CONFLUENCE OF GEODESICS IN LIOUVILLE QUANTUM GRAVITY FOR $\gamma \in (0, 2)$

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We prove that for any metric, which one can associate with a Liouville quantum gravity (LQG) surface for $\gamma \in (0,2)$ satisfying certain natural axioms, its geodesics exhibit the following confluence property. For any fixed point z, a.s. any two γ -LQG geodesics started from distinct points other than z must merge into each other and subsequently coincide until they reach z. This is analogous to the confluence of geodesics property for the Brownian map proven by Le Gall ($Acta \ Math. \ 205 \ (2010) \ 287-360$). Our results apply for the subsequential limits of Liouville first passage percolation and are an important input in the proof of the existence and uniqueness of the LQG metric for all $\gamma \in (0,2)$.

1. Introduction.

1.1. Overview. Fix $\gamma \in (0, 2)$, let $U \subset \mathbb{C}$ and let h be a variant of the Gaussian free field (GFF) on U. The theory of Liouville quantum gravity (LQG) is concerned with the random Riemannian metric

$$(1) e^{\gamma h}(dx^2 + dy^2),$$

on U, where $dx^2 + dy^2$ denotes the Euclidean Riemannian metric tensor. The surface parameterized by this Riemannian metric is called a γ -LQG surface.

LQG was first introduced in the physics literate by Polyakov [24] in the context of bosonic string theory. One reason why LQG surfaces are interesting mathematically is that they arise as the scaling limits of random planar maps. The special case when $\gamma = \sqrt{8/3}$ (called "pure gravity") corresponds to the scaling limit of uniform random planar maps, and other values of γ (sometimes referred to as "gravity coupled to matter") correspond to random planar maps weighted by the partition function of an appropriate critical statistical mechanics model on the map.

The definition of γ -LQG given in (1) does not make literal sense since the GFF is only a distribution, not a function. In particular, the GFF can be integrated against a smooth test function, but it does not have well-defined pointwise values so it cannot be exponentiated. Consequently, one needs a regularization procedure to make rigorous sense of this object. Previously, this has been accomplished for the associated volume form, that is, the γ -LQG area measure. This is a random measure μ_h on U which is a limit of regularized versions of $e^{\gamma h} d^2 z$, where $d^2 z$ denotes Lebesgue measure [6, 11, 25].

It is a long-standing open problem to construct a canonical *metric* associated with a γ -LQG surface, that is, a random metric D_h on U which is obtained, in some sense, by exponentiating the GFF h. The random metric space (U, D_h) (for certain special variants of the GFF) should correspond to the scaling limit of random planar maps, equipped with their graph distance, with respect to the Gromov–Hausdorff topology. Miller and Sheffield [16, 18, 19] constructed

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the LQG metric in the special case when $\gamma = \sqrt{8/3}$ using various special symmetries which are unique to this case. They also showed that, for certain special choices of the pair (U, h), the random metric space (U, D_h) agrees in law with a *Brownian surface*. Brownian surfaces, such as the Brownian map [13, 14] or the Brownian disk [1], are continuum random metric spaces which arise as the scaling limits of uniform random planar maps in the Gromov–Hausdorff topology.

This paper is part of a series of works which aims to construct the γ -LQG metric for all $\gamma \in (0,2)$. Analogously to the γ -LQG measure μ_h , this metric will be defined as the limit of metrics induced by continuous approximations of the GFF. To describe these approximations, we first need to define the γ -LQG dimension exponent $d_{\gamma} > 2$ from [4]. The value of this exponent is not known explicitly, but it can be defined in terms of various discrete approximations of LQG distances (such as random planar maps or Liouville first passage percolation, as discussed just below). Once the γ -LQG metric is constructed, it can be shown that d_{γ} is its Hausdorff dimension [10]. Let

$$\xi := \frac{\gamma}{d_{\gamma}}.$$

For concreteness, we will primarily focus on the whole-plane case (but see Remark 1.2). We say that a random distribution h on $\mathbb C$ is a whole-plane GFF plus a continuous function if there exists a coupling of h with a random continuous function $f:\mathbb C\to\mathbb R$ such that the law of h-f is that of a whole-plane GFF. We similarly define a whole-plane GFF plus a bounded continuous function, except that we also require that f is bounded. Let $p_s(z,w):=\frac{1}{2\pi s}\exp(-\frac{|z-w|^2}{2s})$ be the heat kernel, so that $p_s(z,\cdot)$ approximates a point mass at z when s is small. If h is a whole-plane GFF plus a bounded continuous function, we define the convolution

$$(3) h_{\epsilon}^*(z) := (h * p_{\epsilon^2/2})(z) = \int_{\mathbb{C}} h(w) p_{\epsilon^2/2}^U(z, w) d^2 w \quad \forall z \in \mathbb{C} \ \forall \epsilon > 0,$$

where the integral is interpreted as a distributional pairing. For $z, w \in \mathbb{C}$ and $\epsilon > 0$, we define the ϵ -Liouville first passage percolation (LFPP) metric by ¹

(4)
$$D_h^{\epsilon}(z, w) := \inf_{P: z \to w} \int_0^1 e^{\xi h_{\epsilon}^*(P(t))} |P'(t)| dt,$$

where the infimum is over all piecewise continuously differentiable paths from z to w. It was shown by Ding, Dubédat, Dunlap and Falconet [3] that the laws of the LFPP metrics, rescaled by the median distance across a square, are tight and every subsequential limit induces the Euclidean topology on \mathbb{C} . Subsequently, it was shown in [5, 7] that every subsequential limit can be realized as a measurable function of h, so, in fact, the convergence occurs in probability. Our eventual aim is to show that the subsequential limit of LFPP is unique, so it can be legitimately called $the \ \gamma$ -LQG metric.

The contribution of the present paper is to study the geometry of geodesics for subsequential limits of LFPP. In particular, we will prove that (in contrast to geodesics for smooth metrics) such geodesics satisfy the following property. For any fixed point z, a.s. any two geodesics from distinct points and to z will merge into each other and coincide for an interval of time before reaching z. Hence, the LQG geodesics towards z form a tree-like structure. This property is called *confluence of geodesics*. See Figure 1 for a simulation illustrating this

¹The intuitive reason why we look at $e^{\xi h_{\epsilon}(z)}$ instead of $e^{\gamma h_{\epsilon}(z)}$ to define the metric is as follows: Since μ_h is obtained by exponentiating γh , we can scale LQG areas by a factor of C>0 by adding $\gamma^{-1}\log C$ to the field. By (4), this results in scaling distances by $C^{\xi/\gamma}=C^{1/d_{\gamma}}$ which is consistent with the fact that the "dimension" should be the exponent relating the scaling of areas and distances.

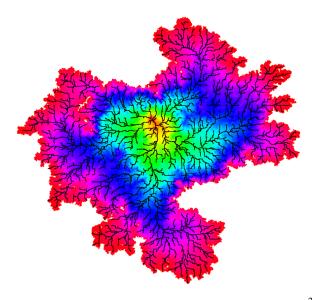


FIG. 1. A metric ball generated from a discrete GFF on a 1024×1024 subset of \mathbb{Z}^2 using the metric where the weight of each path P is given by the sum of $e^{\xi h(x)}$ over those x visited by P where $\xi = 1/\sqrt{6} = \gamma/d_{\gamma}$ for $\gamma = \sqrt{8/3}$ and $d_{\gamma} = 4$. It is believed that these approximations fall into the same universality class as for other forms of LFPP, for example, (4). All of the geodesics from the center of the metric ball to points in the intersection of the metric ball with a grid of spacing 20 (i.e., $(20\mathbb{Z})^2$) are shown.

property. The analogous property for geodesics in the Brownian map was proven in [12] and played an important role in the proof of the uniqueness of the Brownian map in [13, 14] and in the identification of the $\sqrt{8/3}$ -LQG metric with the Brownian map [16–19]. Likewise, the results of this paper are an important tool in the proof that the subsequential limit of LFPP is unique in [8].

Our proofs do *not* use very much external input. Indeed, aside from basic properties of the GFF (as can be found, e.g., in [27] and the introductory sections of [20, 22, 26]), we will only use a few lemmas from [5, 7, 15] which can be easily taken as black boxes. In fact, even the proofs of these external lemmas do not use much beyond basic facts about the GFF.

1.2. Weak LQG metrics. We will not work with LFPP directly in this paper. Instead, we will work in a more general framework involving a random metric which satisfies a list of axioms which are known to be satisfied for every subsequential limit of LFPP (see [5], Theorem 1.2). To state these axioms precisely, we will need some preliminary definitions concerning metric spaces. Let (X, D) be a metric space.

For r > 0, we write $\mathcal{B}_r(A; D)$ for the open ball consisting of the points $x \in X$ with D(x, A) < r. If $A = \{y\}$ is a singleton, we write $\mathcal{B}_r(\{y\}; d_X) = \mathcal{B}_r(y; d_X)$.

For a curve $P:[a,b] \to X$, the *D-length* of *P* is defined by

len(P; D) :=
$$\sup_{T} \sum_{i=1}^{\#T} D(P(t_i), P(t_{i-1})),$$

where the supremum is over all partitions $T: a = t_0 < \cdots < t_{\#T} = b$ of [a, b]. Note that the D-length of a curve may be infinite.

For $Y \subset X$, the *internal metric of D on Y* is defined by

(5)
$$D(x, y; Y) := \inf_{P \subset Y} \operatorname{len}(P; D) \quad \forall x, y \in Y,$$

where the infimum is over all paths P in Y from x to y. Then, $D(\cdot, \cdot; Y)$ is a metric on Y, except that it is allowed to take infinite values.

We say that (X, D) is a *length space* if for each $x, y \in X$ and each $\epsilon > 0$, there exists a curve of D-length at most $D(x, y) + \epsilon$ from x to y.

A continuous metric on an open domain $U \subset \mathbb{C}$ is a metric D on U which induces the Euclidean topology on U, that is, the identity map $(U, |\cdot|) \to (U, D)$ is a homeomorphism. We equip the space of continuous metrics on U with the local uniform topology for functions from $U \times U$ to $[0, \infty)$ and the associated Borel σ -algebra. We allow a continuous metric to satisfy $D(u, v) = \infty$ if u and v are in different connected components of U. In this case, in order to have $D^n \to D$ w.r.t. the local uniform topology, we require that for large enough n, $D^n(u, v) = \infty$ if and only if $D(u, v) = \infty$.

We are now ready to state the axioms under which we will work throughout the rest of the paper. Let $\mathcal{D}'(\mathbb{C})$ be the space of distributions (generalized functions) on \mathbb{C} , equipped with the usual weak topology. For $\gamma \in (0,2)$, a *weak* γ -*LQG metric* is a measurable function $h \mapsto D_h$ from $\mathcal{D}'(\mathbb{C})$ to the space of continuous metrics on \mathbb{C} such that the following is true whenever h is a whole-plane GFF plus a continuous function:

- I. Length space. Almost surely, (\mathbb{C}, D_h) is a length space, that is, the D_h -distance between any two points of \mathbb{C} is the infimum of the D_h -lengths of D_h -continuous paths (equivalently, Euclidean continuous paths) between the two points.
- II. Locality. Let $U \subset \mathbb{C}$ be a deterministic open set. The D_h -internal metric $D_h(\cdot,\cdot;U)$ is a.s. determined by $h|_U$.
 - III. Weyl scaling. Let ξ be as in (2), and, for each continuous function $f: \mathbb{C} \to \mathbb{R}$, define

(6)
$$(e^{\xi f} \cdot D_h)(z, w) := \inf_{P: z \to w} \int_0^{\operatorname{len}(P; D_h)} e^{\xi f(P(t))} dt \quad \forall z, w \in \mathbb{C},$$

where the infimum is over all continuous paths from z to w parameterized by D_h -length. Then, a.s. $e^{\xi f} \cdot D_h = D_{h+f}$ for every continuous function $f: \mathbb{C} \to \mathbb{R}$.

IV. Translation invariance. For each deterministic point $z \in \mathbb{C}$, a.s. $D_{h(\cdot+z)} = D_h(\cdot+z,\cdot+z)$.

V. Tightness across scales. Suppose that h is a whole-plane GFF, and let $\{h_r(z)\}_{r>0,z\in\mathbb{C}}$ be its circle average process. For each r>0, there is a deterministic constant $\mathfrak{c}_r>0$ such that the set of laws of the metrics $\mathfrak{c}_r^{-1}e^{-\xi h_r(0)}D_h(r\cdot,r\cdot)$ for r>0 is tight (w.r.t. the local uniform topology). Furthermore, the closure of this set of laws w.r.t. the Prokhorov topology on continuous functions $\mathbb{C}\times\mathbb{C}\to[0,\infty)$ is contained in the set of laws on continuous metrics on \mathbb{C} (i.e., every subsequential limit of the laws of the metrics $\mathfrak{c}_r^{-1}e^{-\xi h_r(0)}D_h(r\cdot,r\cdot)$ is supported on metrics which induce the Euclidean topology on \mathbb{C}). Finally, there exists $\Lambda>1$ such that, for each $\delta\in(0,1)$,

(7)
$$\Lambda^{-1}\delta^{\Lambda} \le \frac{\mathfrak{c}_{\delta r}}{\mathfrak{c}_r} \le \Lambda\delta^{-\Lambda} \quad \forall r > 0.$$

The following theorem is [5], Theorem 1.2, and is proven building on [3, 7]:

THEOREM 1.1 ([3, 5, 7]). Let $\gamma \in (0, 2)$. For every sequence of positive ϵ 's tending to zero, there is a weak γ -LQG metric D and a subsequence $\{\epsilon_k\}_{k\in\mathbb{N}}$ for which the following is true. Let h be a whole-plane GFF plus a bounded continuous function. Then the rescaled LFPP metrics $\mathfrak{a}_{\epsilon_k}^{-1}D_h^{\epsilon_k}$ from (4) converge in probability to D_h .

In light of Theorem 1.1, to prove theorems about subsequential limits of LFPP, it suffices to prove theorems about weak γ -LQG metrics. A particular advantage of this approach is that the Miller–Sheffield $\sqrt{8/3}$ -LQG metric [16, 18, 19] is a weak $\sqrt{8/3}$ -LQG metric (see [9],

²We refer to [6], Section 3.1, for the basic properties of the circle average process.

Section 2.5, for a careful explanation of why this is the case). So, all of our results also apply to this metric.

Let us now briefly comment on the above axioms. From the perspective that LQG is the random two-dimensional Riemannian manifold obtained by exponentiating the GFF, it is clear that Axioms I, II and III should be true for any reasonable notion of a metric on LQG. It is also not hard to see, at least heuristically, why these axioms should be satisfied for subsequential limits of LFPP.

It is expected that the LQG metric should satisfy a coordinate change formula when we apply a complex affine map which is directly analogous to the coordinate change formula for the LQG measure [6], Proposition 2.1. In particular, it should be the case that for any $a, b \in \mathbb{C}$, $a \neq 0$, a.s.

(8)
$$D_h(a \cdot +b, a \cdot +b) = D_{h(a \cdot +b)+Q \log |a|}(\cdot, \cdot) \quad \text{for } Q = \frac{2}{\nu} + \frac{\gamma}{2}.$$

It is not known at this point that subsequential limits of LFPP satisfy (8) (this will be proven in [8], using the result of the present paper). Axioms IV and V are a substitute for the formula (8), which is why we differentiate between the case of a true (strong) LQG metric and a weak LQG metric. These axioms imply the tightness of various functionals of D_h . For example, if $U \subset \mathbb{C}$ is open and $K \subset \mathbb{C}$ is compact, then the laws of

(9)
$$\left(\mathfrak{c}_{r}^{-1}e^{-\xi h_{r}(0)}D_{h}(rK,r\partial U)\right)^{-1}$$
 and $\mathfrak{c}_{r}^{-1}e^{-\xi h_{r}(0)}\sup_{u,v\in rK}D_{h}(u,v;rU)$

as r varies are tight.

REMARK 1.2 (Metrics associated with other fields). Suppose that D is a weak γ -LQG metric. Then, D_h is defined whenever h is a whole-plane GFF plus a continuous function. It is not hard to see that one can also define the metric if h is equal to a whole-plane GFF plus a continuous function plus a finite number of logarithmic singularities of the form $-\alpha \log |\cdot -z|$ for $z \in \mathbb{C}$ and $\alpha < Q$; see [5], Theorem 1.10 and Proposition 3.17.

We can also define metrics associated with GFF's on proper subdomains of \mathbb{C} . To this end, let $U \subset \mathbb{C}$ be open, and let h be a whole-plane GFF. Due to Axiom II, we can define for each open set $U \subset \mathbb{C}$ the metric $D_{h|_U} := D_h(\cdot, \cdot; U)$ as a measurable function of $h|_U$. We can write $h|_U = \mathring{h}^U + \mathfrak{h}^U$, where \mathring{h}^U is a zero-boundary GFF on U and $\mathring{\mathfrak{h}}^U$ is a random harmonic function on U independent from \mathring{h}^U . In the notation (6), we define

$$(10) D_{hU} := e^{-\xi \mathfrak{h}^U} \cdot D_{h|_U}.$$

It is easily seen from Axioms II and III that $D_{\mathring{h}^U}$ is a measurable function of \mathring{h}^U . Indeed, if we are given an open set $V \subset U$ with $\overline{V} \subset U$, choose a smooth compactly supported bump $f: U \to [0,1]$ which is identically equal to 1 on V. Then, Axiom II applied to the field $h-f\mathfrak{h}^U$ implies that the internal metric of $D_{\mathring{h}^U}$ on V, which equals $D_{h-f\mathfrak{h}^U}(\cdot,\cdot;V)$, is a.s. determined by $(h-f\mathfrak{h}^U)|_V = \mathring{h}^U|_V$. Letting V increase to all of U gives the desired measurability of $D_{\mathring{h}^U}$ w.r.t. \mathring{h}^U . This defines the γ -LQG metric for a zero-boundary GFF. Using Axiom III, we can similarly define the metric for a zero-boundary GFF plus a continuous function on U.

1.3. Main results. Let $\gamma \in (0, 2)$, let D be a weak γ -LQG metric and let h be a whole-plane GFF. For concreteness, we can fix the additive constant for h so that its average on the unit circle is equal to 0, but we note that by Axiom III the geodesics associated with D_h do not depend on how the additive constant is fixed. It is easy to check that D_h is boundedly compact, in the sense that closed, D_h -bounded subsets of $\mathbb C$ are compact [5], Lemma 3.8. By Axiom I and a standard result in metric geometry [2], Corollary 2.5.20, this implies that for

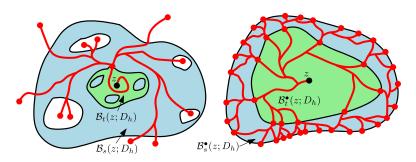


FIG. 2. Left: Theorem 1.3 asserts that for a fixed $z \in \mathbb{C}$, it is a.s. the case that, for each s > 0 and each small enough t > 0 (depending on s), all of the D_h -geodesics from z to points outside of $\mathcal{B}_s(z; D_h)$ coincide until they exit $\mathcal{B}_t(z; D_h)$. Several such D_h -geodesics are shown in red. Right: Theorem 1.4 asserts that a.s. for any $0 < t < s < \infty$ there are only finitely many possibilities for the restriction of a D_h -geodesic from 0 to $\partial \mathcal{B}_s^{\bullet}(0; D_h)$ to the time interval [0, t]. Several such leftmost D_h -geodesics are shown in red.

any two points in \mathbb{C} can be joined by at least one path of minimal D_h -length, which we call a D_h -geodesic.

Recall that for s > 0 and $z \in \mathbb{C}$, we write $\mathcal{B}_s(z; D_h)$ for the D_h -metric ball of radius s centered at z. The simplest form of confluence which we establish in this paper is the following statement; see Figure 2, left, for an illustration.

THEOREM 1.3 (Confluence of geodesics at a point). Almost surely, for each radius s > 0 there exists a radius $t \in (0, s)$ such that any two D_h -geodesics from z to points outside of $\mathcal{B}_s(z; D_h)$ coincide on the time interval [0, t].

We emphasize that the property of D_h -geodesics described in Theorem 1.3 is very different from the behavior of geodesics w.r.t. a smooth Riemannian metric on \mathbb{C} : indeed, in the latter situation geodesics from two different points targeted at z intersect only at z.

It was shown by Le Gall that the analog of Theorem 1.3 holds for geodesics in the Brownian map based at a typical point sampled from the volume measure on the Brownian map [12], Corollary 7.7. Due to the equivalence of Brownian and $\sqrt{8/3}$ -LQG surfaces, this is equivalent to the statement that one a.s. has confluence of $\sqrt{8/3}$ -LQG geodesics at a typical sampled uniformly from the $\sqrt{8/3}$ -LQG area measure. However, Theorem 1.3 has new content even in the case of the $\sqrt{8/3}$ -LQG metric since it gives confluence at a *Lebesgue* typical point, rather than a quantum typical point.

Another form of confluence concerns geodesics across an annulus between two filled D_h -metric balls. For s > 0 and $z \in \mathbb{C}$, we define the *filled metric ball* $\mathcal{B}_s^{\bullet}(z; D_h)$ to be the union of $\overline{\mathcal{B}_s(z; D_h)}$ and the set of points in \mathbb{C} which are disconnected from ∞ by $\overline{\mathcal{B}_s(z; D_h)}$. Each point $x \in \partial \mathcal{B}_s^{\bullet}(z; D_h)$ lies at D_h -distance exactly s from z, so every D_h -geodesic from z to x stays in $\mathcal{B}_s^{\bullet}(z; D_h)$. For some atypical points x there might be many such D_h -geodesics, but there is always a distinguished D_h -geodesic from z to x, called the *leftmost geodesic*, which lies (weakly) to the left of every other D_h -geodesic from 0 to x if we stand at x and look outward from $\mathcal{B}_s^{\bullet}(z; D_h)$ (see Lemma 2.4).

THEOREM 1.4 (Confluence of geodesics across a metric annulus). Fix $z \in \mathbb{C}$. Almost surely, for each $0 < t < s < \infty$ there is a finite set of D_h -geodesics from 0 to $\partial \mathcal{B}_t^{\bullet}(z; D_h)$ such that every leftmost D_h -geodesic from 0 to $\partial \mathcal{B}_s^{\bullet}(z; D_h)$ coincides with one of these D_h -geodesics on the time interval [0, t]. In particular, there are a.s. only finitely many points of $\partial \mathcal{B}_s^{\bullet}(z; D_h)$ which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_s^{\bullet}(z; D_h)$.

See Figure 2, right, for an illustration of the statement of Theorem 1.4. The proof of Theorem 1.4, given in Section 3, is in some sense the core of the paper; see Section 3.1 for an overview of this argument. Theorem 1.3 will be deduced from Theorem 1.4 in Section 4.

We will actually prove much more quantitative versions of Theorem 1.4 which give upper bounds for the number of leftmost D_h -geodesics across an annulus between two filled LQG metric balls which are uniform with respect to the Euclidean size of the LQG metric balls. See Theorems 3.1 and 3.9. This will be important in [8]. However, once we know that every weak LQG metric satisfies the scale invariance property (8) (which will be established in [8]), Theorem 3.1 itself implies an estimate which is uniform w.r.t. the Euclidean size of the LQG metric balls.

We now briefly discuss how the results of this paper are used in [8]. Very roughly, the reason why confluence is useful in this setting is that it allows us to establish near-independence for events which depend on small segments of a D_h -geodesic which are separated from each other. To be more precise, fix distinct points \mathbb{Z} , $\mathbb{W} \in \mathbb{C}$, and let P be the (a.s. unique, see Lemma 2.2 below) D_h -geodesic from \mathbb{Z} to \mathbb{W} . Suppose we are given $0 < t < s < \infty$, and we want to study, for a small $\epsilon > 0$, the conditional law of $h|_{B_{\epsilon}(P(s))}$ given $P|_{[0,t]}$. A priori, this seems to be a very intractable conditioning since we do not know anything about the conditional law of h given $P|_{[0,t]}$. However, we can get around this problem using Theorem 1.4, as follows:

Let $\delta \in (0, 1)$ be a small parameter. If δ is chosen small enough relative to t and then ϵ is chosen small enough relative to δ , then a slightly more quantitative version of Theorem 1.4 tells us that the following is true. With high probability, there is an arc $I \subset \partial \mathcal{B}_{s-\delta}^{\bullet}$ such that (a) all of the leftmost D_h -geodesics from \mathbb{Z} to I agree on the time interval [0, t] and (b) we have $B_{\epsilon}(P(s)) \subset \mathbb{C} \setminus \mathcal{B}_{s}^{\bullet}$ and every point of $B_{\epsilon}(P(s))$ is much closer to I than to $\partial \mathcal{B}_{s-\delta}^{\bullet} \setminus I$ (to keep P(s) away from the endpoints of I, one can use Lemma 3.6 below). Condition (b) tells us that $P(s-\delta) \in I$, and, moreover, this is still the case if we change the behavior of the field in $B_{\epsilon}(P(s))$ in a "reasonable" way. Hence, we can change what happens in $B_{\epsilon}(P(s))$ without changing $P|_{[0,t]}$.

REMARK 1.5 (Confluence for more general fields). We have stated Theorems 1.3 and 1.4 for a whole-plane GFF, but the statements extend immediately to other variants of the GFF via local absolute continuity. The proofs of Theorems 1.3 and 1.4 also work if we replace h by $h - \alpha \log |\cdot|$ for $\alpha \in (-\infty, Q)$, in which case $D_{h^{\alpha}}$ makes sense as a continuous metric on \mathbb{C} . In particular, by taking $\alpha = \gamma$, we get that confluence of geodesics holds a.s. at a "quantum typical point." However, confluence of geodesics does not hold a.s. at all points simultaneously. For example, if z is a point on a D_h -geodesic P which is not one of the endpoints of P, then there are two distinct incoming geodesic "arms" emanating from P which do not coalesce, namely, the segments of P before and after it hits z. That is, the conclusion of Theorem 1.3 does not hold for such a choice of z.

