption without loss of generality). We only need to enlarge and concatenate links of S with at least one endpoint of the link mapped outside X. Therefore all links of S of the form Q(j, s) do not come under this category of links because by Definition 7.4.5 their endpoints do map into X. Other links, however, may have endpoints mapped outside X. Observe, that by the construction all those other links are less than η in diameter and hence have the property (15) of the choice of the constants.

Let t′ be minimal such that f(a_{t′}) ∉ X and t'' > t′ be minimal such that f(a_{t''}) ∈ X. Then f(a_{t′}) ∈ Z_i for some i. Since every component of [a_{t′}, a_{t''}] \ X has image of diameter less than δ (which is less than the distance between any two sets Z_i, Z_j), f(a_t) ∈ Z_i \ X for all t′ ≤ t < t'' . On the other hand, for t' ≤ t ≤ t'' , a_t ∉ U_i. To see this, note that if f(K_i) ∩ Z_i = ∅, then by the above made choices f(U_i) ∩ Z_i = ∅ and if f(K_i) ∩ Z_i ≠ ∅, then f(U_i \ X) ⊂ X by the assumption. Thus, all points a_t, t' ≤ t < t'' are in X \ U_i while all their images f(a_t) are in Z_i \ X, which by the property (14c) of the constant δ implies that these two finite sets of points are at least 3δ distant.

Now, by the proven above, all the A-links of S in the arc [a_{t'−1}, a_{t''}] are of diameters less than η. Hence it follows from the properties (15) and (14c) of the constant δ that f([a_{t'−1}, a_{t''}]) ∩ [a_{t'−1}, a_{t''}] = ∅ and we can remove the points a_t, for t' ≤ t < t'' from the partition A of S. By continuing in the same fashion we obtain a subset ′ ⊂ A such that for the closure of each component C of S \ ′, f(C) ∩ C = ∅ and for both endpoints a and a’ of C, {f(a), f(a’)} ⊂ X. Moreover, as was observed above, the enlarging of links of S does not concern any links of the original bumping simple closed curve of the form Q(j, s), in particular for each j, Q(j, f(s)) is an ′-link of S.
Now we apply a version of the standard argument from the proof of Theorem 7.1.3; here, instead of Theorem 3.1.4 we use the fact that $f$ satisfies the argument principle. Indeed, by Theorem 3.2.2 and Lemma 7.3.2, $\text{ind}(f, S) \geq \sum \text{var}(f, Q(j, j(s))) + 1 \geq m + 1$ contradicting Theorem 7.4.4.

It is possible to use a different approach in Definition 7.4.5 and Theorem 7.4.7. Namely, a version of Definition 7.4.5 could define repelling outside $X$ at a fixed point $p$ as the existence of a family of closed disks $D^j$ with similar properties except now we would require the existence of at least $\text{ind}(f, p)$ pairwise non-homotopic rays outside $X$ from $\infty$ to $p$ (landing on $p$) and the existence of the same number of components of $S^j \setminus X$ non-disjoint from corresponding rays and with non-zero variation on each such component.

Then a version of Theorem 7.4.7 would state that if a positively oriented map with isolated fixed points $f$ repels outside $X$ at each of its fixed points, and the condition (2h) of Theorem 7.4.7 is satisfied, then $X$ must be a point. The proof of this version of Theorem 7.4.7 is almost the same, except for a bit heavier notation needed (now that we have not one, but $\text{ind}(f, p)$ crosscuts with positive variation around each fixed point in $X$). Since for our applications Theorem 7.4.7 suffices we restricted ourselves to this case.

Theorem 7.4.7 implies the following.

**Corollary 7.4.8.** Suppose that $f$ is a positively oriented map with isolated fixed points, and $X \subset \mathbb{C}$ is a non-separating and non-degenerate continuum satisfying condition (2h) stated in Theorem 7.4.7 and such that all fixed points in $X$ belong to $\partial X$. Then either $f$ does not repel outside $X$ at one of its fixed points, or the local index at one of its fixed points is not equal to 1.

Lemma 7.4.9 gives a verifiable sufficient condition for a fixed point $a$ belonging to a locally invariant continuum $X$ to be such that the map $f$ repels outside $X$.

We will apply the lemma in the next section.

**Lemma 7.4.9.** Suppose that $f$ is a positively oriented map, $X \subset \mathbb{C}$ is a non-separating continuum or a point and $p \in \partial X$ is a fixed point of $f$ such that:

(i) there exists a neighborhood $U$ of $p$ such that $f|_U$ is one-to-one and $f(U \cap X) \subset X$,

(ii) there exists a ray $R \subset \mathbb{C}^\infty \setminus X$ from infinity such that $\mathcal{R} = R \cup \{p\}$, $f|_R : R \to R$ is a homeomorphism and for each $x \in R$, $f(x)$ separates $x$ from $\infty$ in $R$,

(iii) there exists a nested sequence of closed disks $D_j \subset U$ with boundaries $S_j$ containing $p$ in their interiors such that $\mathcal{R} \cap \{p\}$ and $f(S_j \setminus \mathcal{R}) \cap D_j = \emptyset$.