1.4. *Basic notation*. We write $\mathbb{N} = \{1, 2, 3, ...\}$ and $\mathbb{N}_0 = \mathbb{N} \cup \{0\}$. For a < b, we define the discrete interval $[a, b]_{\mathbb{Z}} := [a, b] \cap \mathbb{Z}$.

If $f:(0,\infty)\to\mathbb{R}$ and $g:(0,\infty)\to(0,\infty)$, we say that $f(\epsilon)=O_\epsilon(g(\epsilon))$ (resp., $f(\epsilon)=o_\epsilon(g(\epsilon))$) as $\epsilon\to 0$ if $f(\epsilon)/g(\epsilon)$ remains bounded (resp., tends to zero) as $\epsilon\to 0$. We similarly define $O(\cdot)$ and $o(\cdot)$ errors as a parameter goes to infinity.

If $f, g: (0, \infty) \to [0, \infty)$, we say that $f(\epsilon) \leq g(\epsilon)$ if there is a constant C > 0 (independent from ϵ and possibly from other parameters of interest) such that $f(\epsilon) \leq Cg(\epsilon)$. We write $f(\epsilon) \approx g(\epsilon)$ if $f(\epsilon) \leq g(\epsilon)$ and $g(\epsilon) \leq f(\epsilon)$.

We will often specify any requirements on the dependencies on rates of convergence in $O(\cdot)$ and $o(\cdot)$ errors, implicit constants in \leq , etc., in the statements of lemmas/propositions/theorems, in which case we implicitly require that errors, implicit constants, etc., appearing in the proof satisfy the same dependencies.

For $z \in \mathbb{C}$ and r > 0, we write $B_r(z)$ for the Euclidean ball of radius r centered at z. We also define the open annulus

(11)
$$\mathbb{A}_{r_1, r_2}(z) := B_{r_2}(z) \setminus \overline{B_{r_1}(z)} \quad \forall 0 < r_r < r_2 < \infty.$$

- **2. Preliminaries.** In this section we prove some preliminary results which are needed for the main argument in Section 3. In Section 2.1 we prove some basic monotonicity properties of D_h -geodesics and construct the leftmost geodesics appearing in Theorem 1.4. In Section 2.2 we prove a version of the FKG inequality for D_h which says that certain nondecreasing functionals of D_h are positively correlated. In Section 2.3 we state a general lemma from [7, 15] concerning the approximate independence of the restrictions of the GFF to concentric annuli. In Section 2.4 we prove some deterministic geometric lemmas which will be used to control the geometry of filled D_h -metric balls.
- 2.1. Basic properties of LQG metric balls and geodesics. Let D be a weak γ -LQG metric for some $\gamma \in (0, 2)$, and let h be a whole-plane GFF. For s > 0 and $z \in \mathbb{C}$, we define the filled metric ball $\mathcal{B}_s^{\sigma}(z; D_h)$ as in the discussion just above Theorem 1.4.

We recall the definition of a local set of the GFF from [26], Lemma 3.9. Suppose (h, A) is a coupling of h with a random set A. We say that A is a local set for h if, for any open set $U \subset \mathbb{C}$, the event $\{A \cap U \neq \emptyset\}$ is conditionally independent from $h|_{\mathbb{C}\setminus U}$ given $h|_U$. If A is determined by h (which will be the case for all of the local sets we consider), this is equivalent to the statement that A is determined by $h|_U$ on the event $\{A \subset U\}$. For a local set A, we can condition on the pair $(A, h|_A)$; this is, by definition, the same as conditioning on the σ -algebra $\bigcap_{\epsilon>0} \sigma(A, h|_{B_{\epsilon}(A)})$. The conditional law of $h|_{\mathbb{C}\setminus A}$ given $(A, h|_A)$ is that of a zero-boundary GFF on $\mathbb{C}\setminus A$ plus a harmonic function on $\mathbb{C}\setminus A$ which is determined by $(A, h|_A)$.

LEMMA 2.1. Let $z \in \mathbb{C}$ be deterministic. If τ is a stopping time for the filtration generated by $(\overline{\mathcal{B}_s(z;D_h)}, h|_{\overline{\mathcal{B}_s(z;D_h)}})$, then $\overline{\mathcal{B}_{\tau}(z;D_h)}$ is a local set for h. The same is true with $\mathcal{B}_{\mathfrak{s}}^{\bullet}(z;D_h)$ in place of $\overline{\mathcal{B}_s(z;D_h)}$.

PROOF. The lemma in the case of an unfilled metric ball follows from a general fact about local metrics of the Gaussian free field; see [7], Lemma 2.2. We will now treat the case of a filled metric ball. We will use the following criterion from [26], Lemma 3.9: a closed set A coupled with the GFF h is a local set if and only if for each open set $U \subset \mathbb{C}$, the event $\{A \subset U\}$ is conditionally independent from $h|_{\mathbb{C}\setminus U}$ given $h|_U$. For a deterministic radius s, the event $\{\mathcal{B}_s^{\bullet}(z;D_h)\subset U\}$ is the same as the event that $\overline{\mathcal{B}_s(z;D_h)}\subset U$ and each bounded connected component of $U\setminus \overline{\mathcal{B}_s(z;D_h)}$ is contained in U. Hence, $\{\mathcal{B}_s^{\bullet}(z;D_h)\subset U\}$ is determined by the event $\{\overline{\mathcal{B}_s(z;D_h)}\subset U\}$ and the set $\overline{\mathcal{B}_s(z;D_h)}$ on this event. By [26], Lemma 3.9, and the locality of $\overline{\mathcal{B}_s(z;D_h)}$, it follows that $\{\mathcal{B}_s^{\bullet}(z;D_h)\subset U\}$ is determined by $h|_U$, so $\mathcal{B}_s^{\bullet}(z;D_h)$ is a local set. The case of stopping times, which take on only countably many possible values, is immediate from the case of deterministic times. The case of general stopping times follows from the standard strong Markov property argument (i.e., look at the times $2^{-n} \lceil 2^n \tau \rceil$ and send $n \to \infty$) and the fact that local sets behave well under limits [21], Lemma 6.8. \square

We abbreviate

(12)
$$\mathcal{B}_{s}^{\bullet} := \mathcal{B}_{s}^{\bullet}(0; D_{h}).$$

Note that each point $x \in \partial \mathcal{B}_s^{\bullet}$ lies at D_h -distance exactly s from 0. Hence, there is a D_h -geodesic P from 0 to s which stays in \mathcal{B}_s^{\bullet} . For each t < s, such a geodesic satisfies $P(t) \in S$

 $\partial \mathcal{B}_t^{\bullet}$. Indeed, $P(t) \in \partial \mathcal{B}_t(0; D_h)$ by the definition of a geodesic and P has to pass through $\partial \mathcal{B}_t^{\bullet}$ since $\partial \mathcal{B}_t^{\bullet}$ disconnects 0 from x. Conversely, if t < s and P' is a D_h -geodesic from $\partial \mathcal{B}_t^{\bullet}$ to $\partial \mathcal{B}_s^{\bullet}$ (i.e., a path of length s - t between these sets), then the concatenation of any D_h -geodesic from 0 to P'(0) with P' is a path of length s from 0 to $P'(s - t) \in \partial \mathcal{B}_s^{\bullet}$ so is a D_h -geodesic from 0 to P'(s - t).

The goal of this subsection is to prove various monotonicity statements for D_h -geodesics from 0 to $\partial \mathcal{B}_s^{\bullet}$ which, roughly speaking, tell us that such geodesics have to stay in cyclic order. There is some subtlety since there can be multiple D_h -geodesics from 0 to some points of $\partial \mathcal{B}_s^{\bullet}$, and moreover some geodesics might share a nontrivial segment, so one cannot expect to have exact monotonicity.

The starting point of our proofs is the following result from [15].

LEMMA 2.2 ([15]). Almost surely, for each $q \in \mathbb{Q}^2$ there is only one D_h -geodesic from 0 to q.

PROOF. This follows from the proof of [15], Theorem 1.2. We note that the theorem is stated for a strong LQG metric, but the proof does not use the coordinate change axiom. \Box

As a consequence of Lemma 2.2, we get the following lemma which says that D_h -geodesics started from 0 cannot cross D_h -geodesics from 0 to rational points. This will be the main tool in our monotonicity arguments.

LEMMA 2.3. For $q \in \mathbb{Q}^2$, let P_q be the a.s. unique D_h -geodesic from 0 to q. The following holds a.s.: If $q \in \mathbb{Q}^2$, P' is a D_h -geodesic started from 0, and $u \in P_q \cap P'$, then there is a time $s \ge 0$ such that $P_q(s) = P'(s) = u$ and $P_q(t) = P'(t)$ for each $t \in [0, s]$.

Lemma 2.3 says that if a D_h -geodesic started from a rational point and an arbitrary D_h -geodesic started from 0 meet after time 0, then they have to agree for a nontrivial interval of time. We emphasize, however, that we have not yet proven that there are *any* such pairs of D_h -geodesics started from 0 which meet, so Lemma 2.3 does not imply any sort of confluence.

PROOF OF LEMMA 2.3. Suppose $u \in P_q \cap P'$ is as in the statement of the lemma. Since each of P_q and P' is a D_h -geodesic started from 0, the time when each of P_q and P' hits u is equal to $s := D_h(0, u)$. Therefore, $P_q(s) = P'(s) = u$. The concatenation of $P'|_{[0,s]}$ and $P_q|_{[s,D_h(0,q)]}$ is a D_h -geodesic from 0 to q. This D_h -geodesic must coincide with P_q by the uniqueness of D_h -geodesics between rational points, whence $P_q(t) = P'(t)$ for each $t \in [0,s]$. \square

For general points $y \in \partial \mathcal{B}_s^{\bullet}$, there might be many (even infinitely many) D_h -geodesics from 0 to y. However, there are two distinguished geodesics which are in some sense maximal, in the following sense:

LEMMA 2.4. Almost surely, for each s > 0 and each $y \in \partial \mathcal{B}_s^{\bullet}$, there exists a unique leftmost (resp., rightmost) geodesic P_y^- (resp., P_y^+) from 0 to y such that each D_h -geodesic from 0 to y lies (weakly) to the right (resp., left) of P_y^- (resp., P_y^+) if we stand at y and look outward from \mathcal{B}_s^{\bullet} . Moreover, there are sequences of points $q_n^-, q_n^+ \in \mathbb{Q}^2 \setminus \mathcal{B}_s^{\bullet}$ such that the D_h -geodesics from 0 to q_n^{\pm} satisfy $P_{q_n^{\pm}} \to P_y^{\pm}$, uniformly.

See Figure 3, left, for an illustration of the statement and proof of Lemma 2.4. Note that if 0 < t < s, then $P_x^-|_{[0,t]}$ is the leftmost D_h -geodesic from 0 to $P_x^-(t)$, since the concatenation

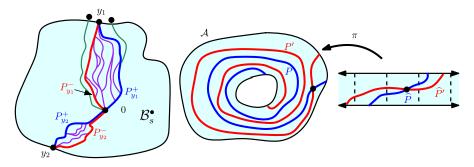


FIG. 3. Left: Two points $y_1, y_2 \in \partial \mathcal{B}_s^{\bullet}$ and their associated leftmost and rightmost D_h -geodesics (red and blue). Other D_h -geodesics from 0 to y_1 and y_2 are shown in purple. We have also shown two D_h -geodesics from 0 to points of $\mathbb{Q}^2 \setminus \mathcal{B}_s^{\bullet}$ (green) which approximate $P_{y_1}^-$ and $P_{y_1}^+$, respectively. Note that $P_{y_1}^-$ and $P_{y_1}^+$ intersect only at their endpoints, whereas $P_{y_2}^-$ and $P_{y_2}^+$ coincide for an initial time interval. Theorem 1.3 implies that a.s. the latter situation holds simultaneously for every $y \in \partial \mathcal{B}_s^{\bullet}$, but this has not been established yet. Middle: Illustration of the statement and proof of Lemma 2.5. Two paths P, P' in an annular region A are shown, with wind(P) slightly smaller than 2 and wind(P') slightly larger than 3. Right: To show that P and P' intersect, we consider their lifts to the universal cover of A chosen so that $\tilde{P}(0)$ lies to the right of $\hat{P}'(0)$. Then, $\hat{P}(1)$ lies to the left of $\hat{P}'(1)$ which forces the paths to intersect.

of $P_{P_x^-(t)}^-$ and $P_x^-|_{[t,s]}$ is a D_h -geodesic which lies (weakly) to the left of P_x^- . We remark that [17], Proposition 2.2, gives an alternative definition and construction of leftmost geodesics, but we will give an independent proof since we need the approximation by geodesics to rational points.

PROOF OF LEMMA 2.4. We will construct P_y^- ; the construction of P_y^+ is analogous. Fix a point $z \in \partial \mathcal{B}_s^{\bullet} \setminus \{y\}$. We can choose a sequence y_n^- of points in the clockwise arc of $\partial \mathcal{B}_s^{\bullet}$ from z to y which converge to y from the left (with "left" defined as in the lemma statement). Since D_h induces the Euclidean topology, for each n we can find $q_n^- \in \mathbb{Q}^2 \setminus \mathcal{B}_s^{\bullet}$ such that $D_h(q_n^-, y_n^-)$ is smaller than half of the distance from q_n^- to any point of $\partial \mathcal{B}_s^{\bullet}$ in the clockwise arc of $\partial \mathcal{B}_s^{\bullet}$ from y to z. Then, the points $P_{q_n^-}(s)$ converge to y from the left. Since each $P_{q_n^-}$ is a geodesic, the Arzéla–Ascoli theorem implies that, after possibly passing to a subsequence, we can arrange that the paths $P_{q_n^-}$'s converge uniformly to a continuous path P_y^- . The path P_y^- is a D_h -geodesic from 0 to y. By Lemma 2.3, no D_h -geodesic from 0 to y can cross any of the $P_{q_n^-}$'s. It follows that each such geodesic lies to the right of P_y^- . \square

We next need to rule out the possibility that some D_h -geodesics started from 0, wind around the origin more than others and, thereby, cross several other D_h -geodesics. It is possible for a D_h -geodesic to wind around 0 an arbitrarily large number of times, so, instead of bounding the amount of winding, we will show that it is approximately the same for all of the D_h -geodesics started from 0 (Lemma 2.6). To make this precise, we need to define the winding number of a curve in an annulus. Let \mathcal{A} be a topological space homeomorphic to a closed annulus, and let $\partial_{\text{in}}\mathcal{A}$ and $\partial_{\text{out}}\mathcal{A}$ be its inner and outer boundaries. Let $\pi: \mathbb{R} \times [0,1] \to \mathcal{A}$ be a universal covering map normalized so that $\pi(\mathbb{R} \times \{0\}) = \partial_{\text{in}}\mathcal{A}$, $\pi(\mathbb{R} \times \{1\}) = \partial_{\text{out}}\mathcal{A}$, and $\pi(x,y) = \pi(x+n,y)$ for each $n \in \mathbb{Z}$ and $(x,y) \in \mathbb{R} \times [0,1]$. For a path $P: [a,b] \to \mathcal{A}$, let $\widehat{P}: [a,b] \to \mathbb{R} \times [0,1]$ be a lift of P with respect to π . We define the winding number of P by

(13) wind(
$$P$$
) := (First coordinate of $\widehat{P}(b)$) – (First coordinate of $\widehat{P}(a)$).

Note that wind(P) does not depend on the choice of lift due to the periodicity of π . We also note that wind(P) is not required to be an integer since we consider general paths, not loops.

LEMMA 2.5. Let A be as above and let P, P' be paths in A from $\partial_{in}A$ to $\partial_{out}A$ which satisfy wind(P)+1 < wind(P'). Then, P and P' must cross each other in the sense that there is a point $w \in P \cap P'$ such that the segments of P and P' before hitting w do not coincide.

PROOF. The proof is completely elementary, but we give the details for the sake of completeness; see Figure 3, middle and right, for an illustration. Let $\pi : \mathbb{R} \times [0, 1] \to \mathcal{A}$ be the universal covering map from above. We parameterize our paths $P, P' : [0, 1] \to \mathcal{A}$.

Let $\widehat{P}:[0,1] \to \mathbb{R} \times [0,1]$ be a lift of P with respect to π started from a point in $[0,1) \times \{0\}$, and let $\widehat{x} \in [0,1)$ be the first coordinate of $\widehat{P}(0)$. We claim that we can choose a lift $\widehat{P}':[0,1] \to \mathbb{R} \times [0,1]$ of P' such that the first coordinate \widehat{x}' of $\widehat{P}'(0)$ lies in $[\widehat{x}-1,\widehat{x})$. Indeed, if the starting point of the lift of P' started from a point in $(0,1] \times \{0\}$ lies to the left of \widehat{x} , then we can take \widehat{P}' to be this lift. Otherwise, we can take \widehat{P}' to be the lift of P' started from a point in $(-1,0] \times \{0\}$, which (by the periodicity of π) is given by translating the lift started from a point in $(0,1] \times \{0\}$ one unit to the left.

By the definition of the winding number, the first coordinates of $\widehat{P}(1)$ and $\widehat{P}'(1)$ are $\widehat{y} = \widehat{x} + \text{wind}(P)$ and $\widehat{y}' = \widehat{x}' + \text{wind}(P')$, respectively. Since wind(P) + 1 < wind(P') and $\widehat{x}' \in [\widehat{x} - 1, \widehat{x})$, it follows that $\widehat{y}' > \widehat{y}$. Since also $\widehat{x}' < \widehat{x}$, it follows that P and P' must cross. \square

LEMMA 2.6. Almost surely, for each $0 < s < s' < \infty$, the winding numbers in $\mathcal{B}_{s'}^{\bullet} \setminus \mathcal{B}_{s}^{\bullet}$ of any two D_h -geodesics from $\partial \mathcal{B}_{s}^{\bullet}$ to $\partial \mathcal{B}_{s'}^{\bullet}$ (i.e., paths of D_h -length s' - s between these boundaries) differ by at most 1.

PROOF. For $q \in \mathbb{Q}^2 \setminus \mathcal{B}_{s'}^{\bullet}$, let P_q be the (a.s., unique) D_h -geodesic from 0 to q. The combination of Lemma 2.3 (applied to a D_h -geodesic from 0 to $\partial \mathcal{B}_{s'}^{\bullet}$, which coincides with the given geodesic on [s,s']), and Lemma 2.5 shows that a.s. the winding number of any D_h -geodesic from $\partial \mathcal{B}_{s'}^{\bullet}$ to $\partial \mathcal{B}_{s'}^{\bullet}$, differs from winding number of each of the $P_q|_{[s,s']}$'s by at most 1. Hence, there is an interval $[a,b] \subset \mathbb{R}$ of length at most one such that wind $(P_q|_{[s,s']}) \in [a,b]$ for every $q \in \mathbb{Q}^2 \setminus \mathcal{B}_{s'}^{\bullet}$. Since the winding number depends continuously on the path, the approximation statement from Lemma 2.4 implies that a.s. for each $y \in \partial \mathcal{B}_{s'}^{\bullet}$, the winding numbers of the leftmost and rightmost geodesics $P_y^-|_{[s,s']}$ and $P_y^+|_{[s,s']}$ lie in [a,b]. Since every D_h -geodesic from y to $\partial \mathcal{B}_s^{\bullet}$ lies between P_y^- and P_y^+ , the winding number of any such geodesic lies between wind $(P_y^-|_{[s,s']})$ and wind $(P_y^+|_{[s,s']})$ so must also lie in [a,b]. \square

The following lemma will be important in the iterative argument used to prove confluence of geodesics in Section 3.4:

LEMMA 2.7. Almost surely, the following is true for each $0 < s < s' < \infty$. Let \mathcal{I} be a finite collection of disjoint arcs of $\partial \mathcal{B}_s^{\bullet}$. For each $I \in \mathcal{I}$, let I' be the set of $y \in \partial \mathcal{B}_{s'}^{\bullet}$ such that the leftmost D_h -geodesic from 0 to y passes through I. Then, each I' is either empty or is a connected arc of $\partial \mathcal{B}_{s'}^{\bullet}$ and the arcs I' for different choices of $I \in \mathcal{I}$ are disjoint.

PROOF. It is obvious that the sets I' for different choices of I are disjoint since each $y \in \mathcal{B}_{s'}^{\bullet}$ gives rise to a unique leftmost geodesic. We need to show that each I' is connected.

Let $\pi: \mathbb{R} \times [0,1] \to \overline{\mathcal{B}_{s'}^{\bullet} \setminus \mathcal{B}_{s}^{\bullet}}$ be a universal covering map normalized so that $\pi(\mathbb{R} \times \{0\}) = \partial \mathcal{B}_{s}^{\bullet}$, $\pi(\mathbb{R} \times \{1\}) = \partial \mathcal{B}_{s'}^{\bullet}$ and $\pi(x,y) = \pi(x+n,y)$ for each $n \in \mathbb{Z}$ and $(x,y) \in \mathbb{R} \times \{0,1\}$. For each D_h -geodesic P from 0 to a point of $\overline{\mathbb{C} \setminus \mathcal{B}_{s'}^{\bullet}}$, let \widehat{P} be the lift of $P|_{[s,s']}$ w.r.t. π which starts from a point in $[0,1) \times \{0\}$. Lemma 2.6 implies that there is an interval $[a,b] \times \{1\}$ of length at most 1 such that $\widehat{P}(s') \in [a,b] \times \{1\}$ for every such D_h -geodesic P. For $q \in \mathbb{Q}^2 \setminus \mathcal{B}_{s'}^{\bullet}$, let P_q be the unique D_h -geodesic from 0 to q. For each such q, the path

 \widehat{P}_q divides $\mathbb{R} \times [0, 1]$ into two connected components. Lemma 2.3 implies that each lift of a

 D_h -geodesic from $\partial \mathcal{B}_s^{\bullet}$ to $\partial \mathcal{B}_{s'}^{\bullet}$ must be in the closure of one of these connected components, that is, it cannot cross \widehat{P}_a .

For $I \in \mathcal{I}$, let \widehat{I} be the subset of $[0,1) \times \{0\}$ such that $\pi(\widehat{I}) = I$, and let $\widehat{I'} \subset [a,b]$ be the subset of $[a,b] \times \{1\}$ such that $\pi(\widehat{I'}) = I'$. By possibly rechoosing π so that $\pi(0,0)$ is a endpoint of one of the intervals in \mathcal{I} , we can arrange that each \widehat{I} is an interval. To show that each I' is connected, it is enough to show that $\widehat{I'}$ is connected. Suppose $(y_1,1), (y_2,1) \in \widehat{I'}$ with $y_1 < y_2$ and that $y_1 < y < y_2$. We claim that $(y,1) \in \widehat{I'}$. To see this, let $P_{y_1}^-$ and $P_{y_2}^-$ be the leftmost D_h -geodesics from 0 to $\pi(y_1)$ and $\pi(y_2)$, respectively, so that $P_{y_1}^-(s), P_{y_2}^-(s) \in I$. If we choose $q_1, q_2 \in \mathbb{Q} \setminus \mathcal{B}_{s'}^{\bullet}$ in such a way that q_1 (resp., q_2) is close to the interior of the counterclockwise arc of $\partial \mathcal{B}_{s'}^{\bullet}$ from $\pi(y_1)$ to $\pi(y)$ (resp., $\pi(y)$ to $\pi(y_2)$), then $y_1 \leq \widehat{P}_{q_1}(s') < \widehat{P}_{q_2}(s') \leq y_2$.

The lifts \widehat{P}_{y_1} and \widehat{P}_{y_2} cannot cross \widehat{P}_{q_1} and \widehat{P}_{q_2} , so, since $\widehat{P}_{y_1}(s)$, $\widehat{P}_{y_2}(s) \in \widehat{I}$, it follows that $\widehat{P}_{q_1}(s)$ and $\widehat{P}_{q_2}(s)$ belong to \widehat{I} . Let P_y be the leftmost D_h -geodesic from \mathcal{B}_s^{\bullet} to $\pi(y)$. Then, \widehat{P}_y cannot cross \widehat{P}_{q_1} or \widehat{P}_{q_2} , so it must be the case that $\widehat{P}_y(y) \in \widehat{I}$. Therefore, $y \in \widehat{I}'$ so \widehat{I}' is connected. \square

2.2. *FKG inequalities*. In this subsection we will prove a variant of the FKG inequality which applies to a weak LQG metric. For the statement, we recall the definition of $e^{\xi f} \cdot D$ for a continuous function f and a metric D from (6). We also recall from Remark 1.2 that D gives rise to a metric associated with a zero-boundary GFF on any domain $U \subset \mathbb{C}$.

PROPOSITION 2.8 (FKG for the LQG metric). Let $\gamma \in (0,2)$, let $U \subset \mathbb{C}$ be an open domain, let h be a zero-boundary GFF on U and let D be a weak γ -LQG metric. Let Φ and Ψ be bounded, real-valued measurable functions on the space of continuous metrics on U which are nondecreasing in the sense that for any two such metrics D_1 , D_2 with $D_1(z,w) \leq D_2(z,w)$ for all $z,w \in U$, one has $\Phi(D_1) \leq \Phi(D_2)$ and $\Psi(D_1) \leq \Psi(D_2)$. Suppose further that Φ and Ψ are a.s. continuous at D_h in the sense that for every (possibly random) sequence of continuous functions $\{f^n\}_{n\in\mathbb{N}}$ which converges to zero uniformly on U, one has $\Phi(e^{\xi f^n} \cdot D_h) \to \Phi(D_h)$ and $\Psi(e^{\xi f^n} \cdot D_h) \to \Psi(D_h)$. Then, $Cov(\Phi(D_h), \Psi(D_h)) \geq 0$.

We will typically apply Proposition 2.8 to functionals of the form

(14)
$$\Phi(D) = \mathbb{I}\left\{\sup_{u \in A, v \in B} D(u, v) \ge c\right\} \quad \text{or} \quad \Phi(D) = \mathbb{I}\left\{\inf_{u \in A, v \in B} D(u, v) \ge c\right\},$$

where c > 0 and $A, B \subset U$. Note that functionals of this form include distances between points and sets as well as diameters of sets (taking A = B). Such functionals Φ are obviously nondecreasing. Moreover, such functionals are a.s. continuous at D_h since the probability that the supremum or infimum in question is exactly equal to c is zero. This can be seen using Axiom III and the fact that adding a smooth compactly supported function to h affects its law in an absolutely continuous way.