Then for a sufficiently large $j$ there exists a component $C$ of $S_j \setminus X$ with $C \cap R \neq \emptyset$ and $\text{var}(f, C) > 0$, so that $f$ repels outside $X$ at $p$.

Observe that here we show that $\text{var}(f, C) > 0$ without any “scrambling” assumptions on $f$.

**Proof.** Choose a Jordan disk $U$ as in (i) so that $(\partial U) \cap R = \{q\}$ is a point and $X \setminus U \neq \emptyset$. Choose $j$ so that $D_j \cup f(D_j) \subset U$. By [BO06] there is a component $C$ of $(\partial D_j) \setminus X$ such that $R$ crosses $C$ essentially (see Definition 3.6.4). Slightly adjusting $D_j$, we may assume that $R \cap \partial D_j$ is finite and each intersection is transversal. Since
$R$ crosses $C$ essentially, $|R \cap C|$ is odd; since $f$ is one-to-one on $\overline{U}$, $|f(C) \cap R|$ is odd as well.

Let $u,v$ be the endpoints of $C$. Observe, that $C$ can be included in a simple closed curve $S$ around $X$ so that $X \subset T(S)$. Since by (i) $f(\overline{U} \cap X) \subset X$, we see that $f(u) \in T(S), f(v) \in T(S)$ and variation $\text{var}(f, C) > 0$ is well-defined (see Definition 2.2.2).

Let us move along $R$ from infinity to $p$ and denote by $w$ be the first point of $C \cap R$ which we meet and by $z$ the last point. Then by (ii), $f(C) \cap [p, z]_R = \emptyset$. Also, by (i) $|f(C) \cap R| = |f(C \cap R)|$ is odd. Since $X \setminus D_j \neq \emptyset$ and $X \cap ([z, w]_C \cup [z, w]_R) = \emptyset$, $X$ is contained in the unbounded component $V_{\infty}$ of the complement to $[z, w]_C \cup [z, w]_R$.

Since $u$ and $v$ belong to $U \cap X$, their images $f(u)$ and $f(v)$ belong to $X$ as well. Hence $f(u), f(v) \in V_{\infty}$.

**Claim A.** $|f(C) \cap [z, w]_R|$ is even.

**Proof of Claim A.** A complementary domain $O$ of $[z, w]_C \cup [z, w]_R$ is called even/odd if there is an arc $J$ from infinity to a point in $O$ so that $J \cap C \cap R = \emptyset$, $J \cap ([z, w]_C \cup [z, w]_R)$ is finite, every intersection is transversal and $|J \cap ([z, w]_C \cup [z, w]_R)|$ is even/odd, respectively. By [OT82] the notion of an even/odd domain is independent of $J$, well-defined and each complementary domain of $[z, w]_C \cup [z, w]_R$ is either even or odd. Since by (iii) $f(\overline{C}) \cap \overline{C} = \emptyset$, $f(\overline{C})$ can only intersect $C \cup R$ at points of $R$. Also, whenever $f(\overline{C})$ meets $[z, w]_R$, it crosses from an even to an odd domain or vice versa. Since both $f(u)$ and $f(v)$ are in the unbounded (and hence even) domain of $C \setminus [f(C) \cup [z, w]_R]$, $|f(C) \cap [z, w]_R|$ is even as desired. \square

Observe that $f(\overline{C})$ is outside $D_j$ and hence is disjoint from $[p, z]_R \subset D_j$. Since $|f(C) \cap R|$ is odd and $|f(C) \cap [p, w]_R| = |f(C) \cap [z, w]_R|$ is even, $|f(C) \cap [w, \infty]_R|$ is odd. Since every intersection is transversal, we can replace $[w, \infty]_R$ by a junction.
For each fixed point \( p \) suppose that \( \sigma \). Recall that \( \Sigma \). Hence points \( u \) and \( v \) can be connected with an arc \( K \) inside \( W \) disjoint from \( R \cup \Sigma \). Then, since \( f_U \) is a homeomorphism, \( f(K) \cap R = \emptyset \). By Lemma 7.3.1 this implies that \( \var(f, C) \geq 0 \) and by the previous paragraph then \( \var(f, C) > 0 \) (basically, we simply choose a junction \( J \) close to \( [w, \infty)_R \) such that \( f(K) \cup J' = \emptyset \) and conclude that \( \var(f, C) = \operatorname{win}(f, C \cup K, w) > 0 \) since \( f \) is a positively oriented map). \( \square \)

It is now easy to see that the following corollary holds.