We expect that Proposition 2.8 is true without the continuity hypothesis, but our proof does not give this.

The basic idea of the proof of Proposition 2.8 is to first prove a version of the FKG inequality for continuous, positively-correlated Gaussian functions using the FKG inequality for finite-dimensional Gaussian vectors [23] and an approximation argument (Lemma 2.9). We then transfer this to the GFF using the white noise decomposition (Lemma 2.10) and finally deduce Proposition 2.8 using Axiom III. These intermediate FKG inequalities appear to be standard results, but we could not find sufficiently general statements in the literature so we will deduce them directly from the FKG inequality for finite-dimensional positive correlated Gaussian vectors [23].

LEMMA 2.9 (FKG for continuous Gaussian functions). Let (X, D) be a locally compact metric space, and let $C(X, \mathbb{R})$ be the space of continuous, real-valued functions on X equipped with the local uniform topology. Let Φ and Ψ be bounded measurable functions from $C(X, \mathbb{R})$ to \mathbb{R} which are nondecreasing in the sense that $\Phi(f) \geq \Phi(g)$ and $\Psi(f) \geq \Psi(g)$ whenever $f(x) \geq g(x)$ for every $x \in X$. Let f be a Gaussian random continuous function on X, and suppose that $Cov(f(x), f(y)) \geq 0$ for every $x, y \in X$. If Φ and Ψ are each a.s. continuous at f, then $Cov(\Phi(f), \Psi(f)) > 0$.

PROOF. We will deduce the lemma from the FKG inequality for finite positively correlated Gaussian vectors [23]. To do this, we will approximate f by a sequence of functions which depend on only finitely many Gaussian random variables (this is how the continuity assumption on Φ and Ψ is used).

Let $\{K_n\}_{n\in\mathbb{N}}$ be an increasing sequence of compact subsets of X whose union is all of X. For $n\in\mathbb{N}$, choose $\delta_n>0$ such that

(15)
$$\mathbb{P}\Big[\sup_{x,y\in K_n: D(x,y)\leq \delta_n} |f(x)-f(y)| \leq 2^{-n}\Big] \geq 1-2^{-n}.$$

Then, choose a finite collection of points $x_1^n, \ldots, x_{N_n}^n$ such that the union of the metric balls $\mathcal{B}_{\delta_n}(x_j^n)$ covers K_n . Let $\phi_1^n, \ldots, \phi_{N_n}^n$ be a partition of unity subordinate to this open cover of K_n , so that each ϕ_j^n is a continuous function on X taking values in [0, 1] and supported in $\mathcal{B}_{\delta_n}(x_j^n)$ and $\sum_{j=1}^{N_n} \phi_j^n \equiv 1$ on K_n .

 $\mathcal{B}_{\delta_n}(x_j^n)$ and $\sum_{j=1}^{N_n} \phi_j^n \equiv 1$ on K_n . Let $f^n := \sum_{j=1}^{N_n} f(x_j^n) \phi_j^n$. Since $\phi_j^n(x) = 0$ whenever $|x - x_j^n| > \delta_n$, the Borel–Cantelli lemma and (15) tell us that a.s. for large enough n,

(16)
$$|f(x) - f^n(x)| \le \sum_{i=1}^{N_n} |f(x) - f(x_j^n)| \phi_j^n(x) \le 2^{-n} \quad \forall x \in K_n.$$

Therefore, $f^n \to f$ uniformly on compact subsets of X.

The function f^n depends only on the values $f(x_1^n), \ldots, f(x_{N_n}^n)$ and increasing one of these values can only increase f^n . Hence, $\Phi(f^n)$ and $\Psi(f^n)$ are nondecreasing functions of the finite-dimensional, positively-correlated Gaussian vector $(f(x_1^n), \ldots, f(x_{N_n}^n))$. By the FKG inequality for finite-dimensional positively correlated Gaussian vectors [23],

$$\mathbb{E}[\Phi(f^n)\Psi(f^n)] \ge \mathbb{E}[\Phi(f^n)]\mathbb{E}[\Psi(f^n)].$$

Since $f^n \to f$ locally uniformly and Φ and Ψ are each bounded and a.s. continuous at f, sending $n \to \infty$ now concludes the proof. \square

LEMMA 2.10 (FKG for the GFF). Let $U \subset \mathbb{C}$ be an open domain, and let h be a zero-boundary GFF on U. Let Φ and Ψ be bounded, real-valued measurable functions on the space of distributions on U which are nondecreasing in the sense that for any continuous, nonnegative function f on U, one has $\Phi(h+f) \geq \Phi(h)$ and $\Psi(h+f) \geq \Psi(h)$. Suppose also that Φ and Ψ are a.s. continuous at h in the sense that for any sequence of (possibly random) functions $\{f^n\}_{n\in\mathbb{N}}$ on U which converge uniformly to zero on compact subsets of U, a.s. $\Phi(h+f^n) \to \Phi(h)$ and $\Psi(h+f^n) \to \Psi(h)$. Then, $Cov(\Phi(h), \Psi(h)) \geq 0$.

PROOF. We will use the white noise decomposition of h to reduce to the statement of Lemma 2.9. Let W be a space-time white noise on $U \times [0, \infty)$ so that for any square integrable functions $\phi, \psi: U \times [0, \infty) \to \mathbb{R}$, the random variables $\int_0^\infty \int_U \phi(w, t) W(d^2w, dt)$ and $\int_0^\infty \int_U \psi(w, t) W(d^2w, dt)$ are centered Gaussian with covariance $\int_0^\infty \int_U \phi(w, t) \psi(w, t) d^2w dt$. Let $p_t(z, w)$ for $z, w \in U$ be the transition density for Brownian motion run up to

time t and stopped upon exiting U, that is, if B is such a Brownian motion started from z, then $p_t(z, w) = \mathbb{P}[B_t \in d^2w, B([0, t]) \in U]$. For $0 \le s < t \le \infty$ and $z \in U$, let

$$h_{s,t}(z) := \sqrt{\pi} \int_{s^2}^{t^2} p_{r/2}(z, w) W(d^2 w, dr).$$

It is easily checked using the Kolmogorov continuity criterion that $h_{s,t}$ for $0 < s \le t \le \infty$ a.s. admits a continuous modification. For s = 0, $h_{0,t}$ does not admit a continuous modification and is instead interpreted as a random distribution. By [25], Lemma 5.4, $h := h_{0,\infty}$ is the zero-boundary GFF on U. Since $p_t(z, w)$ is nonnegative, each $h_{s,t}$ has nonnegative covariances.

For Φ as in the statement of the lemma and t>0, we define a functional on the space of continuous functions $f:U\to\mathbb{R}$ by $\Phi_{h_{0,t}}(f):=\Phi(f+h_{0,t})$. Then, for a fixed realization of $h_{0,t}$, the functional $\Phi_{h_{0,t}}$ is nondecreasing in the sense of Lemma 2.9. Moreover, for a continuous function f on U one has $\Phi_{h_{0,t}}(h_{t,\infty}+f)=\Phi(h+f)$, so $\Phi_{h_{0,t}}$ is a.s. continuous at $h_{t,\infty}$ in the sense of Lemma 2.9. The continuous Gaussian function $h_{t,\infty}$ is positively correlated, Gaussian and independent from $\{h_{0,s}\}_{s\leq t}$. Therefore, if we define $\Psi_{h_{0,t}}$ analogously to $\Phi_{h_{0,t}}$ and apply Lemma 2.9 to $\Phi_{h_{0,t}}(h_{t,\infty})=\Phi(h)$ and $\Psi_{h_{0,t}}(h_{t,\infty})=\Psi(h)$, we get

(17)
$$\mathbb{E}\big[\Phi(h)\Psi(h)|\{h_{0,s}\}_{s\leq t}\big] \geq \mathbb{E}\big[\Phi(h)|\{h_{0,s}\}_{s\leq t}\big]\mathbb{E}\big[\Psi(h)|\{h_{0,s}\}_{s\leq t}\big].$$

By the Kolmogorov zero-one law applied to the independent random variables $\{h_{s,t}: 2^{-k} \le s \le t \le 2^{-k+1}\}$ for $k \in \mathbb{N}$, $\bigcap_{t>0} \sigma(\{h_{0,s}\}_{s \le t})$ is the trivial σ -algebra. We can therefore take a limit as $t \to 0$ in (17) and apply the backward martingale convergence theorem to conclude the proof. \square

PROOF OF PROPOSITION 2.8. By Axiom III, it is a.s. the case that for each continuous function f on U one has $D_{h+f} = e^{\xi f} \cdot D_h$. In particular, if f is nonnegative then $D_{h+f} \geq D_h$. Therefore $g \mapsto \Phi(D_g)$ and $g \mapsto \Phi(D_g)$ are nondecreasing and a.s. continuous at h in the sense of Lemma 2.10. Hence, the proposition statement follows from Lemma 2.10. \square

2.3. Iterating events for the GFF in an annulus. Throughout this subsection, h denotes a whole-plane GFF normalized so that $h_1(0) = 0$. A key tool in our proofs is the following local independence property for events which depend on the GFF in disjoint concentric annuli, which is essentially proven in [15], Section 4; see [7], Lemma 3.1, for the statement we use here. (The statement in [7] also allows the events to depend on a collection of metrics, instead of just the GFF, but we will not need this here since our metrics are determined by the GFF). We recall the definition of the annulus $\mathbb{A}_{r_1,r_2}(z) = B_{r_2}(z) \setminus \overline{B_{r_1}(z)}$ from (11).

LEMMA 2.11 ([7]). Fix $0 < s_1 < s_2 < 1$. Let $\{r_k\}_{k \in \mathbb{N}}$ be a decreasing sequence of positive real numbers such that $r_{k+1}/r_k \le s_1$ for each $k \in \mathbb{N}$, and let $\{E_{r_k}\}_{k \in \mathbb{N}}$ be events such that $E_{r_k} \in \sigma((h - h_{r_k}(0))|_{\mathbb{A}_{s_1r_k, s_2r_k}(0)})$ for each $k \in \mathbb{N}$. For $K \in \mathbb{N}$, let N(K) be the number of $k \in [1, K]_{\mathbb{Z}}$ for which E_{r_k} occurs:

1. For each a > 0 and each $b \in (0, 1)$, there exists $p = p(a, b, s_1, s_2) \in (0, 1)$ and $c = c(a, b, s_1, s_2) > 0$ such that if

(18)
$$\mathbb{P}[E_{r_k}] \ge p \quad \forall k \in \mathbb{N},$$

then

(19)
$$\mathbb{P}[N(K) < bK] \le ce^{-aK} \quad \forall K \in \mathbb{N}.$$

2. For each $p \in (0, 1)$, there exists a > 0, $b \in (0, 1)$, and c > 0, depending only on p, s_1 , s_2 such that if (18) holds, then (19) holds.

We have the following variant of Lemma 2.11 where we explore annuli outward instead of inward.

- LEMMA 2.12. Fix $1 < S_1 < S_2 < \infty$. Let $\{r_k\}_{k \in \mathbb{N}}$ be an increasing sequence of positive real numbers such that $r_{k+1}/r_k \ge S_2$ for each $k \in \mathbb{N}$ and let $\{E_{r_k}\}_{k \in \mathbb{N}}$ be events such that $E_{r_k} \in \sigma((h h_{r_k}(0))|_{\mathbb{A}_{S_1 r_k}, S_2 r_k(0)})$ for each $k \in \mathbb{N}$. For $K \in \mathbb{N}$, let N(K) be the number of $k \in [1, K]_{\mathbb{Z}}$ for which E_{r_k} occurs:
- 1. For each a > 0 and each $b \in (0, 1)$, there exists $p = p(a, b, S_1, S_2) \in (0, 1)$ and $c = c(a, b, S_1, S_2) > 0$ such that if

(20)
$$\mathbb{P}[E_{r_k}] \ge p \quad \forall k \in \mathbb{N},$$

then

(21)
$$\mathbb{P}[N(K) < bK] \le ce^{-aK} \quad \forall K \in \mathbb{N}.$$

- 2. For each $p \in (0, 1)$, there exists a > 0, $b \in (0, 1)$ and c > 0, depending only on p, S_1 , S_2 such that if (18) holds, then (19) holds.
- Lemma 2.12 can be proven using the exact same argument used to prove Lemma 2.11. Alternatively, Lemma 2.12 is an immediate consequence of Lemma 2.11 and the following lemma:
- LEMMA 2.13. Let h be a whole-plane GFF normalized so that $h_1(0) = 0$. Then, the composition of h with the inversion map $z \mapsto 1/z$ has the same law as h.
- PROOF. Let $(\cdot, \cdot)_{\nabla}$ denote the Dirichlet inner product. By the conformal invariance of $(\cdot, \cdot)_{\nabla}$, if f is any smooth function supported on a compact subset of $\mathbb{C} \setminus \{0\}$, one has $(h(1/\cdot), f)_{\nabla} = (h, f(1/\cdot))_{\nabla}$. Therefore, if g is another such function, then

$$\operatorname{Cov}((h(1/\cdot), f)_{\nabla}, (h(1/\cdot), g)_{\nabla}) = (f(1/\cdot), g(1/\cdot))_{\nabla} = (f, g)_{\nabla}$$
$$= \operatorname{Cov}((h, f)_{\nabla}, (h, g)_{\nabla}).$$

Since a Gaussian process is determined by its covariance structure, it follows that the restrictions of h and $h(1/\cdot)$ to any compact subset of $\mathbb{C}\setminus\{0\}$ agree in law modulo additive constant. We know that the additive constants for h and $h(1/\cdot)$ are the same since the condition $h_1(0)=0$ is preserved under inversion. The restrictions of h to compact subsets of $\mathbb{C}\setminus\{0\}$ a.s. determine h since the σ -algebras $\bigcap_{r>0}\sigma(h|_{B_r(0)})$ and $\bigcap_{R>0}\sigma(h|_{\mathbb{C}\setminus B_R(0)})$ are trivial. \square

2.4. Harmonic exposure of boundary intervals. In this section we will prove some deterministic geometric lemmas for general classes of domains in \mathbb{C} . These lemmas will eventually be applied to the complement of a filled D_h -metric ball in Sections 3 and 4. Since we have a rather poor understanding of the geometry of the boundary of such a filled metric ball, it is important that the bounds are uniform over all possible domains satisfying certain mild constraints. The reader may wish to skip this subsection on a first read and refer back to the estimates as they are used. We first prove a universal bound to the effect that most boundary arcs of a simply connected planar domain U containing zero can be disconnected from ∞ by a small set.

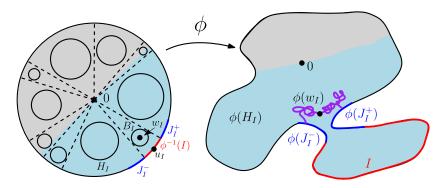


FIG. 4. Illustration of the proof of Lemma 2.14. For each arc $I \in \mathcal{I}$, we consider a point $w_I \in \mathbb{D}$ whose distance to $\phi^{-1}(I)$ is proportional to the length of $\phi^{-1}(I)$. A Brownian motion started from w_I has uniformly positive probability to first exit the half-disk H_I in each of the two arcs J_I^{\pm} on either side of $\phi^{-1}(I)$ before getting too far away from w_I . Using the Beurling estimate, this shows that I can be disconnected from 0 in $\phi(H_I)$ by the union of two paths whose Euclidean diameters are comparable to the Euclidean distance from $\phi(w_I)$ to ∂U . By distortion estimates, this distance is comparable to $|\phi'(w_I)|(1-|w_I|)$. There cannot be too many intervals for which this last quantity is large since $\int_{\mathbb{D}} |\phi'(z)|^2 d^2z = \operatorname{area}(U)$.

LEMMA 2.14. There is a universal constant A > 0 such that the following is true. Let $U \subset \mathbb{C}$ be a bounded simply connected domain containing 0, and view ∂U as a collection of prime ends. Also, let $n \in \mathbb{N}$, and let \mathcal{I} be a collection of $\#\mathcal{I} = n$ arcs of ∂U which intersect only at their endpoints. Then, for C > 0, the number of arcs $I \in \mathcal{I}$, which can be disconnected from 0 in U by a path in U of Euclidean diameter at most $Cn^{-1/2}$, is at least $(1 - AC^{-2} \operatorname{area}(U))n$.

PROOF. See Figure 4 for an illustration of the proof. Let $\phi : \mathbb{D} \to U$ be a conformal map with $\phi(0) = 0$. For $I \in \mathcal{I}$, let r_I be the length of the arc $\phi^{-1}(I)$. Since the arcs $\phi^{-1}(I)$ intersect only at their endpoints, there can be at most eight such arcs with $r_I \geq \pi/4$. By removing these arcs from our collection, we can assume without loss of generality that $r_I \leq \pi/4$ for each $I \in \mathcal{I}$.

For $I \in \mathcal{I}$, let J_I^- (resp., J_I^+) be the arc of $\partial \mathbb{D}$ of length r_I which comes immediately before (resp., after) $\phi^{-1}(I)$ in the counterclockwise direction. Also, let $u_I \in \partial \mathbb{D}$ be the center of $\phi^{-1}(I)$, and let H_I be the half-disk of \mathbb{D} with the property that u_I is the center of the arc $\partial \mathbb{D} \cap \partial H_I$. Since $r_I \leq \pi/4$, we have that J_I^- , $\phi^{-1}(I)$ and J_I^+ intersect only at their endpoints and are contained in ∂H_I .

Define $w_I := (1 - r_I)u_I$. The probability that a Brownian motion started from w_I exits H_I at a point in J_I^- is at least some universal constant a > 0. By the conformal invariance of Brownian motion, a Brownian motion started from $\phi(w_I)$ has probability at least a to exit $\phi(H_I)$ at a point in $\phi(J_I^-)$. On the other hand, the Beurling estimate shows that for R > 1, the probability that such a Brownian motion travels distance $R \operatorname{dist}(\phi(w_I), \partial U)$ without hitting ∂U , is at most a universal constant times $R^{-1/2}$. If we take R to be a sufficiently large universal constant, then this last quantity is smaller than a, so with positive probability, a Brownian motion started from $\phi(w_I)$ exits $\phi(H_I)$ at a point of $\phi(J_I^-)$ and does so before it leaves the ball of radius $R \operatorname{dist}(\phi(w_I), \partial U)$ centered at $\phi(w_I)$. Consequently, there is a path in $\phi(H_I)$ from $\phi(w_I)$ to $\phi(J_I^-)$ with Euclidean diameter at most $2R \operatorname{dist}(\phi(w_I), \partial U)$ which is disjoint from $\phi(H_I)$. Symmetrically, the same statement holds with J_I^+ in place of J_I^- .

Concatenating the paths for J_I^- and J_I^+ gives a path in $\phi(H_I)$ with Euclidean diameter at most $4R \operatorname{dist}(\phi(w_I), \partial U)$ from a point of $\phi(J_I^-)$ to a point of $\phi(J_I^+)$. This path necessarily divides U into at least two connected components. Since the path is disjoint from $\phi(H_I)$ and

 $0 \in \partial \phi(H_I)$, it must disconnect 0 from I in U. Consequently,

(22) I is disconnected from 0 in U by a path of diameter $\leq 4R \operatorname{dist}(\phi(w_I), \partial U)$.

We will now show that $\operatorname{dist}(\phi(w_I), \partial U)$ is small for most $I \in \mathcal{I}$. Set $B_I := B_{r_I/100}(w_I)$. We observe that B_I is contained in the slice of \mathbb{D} bounded by $\phi^{-1}(I)$ and the line segments from the two endpoints of $\phi^{-1}(I)$ to 0; so, since the arcs $\phi(I)$ for $I \in \mathcal{I}$ are disjoint, the balls B_I are disjoint. For $I \in \mathcal{I}$, one has $\operatorname{dist}(w_I, \partial \mathbb{D}) = r_I$. By the Koebe distortion theorem and the Koebe quarter theorem, there is a universal constant $A_0 > 0$ such that

(23)
$$\inf_{z \in B_I} |\phi'(z)| \ge A_0^{-1} |\phi'(w_I)| \quad \text{and} \quad \operatorname{dist}(\phi(w_I), \partial U) \le A_0 |\phi'(w_I)| r_I.$$

Applying the first inequality, then the second inequality, in (23) shows that

(24)
$$\operatorname{area}(U) = \int_{\mathbb{D}} |\phi'(z)|^2 d^2 z \ge \frac{1}{A_0^2} \sum_{I \in \mathcal{I}} |\phi'(w_I)|^2 \frac{\pi r_I^2}{100^2} \ge \frac{\pi}{100^2 A_0^4} \sum_{I \in \mathcal{I}} \operatorname{dist}(\phi(w_I), \partial U)^2.$$

Since $r_I \le 1$ for all I, we can apply the Chebyshev inequality to get that for C > 0,

(25)
$$\#\left\{I \in \mathcal{I} : \operatorname{dist}(\phi(w_I), \partial U) \ge \frac{C}{4R} n^{-1/2}\right\} \le AC^{-2} \operatorname{area}(U)n$$

for $A = 100^2 \pi^{-1} A_0^4 \times 16 R^2$. Combining (22) and (25) concludes the proof.

In the next several lemmas, we will use the following notation. For a domain $U \subset \mathbb{C}$ containing 0, we define

(26)
$$\operatorname{inrad}(U) := \sup\{r > 0 : B_r(0) \subset U\}$$
 and $\operatorname{outrad}(U) := \inf\{R > 0 : U \subset B_R(0)\}.$

Similarly, for a compact connected set $K \subset \mathbb{C}$ containing 0, we define

(27)
$$\operatorname{inrad}(K) := \sup\{r > 0 : B_r(0) \subset K\}$$
 and $\operatorname{outrad}(K) := \inf\{R > 0 : K \subset B_R(0)\}.$

By inverting, we obtain an analog of Lemma 2.15 when we want to disconnect boundary arcs from ∞ instead of from 0.

LEMMA 2.15. There is a universal constant A > 0 such that the following is true: Let $K \subset \mathbb{C}$ be a compact connected set which contains a neighborhood of 0 and whose complement is connected. Also, let $n \in \mathbb{N}$, and let \mathcal{I} be a collection of $\#\mathcal{I} = n$ arcs of ∂K which intersect only at their endpoints. Then, for C > 0, the number of arcs $I \in \mathcal{I}$ which can be disconnected from ∞ in $\mathbb{C} \setminus K$ by a path in $\mathbb{C} \setminus K$ of Euclidean diameter at most $Cn^{-1/2}$ is at least

(28)
$$\left(1 - \frac{A}{C^2} \frac{\operatorname{outrad}(K)^4}{\operatorname{inrad}(K)^2}\right) n.$$

PROOF. Let U be the image of $\mathbb{C} \setminus K$ under the map $\phi: z \mapsto 1/z$. Then, U is contained in the Euclidean ball of radius $1/\operatorname{inrad}(K)$ centered at zero, so $\operatorname{area}(U) \leq \pi/\operatorname{inrad}(K)^2$. Since $\inf_{z \in \mathbb{C} \setminus K} |\phi'(z)| \geq 1/\operatorname{outrad}(K)^2$, the diameter of each subset of $\mathbb{C} \setminus K$ is at most $\operatorname{outrad}(K)^2$ times the diameter of the corresponding subset of U. The statement of the lemma, therefore, follows from Lemma 2.14 applied to U, with $C/\operatorname{outrad}(K)^2$ in place of C. \square

We next prove two lemmas to the effect that a collection of arcs which cover ∂U (or ∂K) must include at least one arc which is sufficiently "exposed," in a certain quantitative sense. These lemmas will be used in the proof of Lemma 4.3. As in the previous lemmas, it is crucial that the bounds be uniform over all choices of U (or K). We will need a particular way of

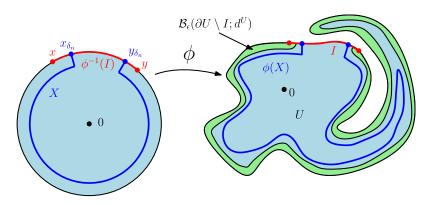


FIG. 5. Illustration of the proof of Lemma 2.16. We start by choosing $I \in \mathcal{I}$ such that the length of $\phi^{-1}(I)$ is bounded below by $2\pi/n$ (i.e., the harmonic measure of I from 0 in U is at least 1/n). We then build a blue "shield" X in \mathbb{D} which separates 0 and part of $\phi^{-1}(I)$ from the rest of $\partial \mathbb{D}$. We show using distortion estimates that the image of this shield under ϕ lies at uniformly positive d^U -distance from $\partial U \setminus I$. Hence, for a small enough $\epsilon > 0$, $\phi(X)$ separates $\mathcal{B}_{\epsilon \text{ inrad}(U)}(\partial U \setminus I; d^U)$ from 0 and from at least one point of I. The right side of the figure shows why we use d^U instead of Euclidean distance. The illustrated domain U has a long "arm" which is very close to I with respect to the ambient Euclidean distance, but it not close to I with respect to d^U .

measuring distances in U which is slightly different from the ordinary Euclidean distance (see Figure 5 for an illustration of why this is needed).

Let $U \subset \mathbb{C}$ be a simply connected domain, and view ∂U as a set of prime ends. If $X \subset U$, we define the *prime end closure* $\mathrm{Cl}'(X)$ to be the set of points in $z \in U \cup \partial U$ with the following property: if $\phi: U \to \mathbb{D}$ is a conformal map, then $\phi(z)$ lies in $\overline{\phi(X)}$. For $z, w \in U \cup \partial U$ we define

(29)
$$d^{U}(z, w) = \inf\{\operatorname{diam}(X) : X \text{ is a connected subset of } U \text{ with } z, w \in \operatorname{Cl}'(X)\},$$

where here diam denotes the Euclidean diameter. Then, d^U is a metric on $U \cup \partial U$ which is bounded below by the Euclidean metric on \mathbb{C} , restricted to $U \cup \partial U$ and bounded above by the internal Euclidean metric on $U \cup \partial U$. Note that d^U is not a length metric. We similarly define $\mathrm{Cl}'(\cdot)$ and d^U in the case when U is an unbounded open subset of \mathbb{C} such that $\mathbb{C} \setminus U$ is compact, in which case we use a conformal map to $\mathbb{C} \setminus \mathbb{D}$ instead of a conformal map to \mathbb{D} .