**Corollary 7.4.10.** Suppose that \( X \subset C \) is a non-separating continuum or a point and \( f : C \to C \) is a positively oriented map with isolated fixed points, and the following conditions hold.

(a) Each fixed point \( p \in X \) is topologically repelling, belongs to \( \partial X \), and has a neighborhood \( U_p \) such that \( f(U_p \cap X) \subset X \) and \( f|_{U_p} \) is a homeomorphism.

(b) For each fixed point \( p \in X \), there exists a ray \( R \subset C^\infty \setminus X \) from infinity landing on \( p \). \( f|_R : R \to R \) is a homeomorphism and for each \( x \in R \), \( f(x) \) separates \( x \) from \( \infty \) in \( R \).

(c) The map \( f \) scrambles the boundary of \( X \). Moreover for every \( i \) either \( f(K_i) \cap Z_i = \emptyset \) or there exists a neighborhood \( U_i \) of \( K_i \) with \( f(U_i \cap X) \subset X \).

Then \( X \) is a (fixed) point.

**Proof.** Let us apply Theorem 7.4.7. To do so, we verify its conditions. The facts that \( X \subset C \) is a non-separating continuum or a point and \( f : C \to C \) is a positively oriented map with isolated fixed points are clearly satisfied. To verify condition (1h) of Theorem 7.4.7, suppose that \( p \in X \) is a fixed point. Then by (a) above \( p \in \partial X \). Moreover, \( p \) is topologically repelling, and so by Lemma 7.4.3 the index at \( p \) is +1.

It remains to verify that \( f \) repels outside \( X \) at \( p \). To do so we apply Lemma 7.4.9. Since \( p \) is topologically repelling, there exists a nested sequence of closed disks \( D_j \subset U \) with boundaries \( \Sigma_j \) containing \( p \) in their interiors, with \( \bigcap D_j = \{ p \} \) and \( f(S_j \setminus X) \cap D_{j+k} = \emptyset \). Hence the condition (iii) of Lemma 7.4.9 is satisfied. The condition (i) of Lemma 7.4.9 immediately follows from (a) above; the condition (ii) of Lemma 7.4.9 immediately follows from (b) above. Hence by Lemma 7.4.9 the map \( f \) repels outside \( X \) at \( p \). Therefore the condition (1h) of Theorem 7.4.7 is satisfied. Condition (2h) of Theorem 7.4.7 is also satisfied (it simply coincides with condition (e) of our corollary), hence by Theorem 7.4.7 \( X \) is a point. \( \square \)

### 7.5. Applications to complex dynamics

We begin by introducing a few facts concerning local dynamics at parabolic and repelling periodic points of a polynomial which were not necessary for stating the results of this section in Chapter 5 but are needed for the proofs. A nice description of this can be found in [Mil00] ([CG93] can also serve as a good source here).

Let \( P : C \to C \) be a complex polynomial, \( J_P \) its Julia set (\( J_P \) is the closure of the set of repelling periodic points of \( P \)) and \( K_P = T(J_P) \) the “filled-in” Julia set. Recall that \( \sigma_d : S^1 \to S^1 \) is defined by \( \sigma_d(\alpha) = d\alpha \mod 1 \), where \( S^1 = \mathbb{R}/\mathbb{Z} \).
is parameterized by \([0,1]\). (This map is conjugate to the map \(z \to z^d\) restricted to the unit circle in the complex plane.) If \(p\) is a periodic point of \(P\) of period \(n\) and \((P^n)'(p) = re^{2\pi i \alpha}\) with \(r \geq 0\), then \(p\) is repelling if \(r > 1\), parabolic if \(r = 1\) and \(\alpha \in \mathbb{Q}\), irrational neutral if \(r = 1\) and \(\alpha \in \mathbb{R} \setminus \mathbb{Q}\) and attracting if \(r < 1\). If \(p\) is a repelling or parabolic fixed point in a non-degenerate component \(Y\) of \(K_P\), then by [DH85a, LP96] there exist \(1 \leq k < \infty\) external rays \(R_{\alpha(i)}\) such that \(\sigma_d|_{\alpha(1),\ldots,\alpha(k)} : \{\alpha(1),\ldots,\alpha(k)\} \to \{\alpha(1),\ldots,\alpha(k)\}\) is a permutation, all \(\alpha(i)\) are of the same minimal period under \(\sigma_d\), for each \(j\) the ray \(R_{\alpha(j)}\) lands on \(p\), and no other external rays land on \(p\).

Components of \(\mathbb{C} \setminus J_P\) are called Fatou domains. There are three types of bounded Fatou domains \(U\). A Fatou domain \(U\) is called an attracting domain if it contains an attracting periodic point, a Siegel domain if it contains an irrational neutral periodic point and a parabolic domain if it is periodic but contains no periodic points. In the latter case there always exists a parabolic periodic point on the boundary of the parabolic Fatou domain. An irrational neutral periodic point inside a Siegel domain is called a Siegel (periodic) point; an irrational neutral periodic point in \(J_P\) is called a Cremer (periodic) point.

Any two distinct parabolic Fatou domains which contain the same parabolic periodic point in their boundaries are separated by two external rays which land at this parabolic point. It is also known that points inside these parabolic domains are attracted by the orbit of this parabolic periodic point while points on the external rays landing at points of this orbit are repelled to infinity. Suppose that \(p\) is a parabolic fixed point, \(P'(p) = e^{2\pi i \xi} r \in \mathbb{Z}, R_0,\ldots,R_{m-1}\) are all external rays landing at \(p\) and \(U_0,\ldots,U_{k-1}\) are all Fatou domains which contain \(p\) in their boundaries. Moreover, suppose that both rays and domains are numbered according to the positive circular order around \(p\). Then combinatorially one can think of the local action of \(P\) on rays and domains at \(p\) as a rotation by \(r/q\). This means that \(P(U_j) = U_{(j+r) \mod k}\) and all rays between \(U_i\) and \(U_{i+1}\) are mapped by \(P\) in an order preserving way onto all rays between \(U_{(i+r) \mod k}\) and \(U_{(i+1+r) \mod k}\).

Before we continue we want to recall the notion of a general puzzle-piece which is first introduced in Definition 5.5.1.

**Definition 5.5.1 (General puzzle-piece).** Let \(P : \mathbb{C} \to \mathbb{C}\) be a polynomial. Let \(X \subset K_P\) be a non-separating subcontinuum or a point such that the following holds.

1. There exists \(m \geq 0\) and \(m\) pairwise disjoint non-separating continua/points \(E_1 \subset X,\ldots,E_m \subset X\).
2. There exist \(m\) finite sets of external rays \(A_1 = \{R_{a_1},\ldots,R_{a_1}\},\ldots,A_m = \{R_{a_m},\ldots,R_{a_m}\}\) with \(i_k \geq 2,1 \leq k \leq m\).
3. We have \(\Pi(A_j) \subset E_j\) (so the set \(E_j \cup (\bigcup_{k=1}^{i_k} R_{a_k}) = E'_j\) is closed and connected).
4. \(X\) intersects a unique component \(C_X\) of \(\mathbb{C} \setminus \bigcup E'_j\).
5. For each Fatou domain \(U\) either \(U \cap X = \emptyset\) or \(U \subset X\).

We call such \(X\) with the continua \(E_i\) and the external rays \(R_{a_k}\) a general puzzle-piece and call the continua \(E_i\) exit continua of \(X\). For each \(k\), the set \(E'_k\) divides the plane into \(i_k\) open sets which we will call wedges (at \(E_k\)); denote by \(W_k\) the wedge which contains \(X \setminus E_k\) (it is well-defined by (4) above).
Let us now see whether condition (1) of Theorem 7.4.7 applies to the polynomial $P$ with fixed parabolic points in a general puzzle-piece $X$ such that $P(X) \cap C_X \subset X$ (loosely, it means that $P(X)$ does not “grow” at points of $C_X$). We need to check that at a fixed parabolic point $p \in X$ the map $P$ repels outside $X$ and the local index is 1.

For convenience, consider $p \in C_X$. Let $L$ be a cycle of parabolic domains containing $p$ in their boundaries. To make a more complete picture, we first observe that either $L \subset X$ or all domains in $L$ are disjoint from $X$. Indeed, suppose that one of the domains in $L$ is contained in $X$. Then, as there is no “growth” of $P(X)$ at $p$, we see that all domains of $L$ are in $X$. Otherwise by condition (5) from Definition 5.5.1 all domains in $L$ are disjoint from $X$.

Suppose first that $P'(p) \neq 1$. There must exist small disks $D$ such that components of $(\partial D) \setminus X$ map outside $D$. Let $Q$ be a crosscut which is a component of $(\partial D) \setminus X$. Choose an external ray $R$ landing at $p$ and crossing $Q$ essentially. By the analysis of dynamics around $p$ (the fact that $P$ locally “rotates”, see, e.g., [Mil00]) and since $P(X) \cap C_X \subset X$, we conclude that $P(Q) \cap R = \emptyset$. However, then $\text{var}(f,Q) = 0$, a contradiction with the existence of a crosscut of positive variation in $(\partial D) \setminus X$.