LEMMA 2.16. Let $U \subset \mathbb{C}$ be a bounded simply connected domain containing 0, and view ∂U as a collection of prime ends. Let $n \in \mathbb{N}$, and let \mathcal{I} be a collection of $\#\mathcal{I} = n$ arcs of ∂U whose union is all of ∂U (not necessarily disjoint). There exists $\epsilon > 0$ depending only on n and the ratio $\operatorname{outrad}(U)/\operatorname{inrad}(U)$ such that, for any choice of U and \mathcal{I} , there exists $I \in \mathcal{I}$ such that the following is true. The arc I is not disconnected from 0 in U by the d^U -neighborhood $\mathcal{B}_{\epsilon \operatorname{inrad}(U)}(\partial U \setminus I; d^U)$.

PROOF. See Figure 5 for an illustration and outline of the proof. By scaling, we can assume without loss of generality that inrad(U) = 1. Let R := outrad(U).

Let $\phi: \mathbb{D} \to U$ be a conformal map which fixes zero. Since $\#\mathcal{I} = n$, there must exist $I \in \mathcal{I}$ such that the $\partial \mathbb{D}$ -length of $\phi^{-1}(I)$ is at least $2\pi/n$. Let x and y be the endpoints of I, in counterclockwise order. Let $\delta_n := 1/(4n)$, and let x_{δ_n} (resp., y_{δ_n}) be the point of I which lies at $\partial \mathbb{D}$ -distance δ_n from x (resp., y) in the counterclockwise (resp., clockwise) direction. Let X be the union of the linear segments $[(1 - \delta_n)x_{\delta_n}, x_{\delta_n}]$ and $[(1 - \delta_n)y_{\delta_n}, y_{\delta_n}]$ and the counterclockwise arc A_{δ_n} of $\partial B_{1-\delta_n}(0)$ from y_{δ_n} to x_{δ_n} (X is shown in blue on the left side of Figure 5).

We claim that there is a constant $\epsilon = \epsilon(n, R) > 0$ such that

(30)
$$d^{U}(\phi(X), \partial U \setminus I) \ge \epsilon.$$

Assuming (30), we conclude the proof as follows. The relation (30) implies that $\phi(X)$ is disjoint from $\mathcal{B}_{\epsilon}(\partial U \setminus I; d^U)$. Since $\mathcal{B}_{\epsilon}(\partial U \setminus I; d^U)$ is connected, by conformally mapping back to \mathbb{D} it follows that there is a connected component of $U \setminus \phi(X)$ which contains 0 and which has some point of I in its prime end closure. Therefore, I is not disconnected from 0 by $\mathcal{B}_{\epsilon}(\partial U \setminus I; d^U)$ (recall that we have assumed that $\operatorname{inrad}(U) = 1$).

Let us now prove (30). Since we have assumed that $\operatorname{inrad}(U) = 1$, we have $\mathbb{D} \subset U$. By the Koebe quarter theorem, $\phi^{-1}(U) = \mathbb{D}$ contains a Euclidean ball of radius at least $|(\phi^{-1})'(0)|/4$ centered at 0. Therefore we must have $|(\phi^{-1})'(0)|/4 \le 1$ so $|\phi'(0)| \ge 1/4$. For $u \in A_{\delta_n}$ we have $|u| = 1 - \delta_n$ so by the Koebe distortion theorem $|\phi'(u)| \ge |\phi'(0)|\delta_n/(2 - \delta_n)^3 \ge \delta_n/(4(2 - \delta_n)^3)$. By the Koebe quarter theorem applied to the restriction of ϕ to $B_{\delta_n}(u)$, we then obtain

(31)
$$\operatorname{dist}(\phi(A_{\delta_n}), \partial U) \ge \frac{\delta_n^2}{16(2 - \delta_n)^3}.$$

We now need to deal with the small linear segments of X. Let $J_{\delta_n} \subset \phi^{-1}(I)$ be the counterclockwise arc of $\partial \mathbb{D}$ from x to x_{δ_n} . The probability that a Brownian motion started from 0 exits \mathbb{D} in J_{δ_n} before hitting the segment $[(1-\delta_n)x_{\delta_n},x_{\delta_n}]$ is at least a universal constant c>0 times δ_n . By the conformal invariance of Brownian motion, the probability that a Brownian motion started from 0 exits U in $\phi(J_{\delta_n})$ before hitting $\phi([(1-\delta_n)x_{\delta_n},x_{\delta_n}])$ is at least $c\delta_n$. Similarly, the probability that a Brownian motion started from 0 exits U in $I\setminus\phi(J_{\delta_n})$ before hitting $\phi([(1-\delta_n)x_{\delta_n},x_{\delta_n}])$ is at least $c\delta_n$.

There is a constant $\zeta = \zeta(n,R) > 0$ such that for any set $Y \subset B_R(0) \setminus B_{1/2}(0)$ with Euclidean diameter at most ζ , the probability that a Brownian motion started from 0 hits Y before exiting $B_R(0)$ is at most $(c/2)\delta_n$. The preceding paragraph implies that neither $\phi(J_{\delta_n})$ nor $I \setminus \phi(J_{\delta_n})$ can be disconnected from 0 in $U \setminus \phi([(1-\delta_n)x_{\delta_n},x_{\delta_n}])$ by a set of Euclidean diameter smaller than ζ . This implies that $d^U(\phi([(1-\delta_n)x_{\delta_n},x_{\delta_n}]),\partial U\setminus I)\geq \zeta$. A symmetric argument shows that the same is true for $\phi([(1-\delta_n)y_{\delta_n},y_{\delta_n}]),\partial U\setminus I)$. Combining this with (31), gives (30) with $\epsilon = \max\{\zeta, \frac{\delta_n^2}{16(2-\delta_n)^3}\}$. \square

We now deduce an analog of Lemma 2.16 for domains which are conformally equivalent to $\mathbb{C}\setminus\overline{\mathbb{D}}$ instead of to \mathbb{D} .

LEMMA 2.17. Let $K \subset \mathbb{C}$ be a compact connected set whose complement is connected, and view ∂K as a collection of prime ends. Let $n \in \mathbb{N}$, and let \mathcal{I} be a collection of $\#\mathcal{I} = n$ (not necessarily disjoint) arcs of ∂K whose union is all of ∂K . There exists $\epsilon > 0$ depending only on n and the ratio $\operatorname{outrad}(K)/\operatorname{inrad}(K)$ such that for any choice of U and \mathcal{I} , there exists $I \in \mathcal{I}$ such that the following is true. The arc I is not disconnected from ∞ in $\mathbb{C} \setminus K$ by the $d^{\mathbb{C} \setminus K}$ -neighborhood $\mathcal{B}_{\epsilon \operatorname{outrad}(U)}(K \setminus I; d^{\mathbb{C} \setminus K})$.

PROOF. Let $\phi(z) := 1/z$, and let $U := \phi(\mathbb{C} \setminus K)$. Then, U is simply connected, and we have $\operatorname{outrad}(U) = 1/\operatorname{inrad}(K)$ and $\operatorname{inrad}(U) = 1/\operatorname{outrad}(K)$. Lemma 2.16 shows that there exists $\epsilon > 0$ as in the statement of the lemma such that for any choice of U and \mathcal{I} , there exists $I \in \mathcal{I}$ such that I is *not* disconnected from ∞ in $\mathbb{C} \setminus K$ by the d^U -neighborhood $\phi(\mathcal{B}_{\epsilon/\operatorname{outrad}(K)}(\partial U \setminus I; d^U))$. We have $|\phi'(z)| \le 1/\operatorname{inrad}(K)^2$ on $\mathbb{C} \setminus B_{\operatorname{inrad}(K)}(0) \supset \mathbb{C} \setminus K$, so for any set $A \subset \mathbb{C} \setminus K$, the Euclidean diameter of $\phi(A)$ is at most $\operatorname{diam}(A)/\operatorname{inrad}(K)^2$. It

follows that

(32)
$$\phi(\mathcal{B}_{\epsilon' \operatorname{inrad}(K)}(K \setminus I; d^{\mathbb{C} \setminus K})) \subset \mathcal{B}_{\epsilon/\operatorname{outrad}(K)}(\partial U \setminus I; d^{U}) \quad \text{for } \epsilon' = \frac{\operatorname{inrad}(K)}{\operatorname{outrad}(K)} \epsilon.$$

Therefore, the statement of the lemma is true with ϵ' in place of ϵ . \square

3. Finitely many leftmost geodesics across an LQG annulus. In this section we will prove a more quantitative version of Theorem 1.4 (Theorem 3.1) which, as we explain just below, immediately implies Theorem 3.1. The main difference between Theorem 1.4 and Theorem 3.1 is that the latter gives a bound for the number of leftmost D_h -geodesics across an LQG annulus which is *uniform in the Euclidean size of the LQG annulus*. This is important since we are only assuming tightness across scales (Axiom V) instead of exact scale invariance. To quantify what Euclidean "scale" we will be working in,³ we will consider the following stopping times for $\{(\mathcal{B}_s^{\bullet}, h|_{\mathcal{B}_s^{\bullet}})\}_{s>0}$:

(33)
$$\tau_{\mathbb{I}} := \inf\{s > 0 : \mathcal{B}_s^{\bullet} \not\subset B_{\mathbb{I}}(0)\} \quad \forall r > 0.$$

We also recall the definition of leftmost geodesics from Lemma 2.4 and the scaling constants $\mathfrak{c}_{\mathbb{F}}$ for $\mathbb{F}>0$ from Axiom V.

THEOREM 3.1. For each t > 0 and $p \in (0,1)$, there exists $N = N(t,p) \in \mathbb{N}$ such that the following is true for each $\mathbb{r} > 0$ and each stopping time τ for $\{(\mathcal{B}_s^{\bullet}, h|_{\mathcal{B}_s^{\bullet}})\}_{s \geq 0}$ such that $\tau \in [\tau_{\mathbb{r}}, \tau_{2\mathbb{r}}]$ a.s. With probability at least p, that there are at most N points of $\partial \mathcal{B}_{\tau}^{\bullet}$ which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau+t\mathfrak{c}_{\mathbb{r}}}^{\bullet}e^{\xi h_{\mathbb{r}}(0)}$.

By Axiom V, typically $\tau_{\mathbb{F}}$ and $\tau_{2\mathbb{F}}$ are each comparable to $\mathfrak{c}_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}$, so typically $\tau + t\mathfrak{c}_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}$ is of the same order of magnitude as τ . Theorem 3.1 is the first result of this paper which gives any sort of confluence of geodesics. Prior to this point, we have not ruled out the possibility that any two distinct D_h -geodesics started from 0 intersect only at their common starting point.

PROOF OF THEOREM 1.4 ASSUMING THEOREM 3.1. By Axiom IV and the translation invariance of the law of h, viewed modulo additive constant, we can assume without loss of generality that z=0. For $0 < t < s < \infty$, let $X_{t,s}$ be the set of points $x \in \partial \mathcal{B}_t^{\bullet}$ which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_s^{\bullet}$. Theorem 3.1 says that for each stopping time τ as in that theorem and each t > 0, a.s. $X_{\tau,\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$ is finite.

We first observe that in the setting of Theorem 3.1, there is a.s. a unique D_h -geodesic from 0 to each point $x \in X_{\tau,\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$. Indeed, consider such a point x, and let $y \in \partial \mathcal{B}^{\bullet}_{\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$ such that the leftmost D_h -geodesic P_y from 0 to y passes through x. By Lemma 2.4, we can find $q \in \mathbb{Q}^2 \setminus \mathcal{B}^{\bullet}_{\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$ such that the distance between P_y and the (a.s. unique) D_h -geodesic P_q from 0 to q w.r.t. the Euclidean uniform metric is smaller than the minimum of |x-x'| over all $x' \in X_{\tau,\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}} \setminus \{x\}$. Since $P_q|_{[0,\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}]}$ is a leftmost D_h -geodesic from 0 to $\partial \mathcal{B}^{\bullet}_{\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$, it follows that P_q must pass through x. By the uniqueness of D_h -geodesics to rational points (Lemma 2.2), it follows that there is only on D_h -geodesic from 0 to x.

Therefore, it is a.s. the case that for each rational $\mathbb{r}>0$ and each rational b>0, the set $X_{\tau_{\mathbb{r}},\tau_{\mathbb{r}}+b\mathfrak{c}_{\mathbb{r}}e^{\xi h_{\mathbb{r}}(0)}}$ is finite and the D_h -geodesic from 0 to each point of this set is unique. We now argue that the same holds for $X_{t,s}$ for all $0< t< s<\infty$ simultaneously. Since

 $^{^{3}}$ We write \mathbb{r} for the fixed Euclidean scale with which we are working. The symbol r is used for other, possibly random, radii which arise in the proof.

 D_h induces the Euclidean topology on $\mathbb C$, it is easily seen that $\mathbb F\mapsto \tau_{\mathbb F}$ is continuous and surjective. Indeed, if this function had an upward jump, then there would be some nontrivial interval of times s during which the Euclidean radius of $\mathcal B^{ullet}_s$ does not increase. Therefore, for any given times $0 < t < s < \infty$, we can find a rational $\mathbb F > 0$ and a small rational b > 0 such that $t < \tau_{\mathbb F} < \tau_{\mathbb F} + b c_{\mathbb F} e^{\xi h_{\mathbb F}(0)} < s$. For every leftmost D_h -geodesic P from 0 to $\partial \mathcal B^{ullet}_s$, the restriction $P|_{[0,\tau_{\mathbb F}+bc_{\mathbb F}e^{\xi h_{\mathbb F}(0)}]}$ is a leftmost D_h -geodesic from 0 to $\partial \mathcal B^{ullet}_{\tau+bc_{\mathbb F}e^{\xi h_{\mathbb F}(0)}}$, so $P(\tau_{\mathbb F} + bc_{\mathbb F}e^{\xi h_{\mathbb F}(0)})$ is one of the points $x \in X_{\tau_{\mathbb F},\tau_{\mathbb F}+bc_{\mathbb F}e^{\xi h_{\mathbb F}(0)}}$. Consequently, P coincides with the unique D_h -geodesic P_x from 0 to x until time $\tau_{\mathbb F}$. Therefore, $P|_{[0,t]}$ is equal to one of the finitely many paths $P_x|_{[0,t]}$ for $x \in X_{\tau_{\mathbb F},\tau_{\mathbb F}+bc_{\mathbb F}e^{\xi h_{\mathbb F}(0)}}$. \square

3.1. Outline of the proof. The rest of this section is devoted to the proof of Theorem 3.1. In fact, we will prove a more quantitative version of the theorem (Theorem 3.9 below) which gives bounds for N in terms of t and p, provided we truncate on a certain global regularity event.

We now outline the proof of Theorem 3.1. Fix a stopping time τ for $\{(\mathcal{B}_s^{\bullet}, h|_{\mathcal{B}_s^{\bullet}})\}_{s\geq 0}$ such that diam $(\mathcal{B}_{\tau}^{\bullet})$ is typically of order $\mathbb{F} > 0$, as in Theorem 3.1. We start by considering an arbitrary finite collection \mathcal{I}_0 of disjoint arcs of $\partial \mathcal{B}_{\tau}^{\bullet}$, chosen in a manner depending only on $(\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}})$. We will show that for each $p \in (0, 1)$, there exists $N \in \mathbb{N}$, which does *not* depend on \mathcal{I}_0 or on \mathbb{F} , such that with probability at least p, there are at most N arcs in \mathcal{I}_0 which are hit by a leftmost D_h -geodesic from 0 to $\partial \mathcal{B}_{\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}^{\bullet}$ (we consider leftmost geodesics because of their monotonicity properties, in particular Lemma 2.7). By taking \mathcal{I}_0 to be a large collection of arbitrarily small arcs, this will show that the set of points of $\partial \mathcal{B}_{\tau}^{\bullet}$ which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}^{\bullet}$ can be covered by at most N arbitrarily small arcs with probability at least p, hence, this set is a.s. finite.

To bound the number of arcs in \mathcal{I}_0 which are hit by a leftmost D_h -geodesic from 0 to $\partial \mathcal{B}^{\bullet}_{\tau+tc_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$, we proceed as follows. Let $n:=\#\mathcal{I}_0$. We first show that there exist exponents $\alpha, \theta > 0$ (depending only on the choice of metric) such that if we condition on $(\mathcal{B}^{\bullet}_s, h|_{\mathcal{B}^{\bullet}_s})$ (and truncate on an appropriate global regularity event), then the following is true:

- A. At least 3n/4 of the arcs $I \in \mathcal{I}_0$ can be disconnected from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$ by a set of Euclidean diameter at most $n^{-1/4}$.
- B. For each arc $I \in \mathcal{I}_0$ which can be disconnected from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$ by a set of Euclidean diameter at most $n^{-1/4}$, it holds with conditional probability at least $1 O_n(n^{-\alpha})$ that no leftmost D_h -geodesic from 0 to a point outside of $\mathcal{B}_{\tau+n^{-\theta}}^{\bullet}$ can pass through I.

Property A follows from the deterministic estimate for general planar domains given in Lemma 2.15 and does not require any information at all about the geometry of $\partial \mathcal{B}_{\tau}^{\bullet}$.

Property B is established in Sections 3.2 and 3.3 by using Weyl scaling (Axiom III) and the Markov property of the GFF to show that with positive conditional probability given $(\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}})$, one can build a "shield" around I which no D_h -geodesic from a point far from I to 0 can pass through. This shield will consist of two concentric annuli of the form $\mathbb{A}_{2r,3r}(z)$ and $\mathbb{A}_{3r,4r}(z)$ for appropriate choices of r > 0 and $z \in \mathbb{C}$ with the following properties. The annulus $\mathbb{A}_{2r,3r}(z)$ disconnects I from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$ and the D_h -distance from any point of $\mathbb{A}_{3r,4r}(z)$ to $\partial \mathcal{B}_{\tau}^{\bullet}$ is smaller than the D_h -distance between the inner and outer boundaries of $\mathbb{A}_{2r,3r}(z)$. Such annuli are illustrated in Figure 7. It is easy to see from the definition of a geodesic (and is explained carefully in the proof of Lemma 3.6) that if such annuli exist, then no D_h -geodesic from 0 to a point outside of $\mathcal{B}_{\tau}^{\bullet} \cup B_{4r}(z)$ can hit I.

We will set things up so that one has a logarithmic number of essentially independent chances to build a shield of the above form, so the probability that no such shield exists decays like a negative power of n. The events needed to build the shield depend on the metric in a

reasonably continuous way, so Axiom V allows us to get a lower bound for the probability that the shield exists which is uniform in \mathbb{F} (see Lemma 3.2). There is some subtlety here since the event that the shield exists depends on both the zero-boundary part of the field outside of $\mathcal{B}^{\bullet}_{\tau}$ and the pair $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$. To deal with this, we will use Lemma 2.12 to find many annuli surrounding I where the harmonic part of $h|_{\mathbb{C}\setminus\mathcal{B}^{\bullet}_{\tau}}$ is under control and only try to build a shield in these "good" annuli (see Lemma 3.2). We will also need to use the FKG inequality (Proposition 2.8) to deal with the behavior of h very close to $\partial \mathcal{B}^{\bullet}_{\tau}$. See Figure 7 for an illustration of this part of the argument.

We will then apply properties A and B iteratively to "kill off" all of the geodesics from all but a constant order number of the arcs in \mathcal{I}_0 . This is done in Section 3.4 and is illustrated in Figure 8. To this end, we define radii s_k and collections of boundary arcs \mathcal{I}_k for $k \in \mathbb{N}$ inductively as follows. We set $s_0 = \tau$, and we define s_k to be a stopping time (to be specified precisely below) which is between s_{k-1} and $s_{k-1} + (\#\mathcal{I}_{k-1})^{-\theta} c_{\mathbb{I}^r} e^{\xi h_{\mathbb{I}^r}(0)}$. For $I \in \mathcal{I}_{k-1}$, we define I' to be the set of points in $\partial \mathcal{B}_{s_k}^{\bullet}$ which are hit by a leftmost D_h -geodesic from 0 which passes through I, and we set $\mathcal{I}_k := \{I' : I \in \mathcal{I}_{k-1}, I' \neq \varnothing\}$. Basic properties of geodesics (see Lemma 2.7) show that \mathcal{I}_k is a collection of disjoint arcs of $\partial \mathcal{B}_{s_k}^{\bullet}$. By applying properties A and B above with s_k in place of τ , we get that typically $\#\mathcal{I}_{k+1} \le \frac{1}{2}\#\mathcal{I}_k$ (Lemma 3.10). Using this, we infer that if N is chosen sufficiently large, in a manner which does not depend on \mathcal{I}_0 , and K is the smallest $k \in \mathbb{N}$ for which $\#\mathcal{I}_k \le N$, then which high probability $s_K \le \tau + t c_{\mathbb{I}^r} e^{\xi h_{\mathbb{I}^r}(0)}$ (Lemma 3.11). We do not prove any bound for K, just for s_K . This shows that with high probability, there are at most N arcs in \mathcal{I}_0 which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau+tc_{\mathbb{I}^r}e^{\xi h_{\mathbb{I}^r}(0)}$. Since the arcs in \mathcal{I}_0 can be made arbitrarily small and N does not depend on \mathcal{I}_0 , this implies Theorem 3.1.

3.2. *Good annuli*. We now define an event for a Euclidean annulus which will eventually be used to build "shields" surrounding boundary arcs of a filled D_h -metric ball through which D_h -geodesics to 0 cannot pass. See Figure 6 for an illustration.

For $\epsilon > 0$, $z \in \mathbb{C}$ and a set $V \subset \mathbb{C}$, we define the collection of Euclidean squares

$$(34) \ \mathcal{S}^{z}_{\epsilon}(V) := \{ [x, x + \epsilon] \times [y, y + \epsilon] : (x, y) \in \epsilon \mathbb{Z}^2 + z, ([x, x + \epsilon] \times [y, y + \epsilon]) \cap V \neq \emptyset \}.$$

Note that $\mathcal{S}^z_{\epsilon}(V)$ depends only on the value of z modulo $\epsilon \mathbb{Z}^2$ and that $\mathcal{S}^z_{\epsilon}(V) - z = \mathcal{S}^0_{\epsilon}(V - z)$.

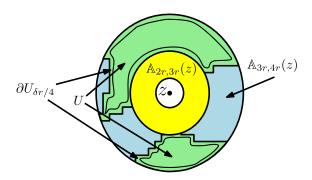


FIG. 6. Illustration of the definitions in Section 3.2. The set $U_r(z) = U_r(z; \delta)$ consists of open subsets U of $\mathbb{A}_{3r,4r}(z)$ such that $\mathbb{A}_{3r,4r}(z) \setminus U$ is a finite union of sets of the form $S \cap \mathbb{A}_{3r,4r}(z)$ for $\delta r \times \delta r$ squares $S \in \mathcal{S}_{\delta r}(\mathbb{A}_{3r,4r}(z))$ (i.e., with corners in $\delta \mathbb{Z}^2$). One such set is shown in light green. For each $U \in \mathcal{U}_r(z)$, $E_r^U(z)$ is the event that: (1) the D_h -distance across the yellow annulus $\mathbb{A}_{2r,3r}(z)$ is bounded below, (2) the D_h -diameter of each of the squares $S \in \mathcal{S}_{\delta r}(\mathbb{A}_{3r,4r}(z)) + z$ is much smaller than the D_h -distance across $\mathbb{A}_{2r,3r}(z)$ and (3) the harmonic part of $h|_U$ is bounded above on the set $U_{\delta r/4} \subset U$ (outlined in black).

For $z \in \mathbb{C}$, r > 0 and $\delta \in (0, 1)$, we define $\mathcal{U}_r(z) = \mathcal{U}_r(z; \delta)$ to be the (finite) set of open subsets U of the annulus $\mathbb{A}_{3r,4r}(z)$ such that $\mathbb{A}_{3r,4r}(z) \setminus U$ is a finite union of sets of the form $S \cap \mathbb{A}_{3r,4r}(z)$ for squares $S \in \mathcal{S}^z_{\delta r}(\mathbb{A}_{3r,4r}(z))$. For $U \in \mathcal{U}_r(z; \delta)$ and $\epsilon > 0$, we define

(35)
$$U_{\epsilon} := \{ u \in U : \operatorname{dist}(z, \partial U) > \epsilon \},$$

where dist denotes Euclidean distance.

For $z \in \mathbb{C}$, r > 0, parameters δ , $c \in (0, 1)$ and A > 0, and $U \in \mathcal{U}_r(z; \delta)$, we let $E_r^U(z) = E_r^U(z; c, \delta, A)$ be the event that the following is true:

- 1. $D_h(\partial B_{2r}(z), \partial B_{3r}(z)) \ge c \mathfrak{c}_r e^{\xi h_r(z)}$.
- 2. One has

(36)
$$\max_{S \in \mathcal{S}_{\delta r}^{z}(\mathbb{A}_{3r,4r}(z))} \sup_{u,v \in S} D_{h}(u,v;\mathbb{A}_{2r,5r}(z)) \leq \frac{c}{100} \mathfrak{c}_{r} e^{\xi h_{r}(z)}.$$

3. Let \mathfrak{h}^U be the harmonic part of $h|_U$. Then, in the notation (35),

(37)
$$\sup_{u \in U_{\delta r/4}} \left| \mathfrak{h}^U(u) - h_r(z) \right| \le A.$$

We also define

(38)
$$E_r(z) = E_r(z; c, \delta, A) := \bigcap_{U \in \mathcal{U}_r(z; \delta)} E_r^U(z).$$

The first two conditions in the definition of $E_r^U(z)$ do not depend on U, so the only difference between $E_r(z)$ and $E_r^U(z)$ is that, for the former event, condition 3 is required to hold for all choices of U simultaneously.