On the other hand, if $P'(p) = 1$ then, as was explained right after the proof of Lemma 7.4.3, $\text{ind}(P,p) > 1$ (again, see, e.g., [Mil00]) which contradicts the
condition (1) of Theorem 7.4.7 as well. Hence in the parabolic case Theorem 7.4.7 cannot be applied “as is” to the polynomial \( P \). Moreover, since at parabolic points the map is not topologically repelling, neither is Corollary 7.4.10 applicable in this case.

The idea allowing us to still deal with parabolic points is that we can change \( P \) inside the parabolic domains in question without compromising the rest of the arguments and modifying these parabolic points to topologically repelling periodic points. The thus constructed new map \( g \) will no longer be holomorphic but will satisfy the conditions of Corollary 7.4.10. We formalize this idea in the following lemma.

**Lemma 7.5.1.** Suppose that \( \{p_j\}, 1 \leq j \leq m \), are all parabolic fixed points of a polynomial \( P \) with \( P'(p_j) = 1 \). Then there exists a positively oriented map \( g_P = g \) which coincides with \( P \) outside the invariant parabolic Fatou domains, is locally one-to-one at each \( p_j \) (hence \( p_j \) is not a critical point of \( g \)) and is such that all the points \( p_j \) are topologically repelling fixed points of \( g \). In particular, \( \text{ind}(g, p_j) = +1 \) for all \( j \), \( 1 \leq j \leq m \).

**Proof.** Let us consider a fixed parabolic point \( p = p_j \). Let \( F_i \) be the invariant Fatou domains containing \( p \) in their boundaries \( B_i \). By a nice recent result of Yin and Roesch [RY08], the boundary \( B_i \) of each \( F_i \) is a simple closed curve and \( P|_{B_i} \) is conjugate to the map \( z \rightarrow z^{d(i)} \) for some integer \( d(i) \geq 2 \). Let \( \psi : F_i \rightarrow \mathbb{D} \) be a conformal isomorphism. Since \( B_i \) is a simple closed curve, \( \psi \) extends to a homeomorphism on \( \mathbb{D} \). Since \( f|_{B_i} \) is conjugate to the map \( z \rightarrow z^{d(i)} \), it now follows that the map \( P|_{\mathbb{D}} \) can be replaced by a map topologically conjugate by \( \psi \) to the map \( g_i(z) = z^{d(i)} \) on the closed unit disk \( \mathbb{D} \) which agrees with \( f \) on \( B_i \). Let \( g \) be the map defined by \( g(z) = P(z) \) for each \( z \in \mathbb{C} \setminus \bigcup F_i \) and \( g(z) = g_i(z) \) when \( z \in F_i \). Then \( g \) is clearly a positively oriented map.

Since by the analysis of the dynamics around parabolic points [Mil00] \( P \) repels points away from \( p \) outside parabolic domains \( F_i \), we conclude, by the construction, that \( p \) is a topologically repelling fixed point of \( g \). Clearly, \( \text{deg}_g(p) = 1 \), hence by Lemma 7.4.3 \( \text{ind}(g, p) = \text{deg}_g(p) = 1 \) as desired. Continuing in this fashion, we can change \( P \) on all invariant parabolic domains with fixed points \( p_j \), \( 1 \leq j \leq m \), in their boundaries to a map \( g \) which satisfies the requirements of the lemma. \( \square \)

We use Lemma 7.5.1 in the proof of the Theorem 7.5.2. Recall, that a fixed point \( x \) of a polynomial \( P \) is said to be non-rotational if there is a fixed external ray landing at \( x \) (it follows that each such point is either repelling or parabolic).

**Theorem 7.5.2.** Let \( P \) be a polynomial with filled-in Julia set \( K_P \) and let \( Y \) be a non-degenerate periodic component of \( K_P \) such that \( P^p(Y) = Y \). Suppose that \( X \subset Y \) is a non-degenerate general puzzle-piece with \( m \geq 0 \) exit continua \( E_1, \ldots, E_m \) such that \( P^p(X) \cap C_X \subset X \) and either \( P^p(E_i) \subset W_i \), or \( E_i \) is a \( P^p \)-fixed point. Then at least one of the following claims holds:

1. \( X \) contains a \( P^p \)-invariant parabolic domain,
2. \( X \) contains a \( P^p \)-fixed point which is neither repelling nor parabolic, or
3. \( X \) has an external ray \( R \) landing at a repelling or parabolic \( P^p \)-fixed point such that \( P^p(R) \cap R = \emptyset \) (i.e., \( P^p \) locally rotates at some parabolic or repelling \( P^p \)-fixed point).
Equivalently, suppose that $Y$ is a non-degenerate periodic component of $K_P$ such that $P^n(Y) = Y$, $X \subset Y$ is a general puzzle-piece with $m \geq 0$ exit continua $E_1, \ldots, E_m$ such that $P^n(X) \cap C_X \subset X$ and either $P^n(E_i) \subset W_i$ or $E_i$ is a $P^n$-fixed point; if, moreover, $X$ contains only non-rotational $P^n$-fixed points and does not contain $P^n$-invariant parabolic domains, then it is degenerate.

**Proof.** We may assume that $p = 1$ and $P(Y) = Y$. We show that if none of the conclusions (1)-(3) hold, then Corollary 7.4.10 applies and therefore $X$ must be a point, contradicting the assumption that $X$ is non-degenerate. However, if $X$ contains parabolic points, we first use Lemma 7.5.1 and replace $X$ by a continuum which “grows” out of $X = X \subset \mathbb{C}$.

For $y \in X$, $g$ is a repelling or parabolic fixed point of $P$ (this is because we assume that claim (2) of Theorem 7.5.2 does not hold and hence all $P^n$-fixed points in $X$ are repelling or parabolic). By Lemma 7.5.1 this implies that $y$ is a topologically repelling fixed point of $g$. Let us show that there exists a small neighborhood $U$ of $y$ such that $g(U \cap X) \subset X$. This is clear if $y \in C_X$ because $g(X) \cap C_X \subset X$ by our assumptions. Assume now that $y \notin C_X$ which means that $\{y\} = E_k$ is one of the exit-continua of $X$. Observe that there a few fixed external rays of $P$ landing at $y$ (the rays are fixed because we assume that the conclusion (3) of Theorem 7.5.2 does not hold and hence all rays which land at $y$ must be fixed). Choose the two rays $R_1, R_2$ which land at $y$ and form the boundary of the wedge $W_k$ at $y$ which contains $X$. Since $g(X) \cap C_X \subset X$ by our assumptions, this implies that in a small neighborhood $U_k$ of $E_k$ the intersection $U_k \cap X$ maps (by $g$ or $P$) into $X$ as desired. This completes the verification of condition (a) of Corollary 7.4.10.

Condition (b) of Corollary 7.4.10 (i.e., the existence of a fixed external ray landing at each fixed point in $X$) follows immediately from the assumption that claim (3) of Theorem 7.5.2 above does not hold.

Let us now check condition (c) of Corollary 7.4.10. Set $g(X) \setminus X = P(X) \setminus X = H$. We may assume that $H \neq \emptyset$ and we can think of $g(X) = P(X)$ as a continuum which “grows” out of $X$. In particular, $m \geq 1$. Fix $k$, $1 \leq k \leq m$. Since $g(X) \cap C_X \subset X$, any component of $H$ whose closure intersects $E_k$ must be contained in one of the wedges at $E_k$ (such wedges are defined in Definition 7.4.10), but not in $W_k$. Let $Z_k$ be the topological hull of the union of all components of $H$ which meet $E_k$ together with $E_k$. Then $Z_k$ is a non-separating continuum. Since either $g(E_i) \subset W_i$ or $E_i$ is a fixed point, the map $g$ scrambles the boundary of $X$ (see Definition 5.4.1). Moreover, if $E_k$ is mapped into $W_k$ then clearly $g(E_k) \cap Z_k = \emptyset$ (because $Z_k \setminus E_k$ is contained in the other wedges at $E_k$ and is disjoint from $W_k$). On the other hand, consider a fixed point $y \in X$ such that $E_k = \{y\}$. Then condition (c) of Corollary 7.4.10 follows from (a) which has already been verified. Hence, by Corollary 7.4.10 we conclude, that $X$ is a point, a contradiction.

Notice that if $X$ is a general puzzle-piece with $m = 0$ in Theorem 7.5.2, then $C_X = \mathbb{C}$. Hence in this case $P(X) \cap C_X \subset X$ implies $P(X) \subset X$ and $X$ is invariant. Thus, a non-separating invariant continuum $X \subset K_P$ is a general puzzle-piece if and only if for every Fatou domain $U$ of $P$ either $U \cap X = \emptyset$, or $U \subset X$.

The proof of the next corollary is left to the reader.
Suppose that $\text{Sh}(\alpha)$ of an angle in the connected case can be defined as the intersection of the closures of all shadows $\text{Sh}(C)$ of all crosscuts $C$ such that $R_\alpha$ crosses $C$ essentially. Corollary 7.5.4 is proved for a somewhat larger class of subcontinua of $J_P$ which includes impressions as an important particular case.

Consider a repelling or parabolic periodic or preperiodic point $x$ and all external rays landing at $x$. Then the union of two such rays and $x$ is said to be a legal cut of the plane. Also, suppose that $x, y \in \partial U$ are two periodic or preperiodic points in the boundary of an attracting or parabolic Fatou domain $U$. By [RY08] there exists an arc $A \subset U$ connecting $x$ and $y$. The union of $A$ and two external rays landing at $x$ and $y$ is also a legal cut of the plane. Finally, call a continuum $Q$ periodic if for some $n > 0$ we have $P^n(Q) \subset Q$.