We think of annuli $\mathbb{A}_{2r,5r}(z)$ for which $E_r(z)$ occurs as "good." We will show in Lemma 3.2 just below that $\mathbb{P}[E_r(z)]$ can be made close to 1 by choosing the parameters δ , c, A appropriately, in a manner which is uniform over the choices of r and z, The reason for separating $E_r(z)$ and $E_r^U(z)$ is that conditioning on $E_r^U(z)$ is easier than conditioning on $E_r(z)$ (see Lemma 3.3 just below).

The occurrence of $E_r^U(z)$ or $E_r(z)$ is unaffected by adding a constant to the field. By this and the locality of D_h (Axiom II), these events are determined by $h|_{\mathbb{A}_{2r,5r}(z)}$, viewed modulo additive constant. This will be important when we apply Lemma 2.12 below.

We will eventually apply condition 3 with U equal to $\mathbb{A}_{3r,4r}(z)$ minus the union of the set of squares in $\mathcal{S}^z_{\epsilon}(\mathbb{A}_{3r,4r}(z))$ which intersect a filled D_h -metric ball centered at 0. Condition 3, together with the Markov property of h, will allow us to show that with uniformly positive conditional probability given $h|_{\mathbb{C}\setminus U}$ and the event $E^U_r(z)$, the D_h -diameter of $U_{\delta r/4}$ is small (see Lemma 3.3). This combined with condition 2 will show that with uniformly positive conditional probability given $h|_{\mathbb{C}\setminus U}$ and $E^U_r(z)$, every point of $\mathbb{A}_{3r,4r}(z)$ lies at D_h -distance strictly smaller than $cr^{\xi}Qe^{\xi h_r(z)}$ from the filled D_h -metric ball. Due to condition 1, this will prevent a D_h -geodesic from crossing $\mathbb{A}_{2r,3r}(z)$ before entering this filled metric ball; see Figure 7 for an illustration of how the events $E^U_r(z)$ will eventually be used.

LEMMA 3.2. For each $p \in (0, 1)$, we can find parameters $c, \delta \in (0, 1)$ and A > 0 such that, in the notation (38), we have $\mathbb{P}[E_r(z)] \ge p$ for each $z \in \mathbb{C}$ and r > 0.

PROOF. By translation invariance and tightness across scales (Axioms IV and V), the laws of the reciprocals of the scaled distances $\mathfrak{c}_r^{-1}e^{-\xi h_r(z)}D_h(\partial B_{2r}(z),\partial B_{3r}(z))$ for $z\in\mathbb{C}$ and r>0 are tight. Therefore, we can find c=c(p)>0 such that for each $z\in\mathbb{C}$ and r>0, condition 1 in the definition of $E_r^U(z)$ occurs with probability at least 1-(1-p)/3. Similarly, Axioms IV and V show that we can find $\delta=\delta(p,c)\in(0,1)$ such that condition 2 in

the definition of $E_r(z)$ occurs with probability at least 1-(1-p)/3. For a given choice of δ , the collection of open sets $\mathcal{U}_r(z;\delta)$ is finite and is equal to $r\mathcal{U}_1(0;\delta)+z$ (here we use the translation by z in (34)). Since \mathfrak{h}^U is continuous away from ∂U , for any fixed choice of $U \in \mathcal{U}_1(0;\delta)$, a.s. $\sup_{u \in U_{\delta/4}} |\mathfrak{h}^U(u)| < \infty$. By combining this with the translation and scale invariance of the law of h, modulo additive constant, we find that there exists A > 0 (depending on δ) such that with probability at least 1-(1-p)/3, condition 3 in the definition of $E_r^U(z)$ holds simultaneously for every $U \in \mathcal{U}_r(z;\delta)$. \square

LEMMA 3.3. For any choice of parameters c, δ, A , there is a constant $\mathfrak{p} = \mathfrak{p}(c, \delta, A) > 0$ such that the following is true. Let r > 0, let $z \in \mathbb{C}$, and let $U \in \mathcal{U}_r(z) = \mathcal{U}_r(z; \delta)$. Also, let $\mathcal{V}(U) \subset \mathcal{U}_r(z)$ be the set of connected components of U. Almost surely,

(39)
$$\mathbb{P}\left[\max_{V\in\mathcal{V}(U)}\sup_{u,v\in V}D_{h}(u,v;\mathbb{A}_{2r,5r}(z))\leq \frac{c}{2}\mathfrak{c}_{r}e^{\xi h_{r}(z)}|h|_{\mathbb{C}\setminus U},E_{r}^{U}(z)\right]\geq \mathfrak{p}.$$

PROOF. By Axiom III, the statement of the lemma does not depend on the choice of additive constant for h, so we can assume without loss of generality that h is normalized so that $h_{r_0}(z_0) = 0$ for some $z_0 \in \mathbb{C}$ and $r_0 > 0$ such that $B_{r_0}(z_0) \cap U = \emptyset$ (we could take $z_0 = z$ and $r_0 = r$, but we find that the proof is clearer if we do not normalize so that $h_r(z) = 0$). For $U \in \mathcal{U}_r(z)$, let \mathring{h}^U be the zero-boundary part of $h|_U$. By the Markov property of the field (see, e.g., [7], Lemma 2.2), we can write $h|_U = \mathring{h}^U + \mathfrak{h}^U$, where \mathring{h}^U is a zero-boundary GFF in U which is independent from $h|_{\mathbb{C}\setminus U}$ (and hence also from the harmonic part \mathfrak{h}^U). We define the metric $D_{\mathring{h}U}$ as in Remark 1.2.

Since conditions 1 and 3 in the definition of $E_r^U(z)$ are determined by $h|_{\mathbb{C}\setminus U}$, the conditional law of \mathring{h}^U given $h|_{\mathbb{C}\setminus U}$ and the event $E_r^U(z)$ is the same as the conditional law of \mathring{h}^U given \mathfrak{h}^U and the event $F_r(z)$ that condition 2 in the definition of $E_r^U(z)$ occurs (note that this condition does not depend on U).

The main idea of the proof is as follows. By Axiom III, we have $D_{h|U} = e^{\xi \mathfrak{h}^U} \cdot D_{\mathring{h}^U}$. Condition 3 in the definition of $E_r^U(z)$ gives an upper bound for \mathfrak{h}^U , so we just need to prove that certain $D_{\mathring{h}^U}$ -distances are very small with positive conditional probability when we condition on $h|_{\mathbb{C}\setminus U}$ and $E_r^U(z)$ (see (40) for the precise event we need). Since \mathring{h}^U is independent from $h|_{\mathbb{C}\setminus U}$ and by the preceding paragraph, we only need to prove a bound for $D_{\mathring{h}^U}$ -distances when we condition on $F_r(z)$.

This can be done as follows. By Axiom III, if we let f be a large bump function supported on a compact subset of U which is close to all of U, then $D_{\mathring{h}^U-f}$ distances are much smaller than $D_{\mathring{h}^U}$ distances. On the other hand, the laws of \mathring{h}^U and \mathring{h}^U-f are mutually absolutely continuous. This shows that $D_{\mathring{h}^U}$ -distances are small with positive probability under the unconditional law of \mathring{h}^U . This last estimate can be made uniform over the choice of $U \in \mathcal{U}_r(z)$ since $\mathcal{U}_r(z) = r\mathcal{U}_1(0) + z$ and $\mathcal{U}_1(0)$ is a finite set (this is why we restrict to domains U which are made up of small squares in a fixed grid). To add in the conditioning on $F_r(z)$, we use the FKG inequality (Proposition 2.8):

Step 1: Reducing to an event for the zero-boundary part. For $U \in \mathcal{U}_r(z)$, define $V_{\delta r/2}$ and $V_{\delta r/4}$ for $V \in \mathcal{V}(U)$ as in (35) but with V in place of U. Note that $U_{\delta r/2} = \bigcup_{V \in \mathcal{V}(U)} V_{\delta r/2}$. Let

(40)
$$G^{U} := \left\{ \max_{V \in \mathcal{V}(U)} \sup_{u, v \in V_{\delta_{r}/2}} D_{\mathring{h}^{U}}(u, v; V_{\delta_{r}/4}) \le \frac{c}{100} e^{-\xi A} \mathfrak{c}_{r} \right\}.$$

By condition 3 in the definition of $E_r^U(z)$ and Axiom III, if $E_r^U(z) \cap G^U$ occurs, then

(41)
$$\max_{V \in \mathcal{V}(U)} \sup_{u,v \in V_{\delta r/2}} D_h(u,v;V_{\delta r/4}) \le \frac{c}{100} \mathfrak{c}_r e^{\xi h_r(z)}.$$

Since any two points of any $V \in \mathcal{V}(U)$ can be joined by a path in V, which is contained in the union of $V_{\delta r/2}$ and at most two squares $S \in \mathcal{S}^z_{\delta r}(\mathbb{A}_{2r,4r}(z))$, combining (41) with condition 2 in the definition of $E^U_r(z)$ shows that if $E^U_r(z) \cap G^U$ occurs, then the event in (39) occurs.

To prove the lemma, it therefore suffices to find $\mathfrak{p} = \mathfrak{p}(c, \delta, A) > 0$ such that a.s. $\mathbb{P}[G^U | h|_{\mathbb{C}\setminus U}, E_r^U(z)] \geq \mathfrak{p}$ for each $U \in \mathcal{U}_r(z)$. By the above discussion about $F_r(z)$ and since G^U is determined by \mathring{h}^U , so is independent from $h|_{\mathbb{C}\setminus U}$, we only need to show that a.s.

(42)
$$\mathbb{P}\left[G^{U}|F_{r}(z)\right] \geq \mathfrak{p} \quad \forall U \in \mathcal{U}_{r}(z).$$

Step 2: Lower bound without conditioning on $F_r(z)$. We first argue that there is a constant $\mathfrak{p} = \mathfrak{p}(c, \delta, A) > 0$ such that a.s.

(43)
$$\mathbb{P}\left[G^{U}\right] \geq \mathfrak{p} \quad \forall U \in \mathcal{U}_{r}(z).$$

By Axioms IV and V and since $U_r(z) = rU_1(0) + z$ and $\#U_1(0)$ depends only on δ , we can find a constant C > 0, depending only on δ , such that for each $z \in \mathbb{C}$, r > 0, and $U \in U_r(z)$,

(44)
$$\mathbb{P}\left[\max_{V\in\mathcal{V}(U)}\sup_{u,v\in\mathcal{V}_{\delta_{r}/2}}D_{\mathring{h}^{U}}(u,v;V_{\delta_{r}/4})\leq C\mathfrak{c}_{r}\right]\geq\frac{1}{2}.$$

Let $U^0 := r^{-1}(U-z) \in \mathcal{U}_1(0)$, and let $g^0 : U^0 \to [0,1]$ be a smooth compactly supported bump function on U^0 which is identically equal to 1 on $U^0_{\delta/4}$. Let $g:U\to [0,1]$ be defined by $g(\cdot)=g^0(r^{-1}(\cdot-z))$. Then, the Dirichlet energy of g equals the Dirichlet energy of g^0 which depends only on U^0 . Let $f:=\xi^{-1}(\log C - \log(ce^{-\xi A}/100))g$. By Axiom III, if the event in (44) occurs, then G^U occurs with \mathring{h}^U-f in place of \mathring{h}^U .

By a standard calculation for the GFF, the laws of \mathring{h}^U and $\mathring{h}^U - f$ are mutually absolutely continuous and the law of the Radon–Nikodym derivative depends only on the Dirichlet energy of f which in turn depends only on U^0 , δ , c, A. It follows that $\mathbb{P}[G^U]$ is bounded below by a constant depending only on U^0 , δ , c, A. Since the number of possibilities for U^0 depends only on δ , by taking the minimum over all such possibilities we get (43) for an appropriate choice of \mathfrak{p} .

Step 3: Adding the conditioning on $F_r(z)$. We now use the FKG inequality to add in the conditioning on $F_r(z)$. Indeed, under the conditional law given $h|_{\mathbb{C}\setminus U}$, both $\mathbb{1}_{G^U}$ and $\mathbb{1}_{F_r(z)}$ are nonincreasing functions of the metric $D_{\hat{h}^U}$ (note that $\mathbb{1}_{F_r(z)}$ depends also on $h|_{\mathbb{C}\setminus U}$, but we can still view it as a function of \hat{h}^U when we condition on a fixed realization of $h|_{\mathbb{C}\setminus U}$). Moreover, it is easily seen that these functions are a.s. continuous at $D_{\hat{h}^U}$ in the sense of Proposition 2.8. In the case of G^U , this follows since the probability that the supremum in (40) is exactly equal to $\frac{c}{100}e^{-\xi A}c_r$ is zero. A similar justification holds for $F_r(z)$. By Proposition 2.8, the events $F_r(z)$ and G^U are positively correlated under the conditional law given $h|_{\mathbb{C}\setminus U}$. Therefore, (43) implies that (42) holds. \square

3.3. Cutting off geodesics from a boundary arc. We will now use the events of the preceding subsection to build "shields" which prevent D_h -geodesics from hitting a given arc of a filled metric ball. Fix parameters $c, \delta \in (0, 1)$ and A > 0, and define $E_r(z) = E_r(z; c, \delta, A)$ as in (38) For $z \in \mathbb{C}$ and $\mathbb{F} > 0$, let $\rho_{\mathbb{F}}^0(z) := \mathbb{F}$ and for $n \in \mathbb{N}$; inductively define

(45)
$$\rho_{\mathbb{F}}^{n}(z) := \inf\{r \ge 6\rho_{\mathbb{F}}^{n-1}(z) : r = 2^{k} \mathbb{F} \text{ for some } k \in \mathbb{Z}, E_{r}(z) \text{ occurs}\}.$$

Since $E_r(z)$ is determined by $h|_{\mathbb{A}_{2r,5r}(z)}$, it follows that $\rho_{\mathbb{T}}^n(z)$ is a stopping time for the filtration generated by $h|_{B_{5r}(z)}$ for $r \geq \mathbb{T}$.

LEMMA 3.4. For each q > 0, we can find parameters $c, \delta \in (0, 1)$ and A > 0 and another parameter $\eta > 0$, all depending on q, such that the following is true. Uniformly over all r > 0 and $z \in \mathbb{C}$,

(46)
$$\mathbb{P}\left[\rho_{\mathbb{F}}^{\lfloor \eta \log C \rfloor}(z) > C\mathbb{F}\right] = O_C(C^{-q}) \quad as \ C \to \infty.$$

PROOF. Since each $E_r(z)$ is determined by $(h - h_r(z))|_{\mathbb{A}_{2r,5r}(z)}$, this follows by combining Lemma 2.12 (applied with $r_k = 8^k \mathbb{F}$, say, $S_1 = 2$, $S_2 = 5$, $a = \log q$, b = 1/2, say, and $K = \lfloor \log_6 C \rfloor$) and Lemma 3.2. \square

LEMMA 3.5. There exists a choice of parameters $c, \delta \in (0, 1)$ and A > 0 and another parameter $\eta > 0$, depending only on the choice of metric D, such that the following is true. For each compact set $K \subset \mathbb{C}$, it holds with probability $1 - O_{\epsilon}(\epsilon^2)$ (at a rate depending on K) that

(47)
$$\rho_{\epsilon \mathbb{r}}^{\lfloor \eta \log \epsilon^{-1} \rfloor}(z) \le \epsilon^{1/2} \mathbb{r} \quad \forall z \in \left(\frac{\epsilon \mathbb{r}}{4} \mathbb{Z}^2\right) \cap B_{\epsilon \mathbb{r}}(\mathbb{r} K).$$

PROOF. This follows from Lemma 3.4 (applied $\epsilon \mathbb{r}$ in place of \mathbb{r} , with $C = \epsilon^{-3/2}$ and with q = 6, say) and a union bound over $O_{\epsilon}(\epsilon^{-2})$ points in $(\frac{\epsilon \mathbb{r}}{4}\mathbb{Z}^2) \cap B_{\epsilon \mathbb{r}}(\mathbb{r}K)$. \square

Henceforth, assume that c, δ , A, and η are as in Lemma 3.5. For $\epsilon > 0$, $\mathbb{r} > 0$ and a compact set $K \subset \mathbb{C}$, let

$$(48) R_{\mathbb{F}}^{\epsilon}(K) := 6 \sup \left\{ \rho_{\epsilon_{\mathbb{F}}}^{\lfloor \eta \log \epsilon^{-1} \rfloor}(z) : z \in \left(\frac{\epsilon_{\mathbb{F}}}{4} \mathbb{Z}^{2} \right) \cap B_{\epsilon_{\mathbb{F}}}(K) \right\} + \epsilon_{\mathbb{F}},$$

so that each of the radii $\rho_{\epsilon_{\mathbb{I}}}^{n}(z)$ for $z \in (\frac{\epsilon_{\mathbb{I}}}{4}\mathbb{Z}^{2}) \cap B_{\epsilon_{\mathbb{I}}}(K)$ and $n \in [1, \eta \log \epsilon^{-1}]_{\mathbb{Z}}$ is determined by $R_{\mathbb{I}}^{\epsilon}(K)$ and $h|_{B_{R_{\mathbb{I}}^{\epsilon}(K)}(K)}$. Lemma 3.5 shows that for each fixed choice of K, $\mathbb{P}[R_{\mathbb{I}}^{\epsilon}(rK) \leq (6\epsilon^{1/2} + \epsilon)\mathbb{I}]$ tends to 1 as $\epsilon \to 0$, at a rate which is uniform in \mathbb{I} .

Recall from Section 2.1 that \mathcal{B}_s^{\bullet} for s > 0 denotes the filled D_h -ball of radius s centered at zero. For s > 0, define

(49)
$$\sigma_{s,\mathbb{I}}^{\epsilon} := \inf\{s' > s : B_{R_{\mathbb{I}}^{\epsilon}(\mathcal{B}_{s}^{\bullet})}(\mathcal{B}_{s}^{\bullet}) \subset \mathcal{B}_{s'}^{\bullet}\},$$

so that $\mathcal{B}^{\bullet}_{\sigma^{\epsilon}_{s,\mathbb{I}}}$ contains $B_{6\rho^{\lfloor\eta\log\epsilon^{-1}\rfloor}(z)}(z)$ for each $z\in B_{\epsilon\mathbb{I}}(\mathcal{B}^{\bullet}_{s})$. Since each $\rho^{\lfloor\eta\log\epsilon^{-1}\rfloor}_{\epsilon\mathbb{I}}(z)$ is a stopping time for the filtration generated by $h|_{B_{5r}(z)}$ for $r\geq\epsilon\mathbb{I}$, it follows that if τ is a stopping time for $\{(\mathcal{B}^{\bullet}_{t},h|_{\mathcal{B}^{\bullet}_{t}})\}_{t\geq0}$, then so is $\sigma^{\epsilon}_{\tau,\mathbb{I}}$. The following lemma will be used to "kill off" the D_h -geodesics from 0 which hit a given boundary arc of a filled D_h -metric ball.

LEMMA 3.6. There exists $\alpha > 0$, depending only on the choice of metric, such that the following is true. Let $\mathbb{r} > 0$, let τ be a stopping time for the filtration generated by $\{(\mathcal{B}_{\mathfrak{s}}^{\bullet}, h|_{\mathcal{B}_{\mathfrak{s}}^{\bullet}})\}_{s\geq 0}$ and let $x\in\partial\mathcal{B}_{\tau}^{\bullet}$ and $\epsilon\in(0,1)$ be chosen in a manner depending only on $(\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}})$. There is an event $G_{x}^{\epsilon}\in\sigma(\mathcal{B}_{\sigma_{\tau,\pi}^{\epsilon}}^{\bullet}, h|_{\mathcal{B}_{\sigma_{\tau,\pi}^{\epsilon}}^{\bullet}})$ with the following properties:

A. If, in the notation (48), we have $R_{\mathbb{I}}^{\epsilon}(\mathcal{B}_{\tau}^{\bullet}) \leq \dim \mathcal{B}_{\tau}^{\bullet}$ and G_{x}^{ϵ} occurs, then no D_{h} -geodesic from 0 to a point in $\mathbb{C} \setminus B_{R_{\mathbb{I}}^{\epsilon}(\mathcal{B}_{\tau}^{\bullet})}(\mathcal{B}_{\tau}^{\bullet})$ can enter $B_{\epsilon_{\mathbb{I}}}(x) \setminus \mathcal{B}_{\tau}^{\bullet}$.

B. There is a deterministic constant $C_0 > 1$ depending only on the choice of metric such that a.s. $\mathbb{P}[G_x^{\epsilon}|\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}}] \geq 1 - C_0 \epsilon^{\alpha}$.

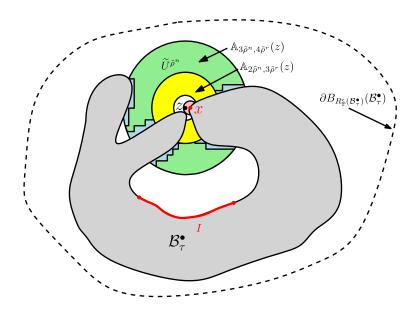


FIG. 7. Illustration of the proof of Lemma 3.6. The point $z \in \frac{\epsilon \pi}{4} \mathbb{Z}^2$ is chosen so that $B_{\epsilon \pi}(x) \subset B_{2\epsilon \pi}(z)$. On the event G_x^{ϵ} defined in (53), there is some $n \in [1, \eta \log \epsilon^{-1}]_{\mathbb{Z}}$ for which the following is true. With $\widetilde{\rho}^n$ as in (51), the sum of the D_h -diameter of the light green region $\widetilde{U}^{\widetilde{\rho}^n}$ and the maximal D_h -diameters of the light blue squares $S \in \mathcal{S}_{\delta\widetilde{\rho}^n}^z(\mathbb{A}_{2\widetilde{\rho}^n},4\widetilde{\rho}^n(z))$ is strictly smaller than the D_h -distance across the yellow annulus $\mathbb{A}_{2\widetilde{\rho}^n},3\widetilde{\rho}^n(z)$. This makes it so that each point in $\mathbb{A}_{3\widetilde{\rho}^n},4\widetilde{\rho}^n(z)$ is D_h -closer to a point of $\mathcal{B}_{\tau}^{\bullet} \cap \mathbb{A}_{3\widetilde{\rho}^n},4\widetilde{\rho}^n(z)$ than it is to $B_{\epsilon \pi}(x)$. Hence, no D_h -geodesic from a point outside of $B_{R_{\tau}^n}(\mathcal{B}_{\tau}^{\bullet})(\mathcal{B}_{\tau}^{\bullet})$ to 0 can enter $B_{\epsilon \pi}(x) \setminus \mathcal{B}_{\tau}^{\bullet}$. The condition that $R_{\tau}^n(\mathcal{B}_{\tau}^{\bullet}) \leq \dim \mathcal{B}_{\tau}^{\bullet}$ ensures that $\mathbb{A}_{3\widetilde{\rho}^n},4\widetilde{\rho}^n(z)$ intersects $\mathcal{B}_{\tau}^{\bullet}$. We can also prevent D_h -geodesics from hitting an arc I of $\partial \mathcal{B}_{\tau}^{\bullet}$ by choosing x so that $B_{\epsilon \pi}(x)$ disconnects I from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$; see Lemma 3.7.

PROOF. See Figure 7 for an illustration of the proof. We will outline the argument just below, after introducing some notation:

Step 1: Setup. We can choose $z \in (\frac{\epsilon \mathbb{F}}{4}\mathbb{Z}^2) \cap B_{\epsilon \mathbb{F}}(\mathcal{B}^{\bullet}_{\tau})$ such that $B_{\epsilon \mathbb{F}}(x) \subset B_{2\epsilon \mathbb{F}}(z)$, in a manner depending only on $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$. Recalling the set of squares $\mathcal{S}^{z}_{\delta r}(\cdot)$ from (34), for r > 0 we define

(50)
$$\widetilde{U}^r := \widetilde{U}^r(z) := \mathbb{A}_{3r,4r}(z) \setminus \bigcup \{ S \in \mathcal{S}^z_{\delta r}(\mathbb{A}_{3r,4r}(z)) : S \cap \mathcal{B}^{\bullet}_{\tau} \neq \emptyset \}.$$

Note that \widetilde{U}^r belongs to the set $\mathcal{U}_r(z)$ of Section 3.2 and \widetilde{U}^r is determined by $(\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}})$. Let $\widetilde{\rho}^0 := \epsilon \mathbb{r}$ and for $n \in \mathbb{N}$, inductively define

(51)
$$\widetilde{\rho}^n = \widetilde{\rho}_{\varepsilon \mathbb{F}}^n(z) := \inf\{r \ge 6\widetilde{\rho}^{n-1} : r = 2^k \mathbb{F} \text{ for some } k \in \mathbb{Z}, E_r^{\widetilde{U}^r}(z) \text{ occurs}\}.$$

In other words, $\widetilde{\rho}^n$ is defined in the same manner as $\rho^n_{\epsilon \mathbb{T}}(z)$ from (45) (with $\epsilon \mathbb{T}$ in place of \mathbb{T}) but with $E_r^{\widetilde{U}^r}(z)$ instead of $E_r(z)$. This means that $E_{\widetilde{\rho}^n}^U(z)$ is only required to occur for $U = \widetilde{U}^{\widetilde{\rho}^n}$ instead of for every $U \in \mathcal{U}_{\widetilde{\rho}^n}(z)$. By this and the definition (48) of $R_{\mathbb{T}}^{\epsilon}(\mathcal{B}^{\bullet}_{\tau})$,

(52)
$$\widetilde{\rho}^n \leq \rho_{\epsilon \mathbb{F}}^n(z) \quad \forall n \in \mathbb{N}_0 \quad \text{and, hence} \quad \widetilde{\rho}^{\lfloor \eta \log \epsilon^{-1} \rfloor} \leq \frac{1}{6} R_{\mathbb{F}}^{\epsilon} (\mathcal{B}_{\tau}^{\bullet}).$$

The reason for considering $\tilde{\rho}^n$ instead of $\rho_{\epsilon \mathbb{F}}^n(z)$ is because we can only condition on $E_r^U(z)$, not on $E_r(z)$, in Lemma 3.3.

Recalling that $\mathcal{V}(\widetilde{U}^{\widetilde{\rho}^n})$ denotes the set of connected components of $\widetilde{U}^{\widetilde{\rho}^n}$, we define

(53)
$$G_{x}^{\epsilon} := \left\{ \exists n \in [1, \eta \log \epsilon^{-1}]_{\mathbb{Z}} \text{ s.t.} \right.$$

$$\max_{V \in \mathcal{V}(\widetilde{U}\widetilde{\rho}^{n})} \sup_{u, v \in V} D_{h}(u, v; \mathbb{A}_{2\widetilde{\rho}^{n}, 5\widetilde{\rho}^{n}}(z)) \leq \frac{c}{2} \mathfrak{c}_{\widetilde{\rho}^{n}} e^{\xi h_{\widetilde{\rho}^{n}}(z)} \right\}.$$

In other words, G_x^{ϵ} is the event that the event of Lemma 3.3 occurs for at least one of the sets $\widetilde{U}^{\widetilde{\rho}^n}$ for $n \in [1, \eta \log \epsilon^{-1}]_{\mathbb{Z}}$:

- To check property A, we will show that if G_{χ}^{ϵ} occurs and n is as in the definition of G_{χ}^{ϵ} , then no D_h -geodesic from a point outside of $\mathbb{C} \setminus B_{R_{\mathbb{T}}^{\epsilon}(\mathcal{B}_{\tau}^{\bullet})}(\mathcal{B}_{\tau}^{\bullet})$ to 0 can cross between the inner and outer boundaries of the annulus $\mathbb{A}_{2\widetilde{\rho}^n,4\widetilde{\rho}^n}(z)$. As we will explain in Step 2 below, the reason for this is that the D_h -distance from any point of $\mathbb{A}_{3\widetilde{\rho}^n,4\widetilde{\rho}^n}(z)$ to $\mathcal{B}_{\tau}^{\bullet}$ is shorter than the D_h -distance across $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$, so it is more efficient to enter $\mathcal{B}_{\tau}^{\bullet}$ before crossing $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$.
- To check property B, we will first apply Lemma 3.3 for a possible realization of $\widetilde{U}^{\widetilde{\rho}^n}$ to get that the conditional probability of G_x^{ϵ} given $z, \epsilon, \widetilde{\rho}^n, \widetilde{U}^{\widetilde{\rho}^n}$, and $h|_{\mathbb{C}\setminus\widetilde{U}\widetilde{\rho}^n}$ is at least \mathfrak{p} . We will then multiply this estimate over all $n \in [1, \eta \log \epsilon^{-1}]$ to get $\mathbb{P}[G_x^{\epsilon}] \ge 1 C_0 \epsilon^{\alpha}$ for α slightly less than $\eta \log(1/(1-\mathfrak{p}))$ and an appropriate choice of $C_0 > 1$.

Since z and \widetilde{U}^r for r>0 are each determined by $(\mathcal{B}^{\bullet}_{\tau},h|_{\mathcal{B}^{\bullet}_{\tau}})$, it follows that each $E^{\widetilde{U}^r}(z)$ is determined by $(\mathcal{B}^{\bullet}_{\tau},h|_{\mathcal{B}^{\bullet}_{\tau}})$ and $h|_{\mathbb{A}_{2r,5r}(z)}$. Hence, $\widetilde{\rho}^n$ is a stopping time for the filtration generated by $h|_{\mathcal{B}_{5r}(z)}$ for $r\geq \epsilon \mathbb{F}$ and $(\mathcal{B}^{\bullet}_{\tau},h|_{\mathcal{B}^{\bullet}_{\tau}})$. By (52) and the definition (49) of $\sigma^{\epsilon}_{\tau,\mathbb{F}}$, we have $B_{5\widetilde{\rho}^n}(z)\subset \mathcal{B}^{\bullet}_{\sigma^{\epsilon}_{\tau,\mathbb{F}}}$. By combining these statements with (53) and the locality of the metric (Axiom II), we get that $G_{\chi}^{\epsilon}\in\sigma(\mathcal{B}^{\bullet}_{\sigma^{\epsilon}_{\tau,\mathbb{F}}},h|_{\mathcal{B}^{\bullet}_{\sigma^{\epsilon}_{\tau,\mathbb{F}}}})$.

Step 2: Proof that G_x^{ϵ} satisfies property A. Assume that $R_{\mathbb{T}}^{\epsilon}(\mathcal{B}_{\tau}^{\bullet}) \leq \dim \mathcal{B}_{\tau}^{\bullet}$ and G_x^{ϵ} occurs. Choose $n \in [1, \eta \log \epsilon^{-1}]_{\mathbb{Z}}$ as in the definition (53) of G_x^{ϵ} . Then,

$$\epsilon \mathbb{r} \leq \widetilde{\rho}^n \leq \frac{1}{6} R_{\mathbb{r}}^{\epsilon} (\mathcal{B}_{\tau}^{\bullet}) \leq \frac{1}{6} \operatorname{diam} \mathcal{B}_{\tau}^{\bullet}.$$

By our choice of z, this means that $\mathbb{A}_{3\widetilde{\rho}^n,4\widetilde{\rho}^n}(z)$ intersects $\mathcal{B}_{\tau}^{\bullet}$ and $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$ disconnects $B_{\epsilon\mathbb{F}}(x)$ from ∞ .

By (50), each point of $\mathbb{A}_{3\widetilde{\rho}^n,4\widetilde{\rho}^n}(z)\setminus \widetilde{U}^{\widetilde{\rho}^n}$ is contained in one of the squares $S\in \mathcal{S}^z_{\delta\widetilde{\rho}^n}(\mathbb{A}_{3\widetilde{\rho}^n,4\widetilde{\rho}^n}(z))$ such that $S\cap \mathcal{B}^{\bullet}_{\tau}\neq\varnothing$. By this and condition 2 in the definition of $E^{\widetilde{U}^{\widetilde{\rho}^n}}_{\widetilde{\rho}^n}(z)$, each point of $\mathbb{A}_{3\widetilde{\rho}^n,4\widetilde{\rho}^n}(z)\setminus \widetilde{U}^{\widetilde{\rho}^n}$ lies at D_h -distance at most $(c/100)\mathfrak{c}_{\widetilde{\rho}^n}e^{\xi h_{\widetilde{\rho}^n}(z)}$ from $\mathcal{B}^{\bullet}_{\tau}$. This together with the definition (53) of G^{ϵ}_{χ} shows that

(54)
$$\sup_{u \in \mathbb{A}_{3\widetilde{\rho}^n, 4\widetilde{\rho}^n}(z)} D_h(u, \mathcal{B}_{\tau}^{\bullet}) < c\mathfrak{c}_{\widetilde{\rho}^n} e^{\xi h_{\widetilde{\rho}^n}(z)}.$$

As a D_h -geodesic from a point outside of $\mathcal{B}_{\tau}^{\bullet}$ to 0 hits $\partial \mathcal{B}_{\tau}^{\bullet}$ exactly once, if such a geodesic hits $B_{\epsilon_{\mathbb{I}}}(x) \setminus \mathcal{B}_{\tau}^{\bullet}$, then it hits $B_{\epsilon_{\mathbb{I}}}(x)$ before entering $\mathcal{B}_{\tau}^{\bullet}$. Therefore, to prove property A, it suffices to consider a path P from a point outside of $\mathbb{C} \setminus B_{R_{\mathbb{I}}^{\epsilon}(\mathcal{B}_{\tau}^{\bullet})}(\mathcal{B}_{\tau}^{\bullet})$ to 0 which enters $B_{\epsilon_{\mathbb{I}}}(x)$ before entering $\mathcal{B}_{\tau}^{\bullet}$ and show that P cannot be a D_h -geodesic.

Since $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$ disconnects $B_{\epsilon \mathbb{I}}(x)$ from ∞ , the path P must cross from the outer boundary of $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$ to the inner boundary of $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$ before hitting $B_{\epsilon \mathbb{I}}(x)$ and, hence, also before hitting $\mathcal{B}^{\bullet}_{\tau}$. By condition 1 in the definition of $E_{\widetilde{\rho}^n}(z)$, each path between the inner and outer boundaries of $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$ has D_h -length at least $c\mathfrak{c}_{\widetilde{\rho}^n}e^{\xi h_{\widetilde{\rho}^n}(z)}$. Hence, the D_h -length of the segment of P after the first time it enters $\mathbb{A}_{2\widetilde{\rho}^n,3\widetilde{\rho}^n}(z)$ must be at least $c\mathfrak{c}_{\widetilde{\rho}^n}e^{\xi h_{\widetilde{\rho}^n}(z)}+\tau$.

But, P must enter $\mathbb{A}_{3\tilde{\rho}^n,4\tilde{\rho}^n}(z)$ before entering $\mathbb{A}_{2\tilde{\rho}^n,3\tilde{\rho}^n}(z)$, so by (54) P must hit a point at D_h -distance strictly smaller than $c\mathfrak{c}_{\tilde{\rho}^n}e^{\xi h_{\tilde{\rho}^n}(z)}$ from $\partial \mathcal{B}^{\bullet}_{\tau}$ before entering $\mathbb{A}_{2\tilde{\rho}^n,3\tilde{\rho}^n}(z)$. Such a point lies at D_h -distance strictly smaller than $c\mathfrak{c}_{\tilde{\rho}^n}e^{\xi h_{\tilde{\rho}^n}(z)} + \tau$ from 0. Combining this with the preceding paragraph shows that P cannot be a D_h -geodesic to 0.

Step 3: Proof that G_x^{ϵ} satisfies property B. For $n \in \mathbb{N}$, let

(55)
$$\widetilde{G}^n := \left\{ \max_{V \in \mathcal{V}(\widetilde{U}^{\widetilde{\rho}^n})} \sup_{u,v \in V} D_h(u,v; \mathbb{A}_{2\widetilde{\rho}^n,5\widetilde{\rho}^n}(z)) \le \frac{c}{2} \mathfrak{c}_{\widetilde{\rho}^n} e^{\xi h_{\widetilde{\rho}^n}(z)} \right\}$$

be the event appearing in the definition (53) of G_x^{ϵ} , so that $G_x^{\epsilon} = \bigcup_{n=1}^{\lfloor \eta \log \epsilon^{-1} \rfloor} \widetilde{G}^n$.

Let \mathfrak{p} be as in Lemma 3.3 with our above choice of c, δ , A. Since we chose these parameters in a manner depending only on the choice of metric, \mathfrak{p} depends only on the choice of metric. Just below, we will show using Lemma 3.3 that a.s.

(56)
$$\mathbb{P}\big[\widetilde{G}^n \big| \widetilde{U}^{\widetilde{\rho}^n}, h \big|_{\mathbb{C} \setminus \widetilde{U}^{\widetilde{\rho}^n}} \big] \ge \mathfrak{p} \quad \forall n \in \mathbb{N}.$$

Before proving (56), we explain why (56) implies property B. Recall that each $\widetilde{\rho}^n$ is a stopping time for the filtration generated by $h|_{B_{5r}(z)}$ for $r \geq 0$ and $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$ and each of the sets \widetilde{U}^r for r > 0 from (50) is determined by $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$. By the locality of D_h (Axiom II), for each $n \in \mathbb{N}$ the event \widetilde{G}^n of (55) is determined by $h|_{B_{6\widetilde{\rho}^n}(z)}$ and $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$.

Since $\widetilde{\rho}^n \geq 6\widetilde{\rho}^{n-1}$ and $\widetilde{U}^{\widetilde{\rho}^n}$ is disjoint from $\mathcal{B}^{\bullet}_{\tau}$, it follows that $h|_{B_{6\widetilde{\rho}^m}(z)}$ and $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$ and, hence, also \widetilde{G}^m for $m \leq n-1$ is determined by $(\widetilde{U}^{\widetilde{\rho}^n}, h|_{\mathbb{C}\setminus\widetilde{U}^{\widetilde{\rho}^n}})$. Hence, we can iterate (56) $\lfloor \eta \log \epsilon^{-1} \rfloor$ times to get that the conditional probability given $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$ that \widetilde{G}^n does not occur for every $n \in [1, \eta \log \epsilon^{-1}]_{\mathbb{Z}}$ is at most $(1 - \mathfrak{p})^{\lfloor \eta \log \epsilon^{-1} \rfloor}$. That is, a.s. $\mathbb{P}[G_x^{\epsilon}|\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}}] \geq 1 - C_0 \epsilon^{\alpha}$ for α slightly smaller than $\eta \log (1/(1 - \mathfrak{p}))$ and an appropriate choice of $C_0 > 1$ depending only on the choice of metric.

It remains to justify (56). To this end, let $\mathfrak{e} > 0$, let $\mathfrak{z} \in \frac{\mathfrak{er}}{4}\mathbb{Z}^2$, let $\mathfrak{r} \geq 6^n\mathfrak{er}$ be a dyadic multiple of \mathfrak{er} and let $\mathfrak{U} \in \mathcal{U}_{\mathfrak{r}}(\mathfrak{z})$. We will study the conditional law given $\{\epsilon = \mathfrak{e}, z = \mathfrak{z}, \widetilde{\rho}^n = \mathfrak{r}, \widetilde{U}^{\widetilde{\rho}^n} = \mathfrak{U}\}$. By Lemma 3.3,

$$\mathbb{P}\left[\max_{V\in\mathcal{V}(\mathfrak{U})}\sup_{u,v\in V}D_{h}(u,v;\mathbb{A}_{\mathfrak{r},5\mathfrak{r}}(\mathfrak{z}))\leq \frac{c}{2}\mathfrak{c}_{\mathfrak{r}}e^{\xi h_{\mathfrak{r}}(\mathfrak{z})}|h|_{\mathbb{C}\backslash\mathfrak{U}},E_{\mathfrak{r}}^{\mathfrak{U}}(\mathfrak{z})\right]\geq \mathfrak{p}.$$

We will now argue that

$$\{\epsilon = \mathfrak{e}, z = \mathfrak{z}, \widetilde{\rho}^n = \mathfrak{r}, \widetilde{U}^{\widetilde{\rho}^n} = \mathfrak{U}\} \in \sigma(h|_{\mathbb{C}\setminus\mathfrak{U}}, \mathbb{1}_{E^{\mathfrak{U}}_{\mathfrak{r}}(\mathfrak{z})}).$$

Recall that $\mathcal{B}^{\bullet}_{\tau}$ is a local set for h (Lemma 2.1). Since $\mathcal{B}^{\bullet}_{\tau}$ is a.s. determined by h, the event $\{\mathcal{B}^{\bullet}_{\tau} \cap \mathfrak{U} = \varnothing\}$ is a.s. determined by $h|_{\mathbb{C}\setminus\mathfrak{U}}$ and moreover $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$ is a.s. determined by $h|_{\mathbb{C}\setminus\mathfrak{U}}$ on $\{\mathcal{B}^{\bullet}_{\tau} \cap \mathfrak{U} = \varnothing\}$.

The points z and ϵ and the sets $\{\widetilde{U}^r\}_{r>0}$ from (50) are all determined by $(\mathcal{B}^{\bullet}_{\tau},h|_{\mathcal{B}^{\bullet}_{\tau}})$. Hence, each of these objects is determined by $h|_{\mathbb{C}\setminus\mathfrak{U}}$ on the event $\{\mathcal{B}^{\bullet}_{\tau}\cap\mathfrak{U}=\varnothing\}$. Each of the events $E^{\widetilde{U}^r}_r(z)$ is determined by $\mathcal{B}^{\bullet}_{\tau}$ and $h|_{\mathbb{A}_{2r,5r}(z)}$. Since $\widetilde{\rho}^n$ is the *smallest* radius $r\geq 6\rho^{n-1}$ which is a dyadic multiple of $\epsilon\mathbb{T}$ for which $E^{\widetilde{U}^r}_r(z)$ occurs, it follows that the event $\{\epsilon=\mathfrak{e},z=\mathfrak{z},\widetilde{\rho}^n=\mathfrak{r}\}$ is determined by $(\mathcal{B}^{\bullet}_{\tau},h|_{\mathcal{B}^{\bullet}_{\tau}})$, $E^{\mathfrak{U}}_{\tau}(z)$ and $h|_{\mathbb{C}\setminus\mathbb{A}_{2\mathfrak{r},5\mathfrak{r}}(z)}$. We have $\mathfrak{U}\subset\mathbb{A}_{2\mathfrak{r},5\mathfrak{r}}(z)$, so $\{\epsilon=\mathfrak{e},z=\mathfrak{z},\widetilde{\rho}^n=\mathfrak{r}\}$ is a.s. determined by $\mathbb{1}_{E_{\tau}(z)}$ and $h|_{\mathbb{C}\setminus\mathbb{U}}$ on the event $\{\mathcal{B}^{\bullet}_{\tau}\cap\mathfrak{U}=\varnothing\}$.

By (50), we have $\mathcal{B}^{\bullet}_{\tau} \cap \mathfrak{U} = \emptyset$ on the event $\{\widetilde{U}^{\widetilde{\rho}^n} = \mathfrak{U}\}$. Combining this with the preceding two paragraphs gives (58). Combining (57) and (58) and using that $\{\mathfrak{z} = z, \widetilde{\rho}^n = \mathfrak{r}, \widetilde{U}^{\widetilde{\rho}^n} = \mathfrak{U}\} \subset E^{\mathfrak{U}}_{\mathfrak{r}}(\mathfrak{z})$ by definition shows that (56) holds. \square

We will most often use the following variant of Lemma 3.6 where we prevent D_h -geodesics from hitting a boundary arc rather than a neighborhood of a point.

LEMMA 3.7. Let α be as in Lemma 3.6. Let $\mathbb{r} > 0$; let τ be a stopping time for the filtration generated by $\{(\mathcal{B}_s^{\bullet}, h|_{\mathcal{B}_s^{\bullet}})\}_{s\geq 0}$. Also, let $\epsilon \in (0, 1)$ and $I \subset \partial \mathcal{B}_{\tau}^{\bullet}$ be an arc, each

chosen in a manner depending only on $(\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}})$, such that I can be disconnected from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$ by a set of Euclidean diameter at most $\epsilon_{\mathbb{F}}$. There is an event $G_I \in \sigma(\mathcal{B}_{\sigma_{\tau,\mathbb{F}}^{\bullet}}^{\bullet}, h|_{\mathcal{B}_{\sigma_{\tau,\mathbb{F}}^{\bullet}}^{\bullet}})$ with the following properties:

- A. If $R_{\mathbb{T}}^{\epsilon}(\mathcal{B}_{\tau}^{\bullet}) \leq \operatorname{diam} \mathcal{B}_{\tau}^{\bullet}$ and G_{I} occurs, then no D_{h} -geodesic from 0 to a point in $\mathbb{C} \setminus \mathcal{B}_{\sigma_{\tau}^{\bullet}}^{\bullet}$. can pass through I.
- B. There is a deterministic constant $C_0 > 1$ depending only on the choice of metric such that a.s. $\mathbb{P}[G_I|\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}}] \geq 1 - C_0 \epsilon^{\alpha}$.

Since I can be disconnected from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$ by a set of Euclidean diameter at most $\epsilon \mathbb{r}$, we can choose a point $x \in \partial \mathcal{B}_{\tau}^{\bullet}$ in a manner depending only on $(\mathcal{B}_{\tau}^{\bullet}, h|_{\mathcal{B}_{\tau}^{\bullet}})$ such that $B_{\epsilon \mathbb{I}}(x)$ disconnects I from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$. Let $G_I := G_x^{\epsilon}$ be the event of Lemma 3.6 for this choice of x. Then, G_I satisfies condition B in the lemma statement. Moreover, $B_{\epsilon \mathbb{F}}(x) \subset$ $B_{R_{\bullet}^{c}(\mathcal{B}_{\bullet}^{\bullet})}(\mathcal{B}_{\bullet}^{\bullet})$ by (48), so each path from a point in the unbounded connected component of $\mathbb{C}\setminus B_{R^{\epsilon}_{\mathbb{T}}(\mathcal{B}^{\bullet}_{\tau})}(\mathcal{B}^{\bullet}_{\tau})$ which first hits $\partial \mathcal{B}^{\bullet}_{\tau}$ at a point of I must pass through $B_{\epsilon_{\mathbb{T}}}(x)$. By (48), $\mathbb{C}\setminus\mathcal{B}^{\bullet}_{\sigma^{\epsilon}_{\tau}}$ is contained in the unbounded connected component of $\mathbb{C}\setminus B_{R^{\epsilon}_{\tau}(\mathcal{B}^{\bullet}_{\tau})}(\mathcal{B}^{\bullet}_{\tau})$. By this and the corresponding condition from Lemma 3.6, we get that G_I satisfies condition A in the lemma statement.

3.4. Proof of Theorem 3.1. Continue to fix parameters c, δ , A, η for which the conclusion of Lemma 3.5 holds. For the rest of the paper we will no longer need to recall the precise definitions of the events $E_r(z)$ and $E_r^U(z)$. Rather, we only need the conclusions of Lemmas 3.5 and 3.7.

We will actually prove a much more quantitative version of Theorem 3.1 (see Theorem 3.9 below) which gives a quantitative bound on how large N needs to be in terms of t and p provided we truncate on a global regularity event, which we now define.

It is shown in [5], Theorem 1.7, that if $\xi = \gamma/d_{\gamma}$ and $Q = 2/\gamma + \gamma/2$, then D_h is a.s. locally χ -Hölder continuous w.r.t. the Euclidean metric for any $\chi \in (0, \xi(Q-2))$. Henceforth, fix such a χ , chosen in a manner depending only on ξ and Q. We also recall the stopping time $\tau_{\mathbb{I}}$ from (33). For $a \in (0, 1)$, we define $\mathcal{E}_{\mathbb{I}}(a)$ to be the event that the following is true:

- 1. $B_{a\mathbb{r}}(0) \subset \mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet}$.
- 2. $\tau_{3\mathbb{r}} \tau_{2\mathbb{r}} \ge a \mathfrak{c}_{\mathbb{r}} e^{-\xi h_{\mathbb{r}}(0)}$. 3. $\mathfrak{c}_{\mathbb{r}}^{-1} e^{-\xi h_{\mathbb{r}}(0)} D_h(u, v) \le (\frac{|u-v|}{\mathbb{r}})^{\chi}$ for each $u, v \in B_{4\mathbb{r}}(0)$ with $|u-v|/\mathbb{r} \le a$.
- 4. In the notation (45), we have $\rho_{\epsilon \mathbb{F}}^{\lfloor \eta \log \epsilon^{-1} \rfloor}(z) \leq \epsilon^{1/2} \mathbb{F}$ for each $z \in (\frac{\epsilon \mathbb{F}}{4} \mathbb{Z}^2) \cap B_{4\mathbb{F}}(0)$ and each dyadic $\epsilon \in (0, a]$.

LEMMA 3.8. For each $p \in (0,1)$, there exists a = a(p) > 0 such that $\mathbb{P}[\mathcal{E}_{\mathbb{F}}(a)] \geq p$ for every $\mathbb{r} > 0$.

PROOF. By Axiom V, if a is chosen sufficiently small, then the probability of each of conditions 1 and 2 is at least 1 - (1 - p)/4. By [5], Proposition 3.18, if a is chosen sufficiently small, then the probability of condition 3 is at least 1 - (1 - p)/4. By Lemma 3.5 and a union bound over dyadic values of ϵ with $\epsilon \in (0, a]$, the probability of condition 4 is at least 1 - (1 - p)/4.

Theorem 3.1 will be an immediate consequence of Lemma 3.8 together with the following quantitative estimate. For the statement, we recall the times $\tau_{\mathbb{I}}$ from (33).

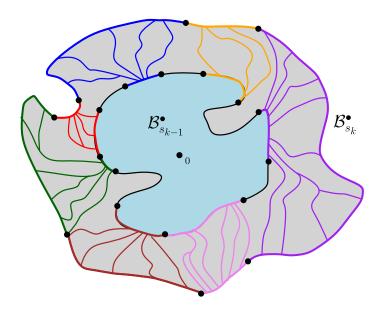


FIG. 8. Illustration of one stage of the iterative procedure used in the proof of Theorem 3.1. The endpoints of the arcs in \mathcal{I}_{k-1} are shown with black dots. Each of the intervals $I \in \mathcal{I}_{k-1}$ with $I' \neq \emptyset$ (i.e., those which are hit by leftmost geodesics from 0 to $\partial \mathcal{B}_{s_k}^{\bullet}$) is assigned a unique color. The corresponding arc $I' \subset \mathcal{B}_{s_k}^{\bullet}$ and some representative leftmost geodesics from I' to I are shown in the same color. Black arcs $I \in \mathcal{I}_{k-1}$ are ones with $I' = \emptyset$. Note that here we have shown some geodesics merging into each other, but we have not yet established whether or not this happens. We show in Lemma 3.10 that typically the number of arcs in \mathcal{I}_k decays geometrically in k.

THEOREM 3.9. For each $a \in (0, 1)$, there is a constant $b_0 > 0$ depending only on a and constants $b_1, \beta > 0$ depending only on the choice of metric D such that the following is true. For each $\mathbb{F} > 0$, each $N \in \mathbb{N}$ and each stopping time τ for $\{(\mathcal{B}_s^{\bullet}, h|_{\mathcal{B}_s^{\bullet}})\}_{s \geq 0}$ with $\tau \in [\tau_{\mathbb{F}}, \tau_{2\mathbb{F}}]$ a.s., the probability that $\mathcal{E}_{\mathbb{F}}(a)$ occurs and there are more than N points of $\partial \mathcal{B}_{\tau}^{\bullet}$ which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau+N-\beta}^{\bullet} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)}$ is at most $b_0 e^{-b_1 N^{\beta}}$.