**Corollary 7.5.4.** Let $P : \mathbb{C} \to \mathbb{C}$ be a complex polynomial and $Q \subset J_P$ be a periodic continuum such that for every legal cut $C$ the set $Q \setminus C$ is contained in one component of $\mathbb{C} \setminus C$. Suppose that $T(Q)$ contains no Siegel or Cremer points. Then $Q$ is degenerate. In particular, if $J_P$ is connected and $Q$ is a periodic impression such that $T(Q)$ contains no Cremer or Siegel points, then $Q$ is a point.

**Proof.** By considering an appropriate power of $P$ we may assume that $Q$ is invariant and non-degenerate. Clearly this implies that $T(Q)$ is invariant too. Suppose that $p' \in Q$ is a fixed point of $P$ and $R_{p'}$ is an external ray landing at $p'$. Then $P(R_{p'})$ also lands on $p'$. If $R_{p'}$ is not fixed, then $C = R_{p'} \cup P(R_{p'})$ is a legal cut. The local dynamics at $p'$ and the fact that $Q$ is invariant imply now that $Q$ has points on either side of $C$, a contradiction with the assumptions on $Q$. Hence each fixed repelling or parabolic point in $Q$ is non-rotational.

Let us show that $Q$ can only intersect the closure of a parabolic or attracting Fatou domain $U$ at one point. Indeed, suppose otherwise and let $x, y \in \partial U \cap Q, x \neq y$. Then there exists an arc $I \subset \partial U$ with endpoints $x, y$, contained in $Q$ because otherwise $Q$ will “shield” some points of $\partial U$ from infinity contradicting the fact that all points of $\partial U$ belong to the closure of the basin of attraction of infinity. By [RY08] we can find, say, periodic points $u, v \in I$ and include them in a legal cut $T$ which will separate some points of $I$ (and hence of $Q$) from other points of $I$, contradiction with our assumptions. Hence $T(Q)$ cannot contain an attracting or parabolic domain $U$ since otherwise, by the above, $Q$ must shield part of $\partial U$ from the basin of attraction of infinity, a contradiction. This implies that $T(Q)$ cannot contain parabolic domains or attracting points. By the assumption $T(Q)$ does not contain Cremer or Siegel points either. Hence by Corollary 7.5.3 $Q$ is a point as desired.

\[ \square \]
In the particular case in the end of the statement we assume that \( J_P \) is connected; the same result in fact holds for all Julia sets but will require the introduction of the notion of the impression for disconnected Julia sets which we avoid here for the sake of simplicity. The verification of the fact that impressions satisfy the conditions of the corollary is straightforward and therefore is left to the reader. In particular, suppose \( R_\alpha \) is a periodic external ray and the topological hull \( T(\text{Imp}(\alpha)) \) of the impression of \( \alpha \) contains only repelling or parabolic periodic points. Then, by Corollary 7.5.4, \( \text{Imp}(\alpha) \) is degenerate.

Note that the assumptions of Corollary 7.5.4 are equivalent to the following. Suppose that \( Q \subset J_P \) is a periodic continuum such that for every legal cut \( C \) the set \( Q \setminus C \) is contained in one component of \( C \setminus C \). As in the proof of Corollary 7.5.4, this implies that \( Q \) can only intersect the boundaries of attracting or parabolic domains at no more than one point. To make the conclusion of the corollary, we need to check that \( T(Q) \) contains no Siegel or Cremer points. We claim that this is equivalent to the following:

1. \( Q \) contains no Cremer point;
2. if the boundary of a Siegel disk is decomposable, then \( Q \) is disjoint from it;
3. if the boundary of a Siegel disk is indecomposable (it is not known if such Siegel disks exist), then \( Q \) intersects it in at most one point.

Indeed if (1) - (3) above are satisfied then, by an argument similar to the proof of Corollary 7.5.4, \( T(Q) \) contains no Cremer or Siegel points. Now, suppose that \( T(Q) \) contains no Cremer or Siegel points. Then by Corollary 7.5.4 \( Q = \{q\} \) is a point. Hence (1) and (3) hold trivially. If \( B \) is the decomposable boundary of a Siegel disk and \( q \in B \), then we may assume that \( B \) and \( Q \) are invariant. It is known [Rog92a, Rog92b] that there exists a monotone map \( p : B \rightarrow S^1 \) and an induced map \( g : S^1 \rightarrow S^1 \) which is an irrational rotation. Since \( Q \) is invariant, \( Q = B \), a contradiction. Hence (2) holds as well.
7.5. APPLICATIONS TO COMPLEX DYNAMICS
Bibliography