In the rest of this section we prove Theorem 3.9. Fix $\mathbb{r} > 0$ and a stopping time τ as in Theorem 3.9. Let \mathcal{I}_0 be a collection of disjoint boundary arcs of $\mathcal{B}^{\bullet}_{\tau}$, chosen in a manner depending only on $(\mathcal{B}^{\bullet}_{\tau}, h|_{\mathcal{B}^{\bullet}_{\tau}})$. We allow arcs to be open, half-open or closed. In particular, the union of all of the arcs of \mathcal{I}_0 is allowed to be all of $\partial \mathcal{B}^{\bullet}_{\tau}$ (this is in fact the typical case we will be interested in). The idea of the proof is to show that, for *any* choice of \mathcal{I}_0 , the probability that $\mathcal{E}_{\mathbb{F}}(a)$ occurs, and there are more than N arcs in \mathcal{I}_0 which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{F}}+N^{-\beta}\mathfrak{c}_{\mathbb{F}}e^{\xi h_{\mathbb{F}}(0)}}$ is at most $b_0e^{-b_1N^{\beta}}$. Taking \mathcal{I}_0 to be a huge collection of tiny arcs will then prove Theorem 3.9; see Figure 8 for an illustration of the setup.

3.4.1. Inductive definition of radii and boundary arcs. We start by inductively defining for each $k \in \mathbb{N}$ a radius $s_k \geq s_{k-1}$ and a finite collection \mathcal{I}_k of disjoint boundary arcs of $\mathcal{B}_{s_k}^{\bullet}$, chosen in a manner depending only on $(\mathcal{B}_{s_k}^{\bullet}, h|_{\mathcal{B}_{s_k}^{\bullet}})$ and satisfying $\#\mathcal{I}_k \leq \#\mathcal{I}_{k-1}$.

Set $s_0 = \tau$, and let \mathcal{I}_0 be as above. Inductively, suppose s_{k-1} and \mathcal{I}_{k-1} have been defined. Let ϵ_{k-1} be the smallest dyadic number with $\epsilon_{k-1} \ge (\#\mathcal{I}_{k-1})^{-1/4}$, and in the notation (49) define

$$(59) s_k := \sigma_{s_{k-1}, \mathbb{F}}^{\epsilon_{k-1}}.$$

For $I \in \mathcal{I}_{k-1}$, let I' be the set of points $x \in \partial \mathcal{B}_{s_k}^{\bullet}$ for which the leftmost D_h -geodesic from 0 to x passes through I. Note that we may have $I' = \emptyset$. Define

(60)
$$\mathcal{I}_k := \{ I' : I \in \mathcal{I}_{k-1}, I' \neq \emptyset \}.$$

We make the following observations about \mathcal{I}_k :

- By Lemma 2.7, \mathcal{I}_k is a collection of disjoint arcs of $\partial \mathcal{B}_{s_k}^{\bullet}$. The arcs in \mathcal{I}_k can be closed, open or half-open.
- \mathcal{I}_k is determined by $(\mathcal{B}_{s_k}^{\bullet}, h|_{\mathcal{B}_{s_k}^{\bullet}})$. Indeed, this is because of Axiom II and the fact that a D_h -geodesic from 0 to $\partial \mathcal{B}_{s_k}^{\bullet}$ cannot exit $\mathcal{B}_{s_k}^{\bullet}$ so must also be a geodesic for the internal metric of D_h on $\mathcal{B}_{s_k}^{\bullet}$.
- If we let $I^{(k)}$ for $I \in \mathcal{I}_0$ be the set of $x \in \partial \mathcal{B}_{s_k}^{\bullet}$ for which the leftmost D_h -geodesic from 0 to x pass through I, then we can equivalently define $\mathcal{I}_k = \{I^{(k)} : I \in \mathcal{I}_0, I^{(k)} \neq \emptyset\}$. Indeed, this is because $I^{(k)}$ is obtained from I by applying the operation $I \mapsto I' k$ times.
- 3.4.2. The cardinalities of the \mathcal{I}_k 's decrease geometrically. To lighten notation, in what follows we abbreviate

(61)
$$n_k := \# \mathcal{I}_k \text{ so that } \epsilon_k \in [n_k^{-1/4}, 2n_k^{-1/4}].$$

For $N \in \mathbb{N}$, we also define

(62)
$$K_N := \min\{k \in \mathbb{N}_0 : s_k > \tau_{3\mathbb{F}} \text{ or } n_k < N\}.$$

We observe that K_N is a stopping time for the filtration generated by $\{(\mathcal{B}_{s_k}^{\bullet}, h|_{\mathcal{B}_{s_k}^{\bullet}})\}_{k \in \mathbb{N}_0}$ and that K_N is a nonincreasing function of N. We will also have occasion to consider a parameter $N_0 \in \mathbb{N}$, which we will eventually choose to be sufficiently large in a manner depending only on a (at several places we will assume that N_0 is sufficiently large that some estimate is true). Most of our estimates will require that $N \geq N_0$.

The key ingredient in the proof of Theorem 3.9 is the following lemma which tells us that the arc counts n_k typically decrease almost geometrically in k:

LEMMA 3.10. Let $k_0 = 0$ and for $j \in \mathbb{N}$; inductively, let k_j be the smallest $k \ge k_{j-1}$ for which $n_k \le \frac{1}{2}n_{k_{j-1}}$. There are constants C > 1 and $\widetilde{\beta} > 0$, depending only on the choice of metric, such that if N_0 is chosen to be sufficiently large, in a manner depending only on a, then for $j \in \mathbb{N}_0$ and M > 1,

(63)
$$\mathbb{P}[k_{j+1} - k_j > M, k_{j+1} \le K_{N_0}, \mathcal{E}_{\mathbb{F}}(a) | \mathcal{B}_{s_{k_j}}^{\bullet}, h|_{\mathcal{B}_{s_{k_j}}^{\bullet}}] \le (n_{k_j}/C)^{-\widetilde{\beta}M}.$$

To prove Lemma 3.10, we will show that, roughly speaking, $n_k \leq \frac{1}{2}n_{k-1}$ with high conditional probability given $(\mathcal{B}_{s_k}^{\bullet}, h|_{\mathcal{B}_{s_k}^{\bullet}})$, then multiply the resulting estimate over M values of k to get (63). Recall that ϵ_k is a negative power of 2 chosen so that $\epsilon_k \in [n_k^{-1/4}, 2n_k^{-1/4})$. The basic idea is that Lemma 2.15 tells us that most of the arcs of $\partial \mathcal{B}_{s_k}^{\bullet}$ can be disconnected from ∞ by sets of diameter smaller than ϵ_k , and Lemma 3.7 tells us that each of these arcs is unlikely to survive to the next step (i.e., one has $I' = \emptyset$ with high probability).

Let us first record what we get from Lemmas 2.15 and 3.7. For $k \in \mathbb{N}_0$, let \mathcal{I}_k^* be the set of "bad" arcs $I \in \mathcal{I}_k$ which *cannot* be disconnected from ∞ in $\mathbb{C} \setminus \mathcal{B}_{s_k}^{\bullet}$ by a set of Euclidean diameter at most $\epsilon_k \mathbb{F}$. By Lemma 2.15 applied with $n = n_k \times \epsilon_k^{-4}$ and $C \times \mathbb{F}/\epsilon_k$ and condition 1 in the definition of $\mathcal{E}_{\mathbb{F}}(a)$, there is a universal constant A > 0 such that on $\mathcal{E}_{\mathbb{F}}(a)$,

(64)
$$#\mathcal{I}_k^* \le A \frac{2^4}{a^2} n_k^{1/2} \quad \forall k \in [0, K_1 - 1]_{\mathbb{Z}}.$$

We henceforth assume that N_0 is chosen sufficiently large that $A\frac{2^4}{a^2}N_0^{-1/2} \le 1/4$, so that the right side of (64) is smaller than $n_k/4$ whenever $n_k \ge N_0$ (which in particular is the case if $k \le K_{N_0} - 1$). Then, on $\mathcal{E}_{\mathbb{T}}(a)$,

(65)
$$\#\mathcal{I}_k^* \le \frac{1}{4} n_k \quad \forall k \in [0, K_{N_0} - 1]_{\mathbb{Z}}.$$

We will now explain how Lemma 3.7 allows us to "kill off" most of the arcs not in \mathcal{I}_k^* . For $I \in \mathcal{I}_k \setminus \mathcal{I}_k^*$, let G_I be the event of Lemma 3.7 with $\tau = s_k$ and $\epsilon = \epsilon_k$, and define a second set of bad arcs

(66)
$$\mathcal{I}_k^{**} := \{ I \in \mathcal{I}_k \setminus \mathcal{I}_k^* : G_I \text{ does not occur} \}.$$

By assertion A of Lemma 3.7, if (in the notation (48)) we have $R_{\mathbb{F}}^{\epsilon_k}(\mathcal{B}_{s_k}^{\bullet}) \leq \operatorname{diam}(\mathcal{B}_{s_k}^{\bullet})$ and G_I occurs, then the interval I' from the definition (60) of \mathcal{I}_{k+1} is empty.

We now want to say that this condition is satisfied for all $k \le K_{N_0} - 1$. Assume that N_0 is chosen sufficiently large that $N_0^{-1/4} \le a/2$. Then, $\epsilon_k \le a$ for each $k \in [0, K_{N_0} - 1]_{\mathbb{Z}}$ so, if $\mathcal{E}_{\mathbb{F}}(a)$ occurs, then condition 4 in the definition of $\mathcal{E}_{\mathbb{F}}(a)$ shows that

By condition 1 in the definition of $\mathcal{E}_{\mathbb{F}}(a)$ and the definition (62) of K_{N_0} , one has $B_{a\mathbb{F}}(0) \subset$ $\mathcal{B}_{s_k}^{\bullet} \subset B_{3\mathbb{F}}(0)$ for each $k \in [0, K_{N_0} - 1]_{\mathbb{Z}}$. In particular, diam $(\mathcal{B}_{s_k}^{\bullet}) \geq a\mathbb{F}$. Therefore, (67) together with the definition (48) of $R_{\mathbb{F}}^{\epsilon_k}(\mathcal{B}_{s_k}^{\bullet})$ shows that if $\mathcal{E}_{\mathbb{F}}(a)$ occurs, then

$$(68) R_{\mathbb{F}}^{\epsilon_k}(\mathcal{B}_{s_k}^{\bullet}) \leq \left(6\epsilon_k^{1/2} + \epsilon_k\right)\mathbb{F} \leq 7n_k^{-1/8}\mathbb{F} \leq \frac{7}{aN_0^{1/8}}\operatorname{diam}(\mathcal{B}_{s_k}^{\bullet}) \forall k \in [0, K_{N_0} - 1]_{\mathbb{Z}}.$$

Hence, if we choose N_0 sufficiently large that $aN_0^{1/8} \ge 7$, then on $\mathcal{E}_{\mathbb{F}}(a)$, we have $R_{\mathbb{F}}^{\epsilon_k}(\mathcal{B}_{s_k}^{\bullet}) \le \operatorname{diam}(\mathcal{B}_{s_k}^{\bullet})$ for each $k \in [0, K_{N_0} - 1]_{\mathbb{Z}}$. By combining this with the preceding paragraph, if $\mathcal{E}_{\mathbb{F}}(a)$ occurs and $k \in [0, K_{N_0} - 1]_{\mathbb{Z}}$, then $I' = \emptyset$ for every $I \in \mathcal{I}_k \setminus (\mathcal{I}_k^* \cup \mathcal{I}_k^{**})$. Therefore,

(69)
$$n_{k+1} = \#\mathcal{I}_{k+1} \le \#\mathcal{I}_k^* + \#\mathcal{I}_k^{**} \quad \forall k \in [0, K_{N_0} - 1]_{\mathbb{Z}}.$$

We will also need assertion B of Lemma 3.7, which tells us that there is an exponent $\alpha > 0$ and a constant $C_0 > 0$ depending only on the choice of metric such that for $k \in \mathbb{N}_0$,

(70)
$$\mathbb{P}\left[G_{I}\big|\mathcal{B}_{s_{k}}^{\bullet},h\big|_{\mathcal{B}_{s_{k}}^{\bullet}}\right] \geq 1 - C_{0}\epsilon_{k}^{\alpha} \quad \forall I \in \mathcal{I}_{k} \setminus \mathcal{I}_{k}^{*}.$$

PROOF OF LEMMA 3.10. By (66) and (70), $\mathbb{E}[\#\mathcal{I}_k^{**}|\mathcal{B}_{s_k}^{\bullet}, h|_{\mathcal{B}_{s_k}^{\bullet}}] \leq C_0 \epsilon_k^{\alpha} \#(\mathcal{I}_k \setminus \mathcal{I}_k^*)$. By Markov's inequality, for $k \in \mathbb{N}_0$,

(71)
$$\mathbb{P}\left[\#\mathcal{I}_{k}^{**} > \frac{1}{4}\#(\mathcal{I}_{k} \setminus \mathcal{I}_{k}^{*})|\mathcal{B}_{s_{k}}^{\bullet}, h|_{\mathcal{B}_{s_{k}}^{\bullet}}\right] \leq 4C_{0}\epsilon_{k}^{\alpha}.$$

For $k \in [k_j, k_{j+1} - 1]_{\mathbb{Z}}$, one has $n_k \ge \frac{1}{2} n_{k_j}$, and, hence, $\epsilon_k \le (n_{k_j}/4)^{-1/4}$. By Lemma 3.7, $G_I \in \sigma(\mathcal{B}_{s_{k+1}}^{\bullet}, h|_{\mathcal{B}_{s_{k+1}}^{\bullet}})$ for each $I \in \mathcal{I}_k$. Since k_j is a stopping time for the filtration generated by $(\mathcal{B}_{s_k}^{\bullet}, h|_{\mathcal{B}_{s_k}^{\bullet}})$ and $\mathcal{I}_k, \mathcal{I}_k^*, \mathcal{I}_k^{**} \in \sigma(\mathcal{B}_{s_{k+1}}^{\bullet}, h|_{\mathcal{B}_{s_{k+1}}^{\bullet}})$, we can set $k = k_j + m$ in (71) and iterate M times to get that

(72)
$$\mathbb{P}\left[k_{j+1} - k_{j} > M \text{ and } \#\mathcal{I}_{k_{j}+m}^{**} > \frac{1}{4} \#\left(\mathcal{I}_{k_{j}+m} \setminus \mathcal{I}_{k_{j}+m}^{*}\right)\right]$$

$$\forall m \in [0, M]_{\mathbb{Z}} |\mathcal{B}_{s_{k_{j}}}^{\bullet}, h|_{\mathcal{B}_{s_{k_{j}}}^{\bullet}}\right] \leq (n_{k_{j}}/C)^{-\alpha M/4}$$

for an appropriate constant C > 1 depending only on the choice of metric.

By (65) and (69), if $\#\mathcal{I}_k^{**} \leq \frac{1}{4}\#(\mathcal{I}_k \setminus \mathcal{I}_k^*)$, $k \leq K_{N_0} - 1$, and $\mathcal{E}_{\mathbb{F}}(a)$ occurs, then

(73)
$$n_{k+1} \le \#\mathcal{I}_k^* + \#\mathcal{I}_k^{**} \le \frac{1}{2}n_k.$$

By the definition of k_{j+1} , if $m \in \mathbb{N}$, $\#\mathcal{I}_{k_j+m}^{**} \leq \frac{1}{4} \#(\mathcal{I}_{k_j+m} \setminus \mathcal{I}_{k_j+m}^*)$, $k_j + m \leq K_{N_0} - 1$, and $\mathcal{E}_{\mathbb{F}}(a)$ occurs, then $k_{j+1} \leq k_j + m$. In other words, the event inside the conditional probability in (72) contains the event $\{k_{j+1} - k_j > M\} \cap \{k_j \leq K_{N_0} - 1\} \cap \mathcal{E}_{\mathbb{F}}(a)$. We therefore get (63) with $\beta = \alpha/4$. \square

3.4.3. *Conclusion of the proof.* We will now apply Lemma 3.10 iteratively to bound how much we need to increase the radius of our metric balls to get down to *N* remaining boundary arcs.

LEMMA 3.11. For each $a \in (0, 1)$, there are constants $b_0, b_1, \beta > 0$ as in the statement of Theorem 3.9 such that for each $\mathbb{F} > 0$ and $N \in \mathbb{N}$,

(74)
$$\mathbb{P}\left[s_{K_N} > \tau + N^{-\beta} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)}, \mathcal{E}_{\mathbb{F}}(a)\right] \leq b_0 e^{-b_1 N^{\beta}}.$$

PROOF. Throughout the proof we write *c* for a constant which is only allowed to depend on the metric and which may change from line to line.

Step 1: Bounding s_{K_N} in terms of the n_k 's. By the first inequality in (68), if $\mathcal{E}_{\mathbb{F}}(a)$ occurs and $k \leq K_{N_0} - 1$, then $R_{\mathbb{F}}^{\epsilon_k}(\mathcal{B}_{s_k}^{\bullet}) \leq (\epsilon_k + 6\epsilon_k^{1/2})_{\mathbb{F}}$. That is, each point of each of the balls $B_{\lfloor \eta \log \epsilon_k^{-1} \rfloor}(z)$ for $z \in (\frac{\epsilon_k r}{4}\mathbb{Z}^2) \cap B_{\epsilon_k r}(\mathcal{B}_{s_k}^{\bullet})$ lies at Euclidean distance at most $(\epsilon_k + 6\epsilon_k^{1/2})_{\mathbb{F}}$ from $\mathcal{B}_{s_k}^{\bullet}$. Since $k \leq K_{N_0} - 1$, we have $\mathcal{B}_{s_k}^{\bullet} \subset B_{3\mathbb{F}}(0)$. By the Hölder continuity condition 3 in the definition of $\mathcal{E}_{\mathbb{F}}(a)$, if we choose N_0 large enough that $(N_0^{-1/4} + 6N_0^{-1/8}) \leq a$, then each such ball is contained in $\mathcal{B}_{s_k+\delta_k}^{\bullet}$ for $\delta_k := (\epsilon_k + 6\epsilon_k^{1/2})^{\chi} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)}$. By the definition (49) of $\sigma_{s_k,\mathbb{F}}^{\epsilon_k}$ and (59), if $\mathcal{E}_{\mathbb{F}}(a)$ occurs (and N_0 is chosen large enough), then

(75)
$$s_{k+1} = \sigma_{s_k, \mathbb{F}}^{\epsilon_k} \le s_k + (\epsilon_k + 6\epsilon_k^{1/2})^{\chi} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)} \le s_k + c n_k^{-\chi/8} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)}$$

$$\forall k \in [0, K_{N_0} - 1]_{\mathbb{Z}},$$

where in the last inequality we recall that $\epsilon_k \approx n_k^{1/4}$. Summing this estimate and recalling that $s_0 = \tau$ shows that on $\mathcal{E}_{\mathbb{T}}(a)$,

(76)
$$s_{K_N} \le \tau + c \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)} \sum_{k=1}^{K_N - 1} n_k^{-\chi/8} \quad \forall N \ge N_0.$$

Step 2: Bounding s_{K_N} in terms of $k_j - k_{j-1}$. Fix $N \ge N_0$. Let k_j for $j \in \mathbb{N}_0$ be as in Lemma 3.10, and let $J = J_N$ be the smallest $j \in \mathbb{N}$ for which $k_j \ge K_N$. By definition, $n_{k_{j-1}} \ge 2n_{k_j}$ for each $j \in \mathbb{N}_0$. By iterating this, we get that $n_{k_j} \ge 2^{J-j-1}n_{k_{J-1}} \ge 2^{J-j-1}N$ for each $j \in [0, J-1]_{\mathbb{Z}}$. Since $n_k \ge n_{k_j}/2$ for $k \in [k_j, k_{j+1}-1]_{\mathbb{Z}}$, it follows that

(77)
$$n_{k_j} \ge 2^{J-j-2} N \quad \forall j \in [0, J-1]_{\mathbb{Z}} \ \forall k \in [k_j, k_{j+1}-1]_{\mathbb{Z}}.$$

Since $k_J - 1 \ge K_N - 1$, we can plug this into (76) to get

(78)
$$s_{K_N} \le \tau + cN^{-\chi/8} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)} \sum_{j=0}^{J-1} (k_{j+1} - k_j) 2^{-(J-j)}.$$

Step 3: Bounding $k_j - k_{j-1}$. By taking unconditional expectations of both sides of (63) from Lemma 3.10 and then applying (77), we get that for M > 0 and $j \in \mathbb{N}_0$,

$$\begin{split} \mathbb{P}\big[k_{j+1} - k_j > M, \, j \leq J - 1, \, \mathcal{E}_{\mathbb{F}}(a)\big] &\leq \mathbb{E}\big[(n_{k_j}/C)^{-\widetilde{\beta}M}\big] \\ &\leq (4C)^{\widetilde{\beta}M} N^{-\widetilde{\beta}M} \mathbb{E}\big[2^{-(J-j)\widetilde{\beta}M}\big]. \end{split}$$

Now, set $M = N^{\chi/16}$ and sum over all $j \in \mathbb{N}_0$ to get

(79)
$$\mathbb{P}\left[\exists j \in [0, J-1]_{\mathbb{Z}} \text{s.t.} k_{j+1} - k_{j} > N^{\chi/16}, \mathcal{E}_{\mathbb{F}}(a)\right]$$

$$\leq (4C)^{\widetilde{\beta}N^{\chi/16}} N^{-\widetilde{\beta}N^{\chi/16}} \mathbb{E}\left[\sum_{j=0}^{J-1} 2^{-(J-j)\widetilde{\beta}N^{\chi/16}}\right]$$

$$\leq (4C)^{\widetilde{\beta}N^{\chi/16}} N^{-\widetilde{\beta}N^{\chi/16}},$$

with the implicit constant depending only on the choice of metric. This last quantity is bounded above by $b_0 e^{-b_1 N^{\beta}}$ for constants $b_0, b_1, \beta > 0$ satisfying the conditions in Theorem 3.9. We can arrange that $\beta < \chi/16$. Hence, if $\mathcal{E}_{\mathbb{F}}(a)$ occurs, then except on an event of probability at most $b_0 e^{-b_1 N^{\beta}}$,

(80)
$$k_{j+1} - k_j \le N^{\chi/16} \quad \forall j \in [0, J-1]_{\mathbb{Z}}.$$

Plugging this bound into (78) shows that if $N \ge N_0$ and $\mathcal{E}_{\mathbb{T}}(a)$ occurs, then except on an event of probability at most $b_0 e^{-b_1 N^{\beta}}$,

(81)
$$s_{K_N} \le \tau + cN^{-\chi/16} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)} \sum_{j=0}^{J-1} 2^{-(J-j)} \le \tau + cN^{-\beta} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)}.$$

After possibly shrinking β (in order to absorb the c in (81) into a small power of N^{β}) and increasing b_0 (to make it so that $b_0e^{-b_1N^{\beta}} > 1$ for $N < N_0$), we obtain the statement of the lemma. \square

PROOF OF THEOREM 3.9. Let b_0 , b_1 , $\beta > 0$ be as in Lemma 3.11. By condition 2 in the definition of $\mathcal{E}_{\mathbb{T}}(a)$ and since $\tau \leq \tau_{2\mathbb{T}}$, if N is sufficiently large that $N^{-\beta} < a$, then $s_{K_N} \leq \tau + N^{-\beta}\mathfrak{c}_{\mathbb{T}}e^{\xi h_{\mathbb{T}}(0)}$ implies that $s_{K_N} < \tau_{3\mathbb{T}}$. By the definition (62) of K_N , this implies that $\#\mathcal{I}_{K_N} \leq N$. Hence, Lemma 3.11 implies that if $\mathcal{E}_{\mathbb{T}}(a)$ occurs, then, except on an event of probability at most $b_0 e^{-b_1 N^{\beta}}$, the number of $I \in \mathcal{I}_0$ for which there is a leftmost D_h -geodesic from 0 to $\partial \mathcal{B}_{\tau+N^{-\beta}\mathfrak{c}_{\mathbb{T}}e^{\xi h_{\mathbb{T}}(0)}}^{\bullet}$ which passes through I is at most N. This holds regardless of the initial choice of \mathcal{I}_0 .

Let X_N be the set of points on $\partial \mathcal{B}_{\tau}^{\bullet}$ which are hit by a leftmost D_h -geodesic from 0 to $\partial \mathcal{B}_{\tau+N^{-\beta}c_{\mathbb{T}}e^{\xi h_{\mathbb{T}}(0)}}^{\bullet}$. For a given $n_0 \in \mathbb{N}$, we choose \mathcal{I}_0 to be a collection of n_0 disjoint half-open arcs of $\partial \mathcal{B}_{\tau}^{\bullet}$ which cover $\partial \mathcal{B}_{\tau}^{\bullet}$ and such that the harmonic measure from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau}^{\bullet}$ of each $I \in \mathcal{I}_0$ is $1/n_0$. Applying the preceding paragraph to this choice of \mathcal{I}_0 , we find that if $\mathcal{E}_{\mathbb{T}}(a)$ occurs, then, except on an event of probability at most $b_0 e^{-b_1 N^{\beta}}$, the set X_N can be covered by at most N boundary arcs of $\mathcal{B}_{\tau}^{\bullet}$ which each have harmonic measure from ∞ at most $1/n_0$. Since n_0 can be made arbitrarily large, this implies that on $\mathcal{E}_{\mathbb{T}}(a)$ it holds, except on en event of probability at most $b_0 e^{-b_1 N^{\beta}}$, that $\#X_N \leq N$. \square

4. Reducing to a single geodesic. In this section we will prove Theorem 1.3. Unlike in the case of Theorem 3.1, we do not prove a version which is uniform over the Euclidean scale since we do not need such a statement in [8]. The only result from Section 3 which we need in this section is Theorem 3.1.

To deduce Theorem 1.3 from Theorem 3.1 we need to reduce from finitely many points hit by geodesics to one point. To accomplish this, in Section 4.1 we show that if $\tau_{\mathbb{I}}$ is as in (33), A > 1 is fixed and $I \subset \partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}$ is an arc chosen in a manner depending only on $(\mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}, h|_{\mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}})$, then with positive conditional probability given $(\mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}, h|_{\mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}})$, every D_h -geodesic from 0 to a point outside of $B_{A\mathbb{I}}(0)$ passes through I (see Lemma 4.1 for a precise statement). This is

proven using the Markov property of the GFF and the FKG inequality via a similar argument to the one in Section 3.3. The argument is significantly simpler, however, since the estimate we need is much less quantitative than the one in Lemma 3.7. See Figure 9 for an illustration of the argument.

In Section 4.2 we combine Theorem 3.1 with the result of Section 4.1, applied with I equal to the set of points in $x \in \partial \mathcal{B}_{\tau_{2r}}^{\bullet}$ such that the leftmost D_h -geodesic from 0 to x hits some specified point of $\partial \mathcal{B}_{\tau_r}^{\bullet}$, to show that the following is true. There is an M>1 such that with positive conditional probability given $(\mathcal{B}_{\tau_r}^{\bullet}, h|_{\mathcal{B}_{\tau_r}^{\bullet}})$, all of the leftmost D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau_{Mr}}^{\bullet}$ hit the same point of $\partial \mathcal{B}_{\tau_r}^{\bullet}$. Using the uniqueness of geodesics to rational points, we can improve this to say that, with positive probability, all of the geodesics from 0 to $\partial \mathcal{B}_{\tau_{Mr}}^{\bullet}$ pass through the same point of $\partial \mathcal{B}_{\tau_r}^{\bullet}$. Using a zero-one law argument, we then conclude the proof of Theorem 1.3.

4.1. Killing off all but one geodesic with positive probability. The following lemma allows us to kill off all of the geodesics which do not hit a specified boundary arc $I \subset \partial \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$. For the statement, we recall the definition of the internal diameter metric d^U for $U \subset \mathbb{C}$ a connected open set such that $\mathbb{C} \setminus U$ is compact, as in the discussion just after (29).

LEMMA 4.1. For each A>1, $\mathbb{r}>0$, $\epsilon\in(0,(A-1)/100)$ and $p\in(0,1)$, there exists $\mathfrak{p}=\mathfrak{p}(A,\epsilon,p)>0$ such that the following is true. Let $\tau_{\mathbb{r}}$ be as in (33). Let $I\subset\partial\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}$ be a closed boundary arc, chosen in a manner depending only on $(\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}},h|_{\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}})$, with the property that the $d^{\mathbb{C}\setminus\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}}$ -neighborhood $\mathcal{B}_{\epsilon\mathbb{r}}(\partial\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}\setminus I;d^{\mathbb{C}\setminus\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}})$ does not disconnect I from ∞ in $\mathbb{C}\setminus\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}$. With probability at least p, it holds with conditional probability at least p given $(\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}},h|_{\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}})$ that every D_h -geodesic from 0 to a point of $\mathbb{C}\setminus B_{A\mathbb{r}}(0)$ passes through I.

The proof of Lemma 4.1 is similar to that of Lemma 3.7 but simpler, since we do not need a quantitative bound on probabilities, so we only need to define one event rather than defining an event in every Euclidean annulus. See Figure 9 for an illustration of the proof.

To lighten notation, write

(82)
$$\mathcal{B}^* := \mathcal{B}_{\epsilon \mathbb{r}/4} (\partial \mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet} \setminus I; d^{\mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet}}).$$

By Axiom V, we can find $c = c(A, \epsilon, p) > 0$ such that with probability at least 1 - (1 - p)/3, each path in $B_{A\mathbb{r}}(0)$ with Euclidean diameter at least $\epsilon \mathbb{r}/4$ has D_h -length at least $c\mathfrak{c}_{\mathbb{r}}e^{\xi h_{\mathbb{r}}(0)}$. By the definition of \mathcal{B}^* , each path in $\mathbb{C} \setminus \mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}$ from $\partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}} \setminus I$ to a point of $\mathbb{C} \setminus (\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}} \cup \mathcal{B}^*)$ has Euclidean diameter at least $\epsilon \mathbb{r}/4$. Hence, with probability at least 1 - (1 - p)/3,

$$(83) D_h(\partial \mathcal{B}_{\tau_{\mathbb{I}^r}}^{\bullet} \setminus I, \mathbb{C} \setminus (\mathcal{B}_{\tau_{\mathbb{I}^r}}^{\bullet} \cup \mathcal{B}^*); \mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{I}^r}}^{\bullet}) \ge c \mathfrak{c}_{\mathbb{I}^r} e^{\xi h_{\mathbb{I}^r}(0)}.$$

Define the collection of $\delta \mathbb{r} \times \delta \mathbb{r}$ squares $S_{\delta \mathbb{r}}(B_{A\mathbb{r}}(0)) = S_{\delta \mathbb{r}}^0(B_{A\mathbb{r}}(0))$ with corners in $\delta \mathbb{r} \mathbb{Z}^2$ as in (34) with z = 0. Again, by Axiom V we can find $\delta = \delta(c, A, \epsilon) \in (0, \epsilon/100)$ such that with probability at least 1 - (1 - p)/3,

(84)
$$\sup_{u,v \in S} D_h(u,v) \le \frac{c}{100} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)} \quad \forall S \in \mathcal{S}_{\delta\mathbb{F}} (B_{A\mathbb{F}}(0)).$$

Let $\mathcal{U}_{\mathbb{F}}$ be the (finite) set of subdomains U of $B_{A\mathbb{F}}(0)$ such that $B_{A\mathbb{F}}(0) \setminus U$ is a finite union of sets of the form $S \cap B_{A\mathbb{F}}(0)$ for $S \in \mathcal{S}_{\delta\mathbb{F}}(B_{A\mathbb{F}}(0))$. For $U \in \mathcal{U}_r$, let \mathfrak{h}^U be the harmonic part of $h|_U$. Also, let $U_{\delta\mathbb{F}/4}$ be the set of points in U which lie at Euclidean distance at least $\delta\mathbb{F}/4$ from ∂U . Since there are only finitely many sets in $\mathcal{U}_{\mathbb{F}}$ and by the translation and scale invariance of the law of h, modulo additive constant, we can find $C = C(\delta, A, \epsilon) > 0$ such that with probability at least 1 - (1 - p)/3, it holds simultaneously for each $U \in \mathcal{U}_{\mathbb{F}}$ that

(85)
$$\sup_{u \in U_{\delta_{\mathbb{F}}/4}} \left| \mathfrak{h}^{U}(u) - h_{\mathbb{F}}(0) \right| \leq C.$$

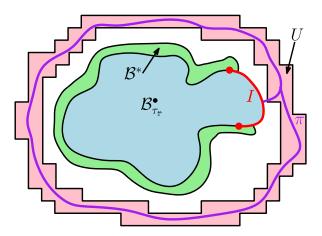


FIG. 9. Illustration of the proof of Lemma 4.1. On the high-probability event $E^U_{\mathbb{T}}$, the D_h -distance from the $d^{\mathbb{C}\setminus \mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}}$ -neighborhood $\mathcal{B}^*=\mathcal{B}_{\in\mathbb{T}/4}(\partial\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}\setminus I;d^{\mathbb{C}\setminus \mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}})$ to $\partial\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}\setminus I$ is bounded below, the D_h -diameter of ever $\delta\mathbb{T}\times\delta\mathbb{T}$ square with corners in $\delta\mathbb{T}\mathbb{Z}^2$ which intersects $B_{A\mathbb{T}}(0)$ is small, and the absolute value of the harmonic part of $h|_U$ is bounded above on the set $U_{\delta\mathbb{T}/4}$ of points in U lying at distance at least $\delta\mathbb{T}/4$ from ∂U . In order to force every D_h -geodesic from 0 to a point of $\mathbb{C}\setminus B_{A\mathbb{T}}(0)$ to go through I, we consider a path π started from I which disconnects $\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}$ from $\partial B_{A\mathbb{T}}(0)$ and let U be the set of $\delta\mathbb{T}\times\delta\mathbb{T}$ squares as above which intersect π but not $\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}$. Using the Markov property of the field and the FKG inequality (Lemma 4.2), we see that with uniformly positive conditional probability given $(\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}, h|_{\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}})$, the D_h -distance from every point of U to I is smaller than the D_h -distance from $\partial\mathcal{B}^*$ to $\partial\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}\setminus I$. This forces every D_h -geodesic to 0 which crosses U to enter $\mathcal{B}^{\bullet}_{\mathbb{T}_{\mathbb{T}}}$ at a point of I.

For a given choice of $U \in \mathcal{U}_{\mathbb{F}}$, let $E_{\mathbb{F}}^U$ be the event that (83), (84) and (85) all hold, so that

(86)
$$\mathbb{P}\bigg[\bigcap_{U\in\mathcal{U}_{\mathbb{F}}}E_{\mathbb{F}}^{U}\bigg]\geq p.$$

The reason for considering $E_{\mathbb{r}}^U$ instead of $\bigcap_{U \in \mathcal{U}_{\mathbb{r}}} E_{\mathbb{r}}^U$ is the same as in Section 3.2[it is easier to condition on $E_{\mathbb{r}}^U$, as explained in the following lemma:

LEMMA 4.2. There is a constant $\mathfrak{p} = \mathfrak{p}(A, \epsilon, p) > 0$ such that the following is true. Also, let $\mathcal{V}(U)$ denote the set of connected components of U. On the event that $U \cap (\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \cup \mathcal{B}^{*}) = \emptyset$, a.s.

(87)
$$\mathbb{P}\left[\max_{V \in \mathcal{V}(U)} \sup_{u,v \in V} D_h(u,v) \leq \frac{c}{2} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(z)} |h|_{\mathbb{C} \setminus U}, E_{\mathbb{F}}^U\right] \geq \mathfrak{p}.$$

PROOF. Since $\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$ is a local set for h and \mathcal{B}^* is determined by $(\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}, h|_{\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}})$, the event $\{U \cap (\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \cup \mathcal{B}^*) = \varnothing\}$ is determined by $h|_{\mathbb{C}\setminus U}$. Moreover, the intersection of this event with the events in each of (83) and (85) is also determined by $h|_{\mathbb{C}\setminus U}$. Hence, on the event $U \cap (\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \cup \mathcal{B}^*) = \varnothing$, the conditional law of $h|_{U}$ given $h|_{U}$ and $E_{\mathbb{F}}^{U}$ is the same as its conditional law given $h|_{U}$ and the event from (84). With this observation in hand, the lemma follows from the Markov property of the GFF and the FKG inequality (Proposition 2.8) via exactly the same argument used to prove Lemma 3.3. \square

PROOF OF LEMMA 4.1. Step 1: Choosing a random domain U. We first choose the domain U to which we will apply Lemma 4.2. The choice will depend on $\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$ and I which is why we need a lower bound for the probability of the intersection of all of the $E_{\mathbb{F}}^U$'s in (86).

Since $\mathcal{B}_{\epsilon \mathbb{F}}(\partial \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \setminus I; d^{\mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}})$ does not disconnect I from ∞ and $\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \subset \overline{B_{\mathbb{F}}(0)}$, we can choose, in a manner depending only on $\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$ and I, a path π in $B_{(A-\epsilon)\mathbb{F}}(0) \setminus \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$ which starts

from a point of I, disconnects $\partial B_{\mathbb{I}}(0)$ from $\partial B_{(A-\epsilon)\mathbb{I}}(0)$ and lies at $d^{\mathbb{C}\setminus\mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}}$ -distance at least $\epsilon_{\mathbb{I}}$ from $\partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}\setminus I$. This path is shown in purple in Figure 9. Let U be the interior of the union of all of the $\delta_{\mathbb{I}} \times \delta_{\mathbb{I}}$ squares $S \in \mathcal{S}_{\delta_{\mathbb{I}}}(B_{A\mathbb{I}}(0))$ which intersect π but do not intersect $\mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}$. Then, $U \in \mathcal{U}_{\mathbb{I}}$, as defined just above (85).

By definition, $U \cap \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} = \varnothing$. We claim that also $U \cap \mathcal{B}^* = \varnothing$. Indeed, each of the $\delta \mathbb{F} \times \delta \mathbb{F}$ squares S in the union defining U is contained in $\mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$ and has Euclidean diameter at most $\sqrt{2}\delta \mathbb{F} < \epsilon \mathbb{F}/4$. If one of these squares intersected \mathcal{B}^* , then by the triangle inequality and the definition 29 of $d^{\mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}}$, the $d^{\mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}}$ -distance from π to $\partial \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \setminus I$ would be at most $\epsilon \mathbb{F}/2$, contrary to the definition of π .

Hence, $U \cap (\mathcal{B}_{\tau_{\mathbb{I}}}^{\bullet} \cup \mathcal{B}^{*}) = \emptyset$. Since $\mathcal{B}_{\tau_{\mathbb{I}}}^{\bullet}$ is a local set for h (Lemma 2.1) and π is determined by $(\mathcal{B}_{\tau_{\mathbb{I}}}^{\bullet}, h|_{\mathcal{B}_{\tau_{\mathbb{I}}}^{\bullet}})$, for $\mathfrak{U} \in \mathcal{U}_{\mathbb{I}}$ the event $\{U = \mathfrak{U}\}$ is determined by $h|_{\mathbb{C}\setminus\mathfrak{U}}$. Therefore, the bound (87) of Lemma 4.2 holds a.s. for our (random) choice of U.

Step 2: Bounding conditional probabilities. By (86), we have $\mathbb{P}[E_{\mathbb{T}}^U] \geq p$, so Markov's inequality (applied to the random variable $\mathbb{P}[(E_{\mathbb{T}}^U)^c | \mathcal{B}_{\tau_{\mathbb{T}}}^{\bullet}, h|_{\mathcal{B}_{\tau_{\mathbb{T}}}^{\bullet}}]$, which has mean at most 1-p), implies that

(88)
$$\mathbb{P}\left[\mathbb{P}\left[E_{\mathbb{I}}^{U} \middle| \mathcal{B}_{\tau_{\mathbb{I}}}^{\bullet}, h \middle|_{\mathcal{B}_{\tau_{\mathbb{I}}}^{\bullet}}\right] \ge 1 - (1-p)^{1/2}\right] \ge 1 - (1-p)^{1/2}.$$

By the locality of $h|_{\mathcal{B}_{\mathbb{T}_{\mathbb{F}}}^{\bullet}}$ (Lemma 2.1), $\sigma(\mathcal{B}_{\mathbb{T}_{\mathbb{F}}}^{\bullet}, h|_{\mathcal{B}_{\mathbb{T}_{\mathbb{F}}}^{\bullet}}) \subset \sigma(h|_{\mathbb{C}\setminus U})$. By Lemma 4.2 and the preceding sentence, the conditional probability given $(\mathcal{B}_{\mathbb{T}_{\mathbb{F}}}^{\bullet}, h|_{\mathcal{B}_{\mathbb{T}_{\mathbb{F}}}^{\bullet}})$ and $E_{\mathbb{F}}^{U}$ that

(89)
$$\max_{u,v \in \mathcal{V}(U)} \sup_{u,v \in V} D_h(u,v) \le \frac{c}{2} \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(z)}$$

is at least \mathfrak{p} , where $\mathcal{V}(U)$ is as in Lemma 4.2. By this and (86) and since p can be made arbitrarily close to 1, to conclude the proof of Lemma 4.1 we only need to show that if $E_{\mathbb{F}}^U$ occurs and (89) holds, then every D_h -geodesic from 0 to a point of $\mathbb{C} \setminus B_{A\mathbb{F}}(0)$ passes through I. This will be accomplished via a similar argument to the proof of Lemma 3.6, as we now explain.

Step 3: Preventing D_h -geodesics from hitting $\partial \mathcal{B}_{\tau_{\Gamma}}^{\bullet} \setminus I$. By the definitions of π and U, each square $S \in \mathcal{S}_{\delta \Gamma}(B_{A\Gamma}(0))$ hit by π which is *not* included in the union defining U intersects I. Note that such a square cannot intersect $\partial \mathcal{B}_{\tau_{\Gamma}}^{\bullet} \setminus I$ without intersecting I since then the $d^{\mathbb{C} \setminus \mathcal{B}_{\tau_{\Gamma}}^{\bullet}}$ -distance from π to $\partial \mathcal{B}_{\tau_{\Gamma}}^{\bullet} \setminus I$ would be smaller than $2\sqrt{2}\delta_{\Gamma} < \epsilon_{\Gamma}$, contrary to the definition of π . Since π is connected, it follows that each connected component of U shares a boundary point with a square $S \in \mathcal{S}_{\delta \Gamma}(B_{A\Gamma}(0))$ which intersects I. The event in (84) shows that this square has D_h -diameter at most $\frac{c}{100}\mathfrak{c}_{\Gamma}e^{\xi h_{\Gamma}(0)}$. This together with the event in (89) shows that

(90)
$$\sup_{u \in U} D_h(u, I) < c \mathfrak{c}_{\mathbb{F}} e^{\xi h_{\mathbb{F}}(0)}.$$

By our choice of U, each path P from a point outside of $B_{A\mathbb{r}}(0)$ to 0, which first hits $\partial \mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet}$ at a point not in I, must pass through U and then must subsequently cross from a point of $\mathbb{C}\setminus(\mathcal{B}_*\cup\mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet})$ to $\partial \mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet}\setminus I$. By (90), the D_h -distance from the first point of U hit by P to 0 is strictly smaller than $\tau_{\mathbb{r}}+c\mathfrak{c}_{\mathbb{r}}e^{\xi h_{\mathbb{r}}(0)}$. On the other hand, (83) shows that the D_h -length of the segment of P, which crosses from $\mathbb{C}\setminus(\mathcal{B}_*\cup\mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet})$ to $\partial \mathcal{B}_{\tau_{\mathbb{r}}}^{\bullet}\setminus I$, is at least $c\mathfrak{c}_{\mathbb{r}}e^{\xi h_{\mathbb{r}}(0)}$, so the D_h -length of the segment of P after it first hits U is at least $\tau_{\mathbb{r}}+c\mathfrak{c}_{\mathbb{r}}e^{\xi h_{\mathbb{r}}(0)}$. Therefore, P cannot be a D_h -geodesic. \square

4.2. *Killing off all but one geodesic almost surely*. By combining Theorem 3.1 and Lemma 4.1, we get the following lemma:

LEMMA 4.3. There exists M > 1 and $\mathfrak{q} > 0$, depending only on ξ , such that for each $\mathfrak{r} > 0$, it holds with probability at least \mathfrak{q} that there is a single point $x \in \partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{I}}}$ which is hit by every leftmost D_h -geodesic from 0 to $\partial \mathcal{B}^{\bullet}_{\tau_{M_{\mathbb{I}}}}$.

PROOF. By Theorem 3.1 applied with $\tau = \tau_{\mathbb{F}}$ (and Axiom V to lower-bound $\tau_{2\mathbb{F}} - \tau_{\mathbb{F}}$), there is an $N \in \mathbb{N}$ such that, for each $\mathbb{F} > 0$, it holds with probability at least 7/8 that there are only N points of $\partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{F}}}$ which are hit by leftmost D_h -geodesics from 0 to $\partial \mathcal{B}^{\bullet}_{\tau_{2\mathbb{F}}}$. Let $X \subset \partial \mathcal{B}^{\bullet}_{\tau_{\mathbb{F}}}$ be the set of such points, and note that X is determined by $(\mathcal{B}^{\bullet}_{\tau_{2\mathbb{F}}}, h|_{\tau_{2\mathbb{F}}}^{\bullet})$. For $x \in X$, let I_X be the set of $y \in \partial \mathcal{B}^{\bullet}_{\tau_{2\mathbb{F}}}$ for which there is a leftmost D_h -geodesic from 0 to y which passes through x. By Lemma 2.7, the I_X 's are disjoint arcs of $\partial \mathcal{B}^{\bullet}_{\tau_{2\mathbb{F}}}$ and by the definition of X their union is all of $\partial \mathcal{B}^{\bullet}_{\tau_{2\mathbb{F}}}$.

By Axiom V, there is an $a \in (0, 1)$ such that with probability at least 7/8, $B_{a\mathbb{F}}(0) \subset \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$. If this is the case, then in the notation (27) we have outrad $(\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet})/\operatorname{inrad}(\mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}) \leq 1/a$. Lemma 2.17 applied with $K = \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}$ and $\mathcal{I} = \{I_x\}_{x \in X}$ therefore shows that there exists $\epsilon = \epsilon(a, N) > 0$ such that whenever $\#X \leq N$ and $B_{a\mathbb{F}}(0) \subset \mathcal{B}_{\tau_{2\mathbb{F}}}^{\bullet}$ (which happens with probability at least 3/4), there exists $x \in X$ such that the following is true. The arc I_x is not disconnected from ∞ in $\mathbb{C} \setminus \mathcal{B}_{\tau_{2\mathbb{F}}}^{\bullet}$ by $\mathcal{B}_{\epsilon\mathbb{F}}(\partial \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet} \setminus I; d^{\mathbb{C} \setminus \mathcal{B}_{\tau_{\mathbb{F}}}^{\bullet}})$. Choose such an x in a manner depending only on $(\mathcal{B}_{\tau_{2\mathbb{F}}}^{\bullet}, h|_{\tau_{2\mathbb{F}}}^{\bullet})$.

By Lemma 4.1 (applied with $2\mathbb{r}$ in place of \mathbb{r} and with A=3/2), we can find $\mathfrak{p}=\mathfrak{p}(\epsilon)>0$ such that with probability at least 3/4, it holds with conditional probability at least \mathfrak{p} given $(\mathcal{B}_{\tau_{2\mathbb{r}}}^{\bullet}, h|_{\mathcal{B}_{\tau_{2\mathbb{r}}}}^{\bullet})$ that every D_h -geodesic from 0 to a point of $\mathbb{C}\setminus B_{3\mathbb{r}}(0)$ passes through I_x . By the definition of I_x and the estimates for the probabilities of the events above, it holds probability at least $\mathfrak{p}/4$ that every leftmost D_h -geodesic from 0 to a point of $\mathbb{C}\setminus B_{3\mathbb{r}}(0)$ passes through x.

By Axiom V, we can find M > 4 such that for every $\mathbb{r} > 0$, it holds with probability at least $1 - \mathfrak{p}/8$ that $\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}} \supset B_{4\mathbb{r}}(0)$. Then, with probability at least $\mathfrak{p}/4 - \mathfrak{p}/8 = \mathfrak{p}/8$, each leftmost D_h -geodesic from 0 to $\partial \mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$ passes through x. Hence, the statement of the lemma is true with $\mathfrak{q} = \mathfrak{p}/8$. \square

We now upgrade from the statement that all *leftmost* D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau_{Mr}}^{\bullet}$ hit the same point of $\partial \mathcal{B}_{\tau_{mr}}^{\bullet}$ to the statement that *all* D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau_{Mr}}^{\bullet}$ coincide until they hit $\partial \mathcal{B}_{\tau_{mr}}^{\bullet}$.

LEMMA 4.4. There exists M > 1 and $\mathfrak{q} > 0$, depending only on ξ , such that for each $\mathfrak{r} > 0$, it holds with probability at least \mathfrak{q} that any two D_h -geodesics from 0 to a point of $\mathbb{C} \setminus \mathcal{B}_{\tau_{Mr}}^{\bullet}$ coincide on the time interval $[0, \tau_{\mathfrak{r}}]$.

PROOF. By Lemma 4.3, there exists M and \mathfrak{q} as in the statement of the lemma such that for each $\mathbb{r}>0$, it holds with probability at least \mathfrak{q} that there is a single point $x\in\partial\mathcal{B}^{\bullet}_{\tau_{\mathbb{r}}}$ which is hit by every leftmost D_h -geodesic from 0 to $\partial\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$. Henceforth, assume that this is the case. By the uniqueness of geodesics to points in \mathbb{Q}^2 (Lemma 2.2), every D_h -geodesic from 0 to a point in $\mathbb{Q}^2\setminus\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$ must pass through x. By Lemma 2.4, leftmost and rightmost geodesics can be approximated by geodesics to points in $\mathbb{Q}^2\setminus\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$, so it follows that every rightmost geodesic and every leftmost geodesic from 0 to $\partial\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$ must pass through x. This implies that a.s. every D_h -geodesic from 0 to $\partial\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$ must pass through x, so the restriction of any such geodesic to the time interval $[0, \tau_{\mathbb{r}}]$ is a D_h -geodesic from 0 to x.

We now argue that there is only one D_h -geodesic from 0 to x, so that any two D_h -geodesics from 0 to $\partial \mathcal{B}_{\tau_{M_{\mathbb{F}}}}^{\bullet}$ coincide on $[0, \tau_{\mathbb{F}}]$. To this end, choose a point $q \in \mathbb{Q}^2 \setminus \mathcal{B}_{\tau_{M_{\mathbb{F}}}}^{\bullet}$. By Lemma 2.2, the D_h -geodesic from 0 to q is a.s. unique. If there were more than one

 D_h -geodesic from 0 to x, then, by concatenating such geodesics with a fixed geodesic from x to q, we would get multiple distinct geodesics from 0 to q, so there must be only one D_h -geodesic from 0 to x. \square

We now conclude the proof via a zero-one law argument.

PROOF OF THEOREM 1.3. By Axiom IV and the translation invariance of the law of h, viewed modulo additive constant, we can assume without loss of generality that z=0. Fix M and \mathfrak{q} as in Lemma 4.4, and for $\mathfrak{r}>0$ let $F_{\mathbb{r}}$ be the event that any two D_h -geodesics from 0 to a point of $\partial \mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}$ coincide on the time interval $[0, \tau_{\mathbb{r}}]$. Then, $\mathbb{P}[F_{\mathbb{r}}] \geq \mathfrak{q}$ and $F_{\mathbb{r}}$ is determined by $(\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}, h|_{\mathcal{B}^{\bullet}_{\tau_{M\mathbb{r}}}})$ and hence by $h|_{B_{M\mathbb{r}}(0)}$. Since the tail σ -algebra $\bigcap_{\mathbb{r}>0} \sigma(h|_{B_{\mathbb{r}}(0)})$ is trivial, there a.s. exists arbitrarily small values of $\mathbb{r}>0$ for which $F_{\mathbb{r}}$ occurs.

Since each LQG metric ball centered at 0 contains a Euclidean ball centered at 0, it is a.s. the case that for each s>0, there exists $\mathbb{r}>0$ for which $\mathcal{B}_{\tau_{M\mathbb{r}}}^{\bullet}\subset\mathcal{B}_{s}(0;D_{h})$ and $F_{\mathbb{r}}$ occurs. Then, any two D_{h} -geodesic from 0 to a point outside of $\mathcal{B}_{s}(0;D_{h})$ coincide on $[0,\tau_{\mathbb{r}}]$, so the theorem is true with $t=\tau_{\mathbb{r}}$. \square

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