Index

accessible point, 29
acyclic, 16
allowable partition, 28
B, 35
B^\infty, 35
boundary scrambling
  for dendrites, 48
  for planar continua, 49
branched covering map, 66
branchpoint of f, 73
bumping
  arc, 25
  simple closed curve, 25
C(a, b), 38
Carathéodory Loop, 27
chain of crosscuts, 28
  equivalent, 28
channel, 29
dense, 29, 55
completing a bumping arc, 25
C, 1
C^\infty, 1
canonical
  external ray, 51
continuum
  decomposable, 47
  indecomposable, 47
conv\mathcal{E}(K), 35
conv\mathcal{N}(B \cap K), 35
convex hull
  Euclidean, 35
  hyperbolic, 35
clockwise order
  on an arc in a simple closed curve, 13
critical point, 72
crosscut, 25
shadow, 25
cutpoint, 49
C_X, 51, 83
D, 38
defines variation near X, 48
degree, 13
deg_f(a), 73
degree(g), 13
degree(f_p), 16
\partial boundary operator, 2
D^\infty, 27
dendrite, 6
domain
  attracting, 50
  Fatou, 50
  parabolic, 50
  Siegel, 50
embedding
  orientation preserving, 19
essential crossing, 29
\mathcal{E}_1, 28
exit continuum, 49
for a general puzzle-piece, 51
external ray, 29
  end of, 29
  essential crossing, 29
  landing point, 29
  non smooth, 51
  smooth, 51
fixed point
  for positively oriented maps, 61
  non-rotational, 51
(f, X, \eta), 48
G, 38
g, 38
\mathfrak{g}, 35
\Gamma, 38
gap, 37, 66
general puzzle-piece, 51
great circle
  hyperbolic, 38
geometric outchannel
  negative, 48
  positive, 48
hull
  hyperbolic, 37
Index

- topological, 1
- hyperbolic
  - geodesic, 35
  - halfplane, 35
  - hyperbolic, 38
- id identity map, 13
- $\text{Im}(E_t)$, 29
- impression, 29
- index, 13
  - fractional, 13
- $I=V+1$ for Carathéodory Loops, 28
- Index=Variation+1 Theorem, 20
- local, 50
  - $\text{ind}(f, x)$, 50
  - $\text{ind}(f, A)$, 19
  - $\text{ind}(f, g)$, 13
  - $\text{ind}(f, [a,b])$, 13
  - $\text{ind}(f, S)$, 19
- $J_P$, 50
- Julia set, 50
- junction, 14
- Jørgensen Lemma, 39
- $K_P$, 50
- $K$, 37
- $K^\ast$, 41
- $K^\ast$-chord, 37
- $K^\ast\ast$, 43
- $K^\ast P$, 37
- $K^\ast P^\ast$, 41
- Kulkarni-Pinkall Lemma, 36
  - Partition, 37
- lamination, 66
  - degenerate, 66
  - invariant, 66
- landing point, 29
- leaf, 66
- link, 25
- local degree, 73
- local index, 50
- Lollipop Lemma, 23
- map, 13
  - confluent, 16
  - light, 16
  - monotone, 16
  - negatively oriented, 16
  - on circle of prime ends, 32
  - oriented, 16
  - perfect, 16
  - positively oriented, 16
  - maximal ball, 35
- Maximum Modulus Theorem, 31
- monotone-light decomposition of a map, 16
- narrow strip, 56
- natural retraction of dendrites, 62
- non-separating, $\xi$
- order
  - on subarc of simple closed curve, 13
  - orientation preserving
    - embedding, 19
  - outchannel, 48
  - geometric, 48
  - uniqueness, 57
- periodic point
  - Cremer, 50
  - parabolic, 50
  - repelling, 50
  - Siegel, 50
  - weakly repelling, 49
  - point of period two
    - for oriented maps, 62
  - positively oriented arc, 75
  - Pr$(E_1)$, 29
- prime end, 28
  - channel, 29
  - impression, 29
  - principal continuum, 29
  - principal continuum, 29
- $R$, 1
  - repels outside $X$ at $p$, 74
  - $R_t$, 29
- shadow, 25
  - Sh$(A)$, 25
  - smallest ball, 35
  - standing hypothesis, 48
- topological
  - Julia set, 66
  - polynomial, 66
  - topological hull, 1
  - topologically
    - attracting, 73
    - repelling, 73
  - $T(X)$, 1
  - $T(X)_a$, 43
- $U^\infty$, 1
- unlinked, 66
- val$y(x)$, 49
- valence, 49
  - $\text{var}(f, A)$, 27
- variation
  - for crosscuts, 25
  - of a simple closed curve, 15
- on an arc, 14
  - on finite union of arcs, 15
  - $\text{var}(f, A, S)$, 14
  - $\text{var}(f, S)$, 15
weakly repelling, 49, 63
wedge (at an exit continuum), 51
\( \text{win}(g, S^1, u) \), 13