Non-extinction of a Fleming-Viot particle model

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Abstract We consider a branching particle model in which particles move inside a Euclidean domain according to the following rules. The particles move as independent Brownian motions until one of them hits the boundary. This particle is killed but another randomly chosen particle branches into two particles, to keep the population size constant. We prove that the particle population does not approach the boundary simultaneously in a finite time in some Lipschitz domains. This is used to prove a limit theorem for the empirical distribution of the particle family.

Keywords Brownian motion · Branching particle system

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1 Introduction

The paper is concerned with a branching particle system $\mathbf{X}_t = (X_t^1, \dots, X_t^N)$ in which individual particles X^j move as N independent Brownian motions and die

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when they hit the complement of a fixed domain $D \subset \mathbb{R}^d$. To keep the population size constant, whenever any particle X^j dies, another one is chosen uniformly from all particles inside D, and the chosen particle branches into two particles. Alternatively, the death/branching event can be viewed as a jump of the *j*-th particle. See Sect. 5 for a more detailed description of the construction.

Let τ_k be the time of the *k*-th jump of \mathbf{X}_t . Since the distribution of the hitting time of ∂D by Brownian motion has a continuous density, only one particle can hit ∂D at time τ_k , for every *k*, a.s. The construction of the process is elementary for all $t < \tau_{\infty} = \lim_{k \to \infty} \tau_k$. However, there is no obvious way to continue the process \mathbf{X}_t after the time τ_{∞} if $\tau_{\infty} < \infty$. Hence, the question of the finiteness of τ_{∞} is interesting. Theorem 1.1 in [10] asserts that $\tau_{\infty} = \infty$, a.s., for every domain *D*. Unfortunately, the proof of that theorem contains an irreparable error (see Example 5.7 below). The cited theorem might be true but it appears to be much harder to prove than the original incorrect argument might have suggested. Example 5.7 given below shows that result cannot be generalized to arbitrary Markov processes. We will show in Remark 5.6 that the other main results in [10], i.e., Theorems 1.3 and 1.4 hold true and an argument showing that $\tau_{\infty} = \infty$, a.s., in domains satisfying the internal ball condition is implicit in the proof of Theorem 1.4 of [10].

In this article, we will prove that $\tau_{\infty} = \infty$, a.s., if the domain $D \subset \mathbb{R}^d$ is Lipschitz with a Lipschitz constant *c* depending on *d* and the number *N* of particles—see Theorem 5.4 and Remark 5.5 below. In addition, we prove theorems on existence and the form of the stationary distribution of the process \mathbf{X}_t , generalizing those in [10]—see Sect. 7.

We use this attempt to rectify an error in an earlier paper to introduce two new techniques. In the end, these techniques may have greater interest or significance than the main theorems. The first technique, developed in Sect. 4, is the construction of a process of Brownian excursions in a cone, with all excursions starting at the vertex. Such a process exists only in cones with certain angles. The construction is combined with a coupling argument to provide a "lower bound" for X_t , in an appropriate sense. The process constructed from Brownian excursions is simpler to analyze than X_t .

The second technique is a new type of boundary Harnack principle (see Sect. 3). The standard boundary Harnack principle compares two functions satisfying a PDE with the same operator, for example, Laplacian, and different boundary conditions. Our new version of the boundary Harnack principle compares a harmonic function with a function u satisfying $\Delta u = -1$. The reason for proving the new form of the boundary Harnack principle is that it allows one to compare certain probabilities and expectations, and then use a method of proof that goes back at least to Davis [12]. The "new boundary Harnack principle" has been proved independently by Atar et al. [3], together with a number of other interesting theorems. We include a full proof of the new boundary Harnack principle because it is different from that in [3], and ours is amenable to generalizations that will be the subject of a forthcoming article.

Both techniques mentioned above—the Brownian excursion process and the boundary Harnack principle—are limited to Lipschitz domains and, moreover, the Lipschitz constant has to satisfy a certain inequality. A natural question arises whether such special Lipschitz domains are the largest natural family of sets where our results hold. It turns out that they are not. In the last section of the paper we will show that, for the two particle process, $\tau_{\infty} = \infty$, a.s., in all polyhedral domains, with arbitrary angles between the faces of the boundary. Unfortunately, our method cannot be easily adapted to the multiparticle case, so we leave this generalization as an open problem.

For some related results on Fleming-Viot type models in smooth domains, see [16] and references therein. The discrete version of the model is studied in [2]; see also references in that paper.

2 Preliminaries

For $y = (y^1, \ldots, y^d) \in \mathbb{R}^d$, let |y| denote the Euclidean norm of y and let $\tilde{y} = (y^1, \ldots, y^{d-1})$. We will denote the open ball with center x and radius r by B(x, r). The closure of a set A will be denoted \overline{A} and its interior will be denoted IntA. All constants, typically denoted by c with or without subscript, are assumed to be strictly positive and finite.

A function $F : \mathbb{R}^{d-1} \to \mathbb{R}$ is called Lipschitz if there exists a constant L such that

$$|F(x) - F(y)| \le L|x - y|, \quad x, y \in \mathbb{R}^{d-1}.$$

Any constant L satisfying the above condition will be called a Lipschitz constant of F.

Consider a bounded connected open set $D \subset \mathbb{R}^d$, $d \ge 2$. We will call D a Lipschitz domain with Lipschitz constant L if ∂D can be covered by a finite number of open balls B_1, \ldots, B_n such that for every $i = 1, \ldots, n$, there exists a Lipschitz function $F_i : \mathbb{R}^{d-1} \to \mathbb{R}$ with Lipschitz constant L, and an orthonormal coordinate system CS_i such that

$$D \cap B_i = \left\{ (y^1, \dots, y^d) \text{ in } CS_i : y^d > F_i(\widetilde{y}) \right\} \cap B_i.$$

The following Harnack principles can be found in [5].

Theorem 2.1 (Harnack inequality)

(a) Suppose 0 < r < R. There exists c = c(r, R, d) such that if u is nonnegative and harmonic in $B(0, R) \subset \mathbb{R}^d$ and $x, y \in B(0, r)$, then

$$u(x) \le c \, u(y).$$

(b) Suppose that $D \subset \mathbb{R}^d$ is a domain and $x, y \in D$ can be connected by a curve $\gamma \subset D$ such that $\inf_{z \in \gamma} \operatorname{dist}(z, \partial D) \geq R$. There exists $c = c(\gamma, R, d)$ such that if u is nonnegative and harmonic in D, then

$$u(x) \le c \, u(y).$$

Theorem 2.2 (Boundary Harnack principle) Suppose D is a connected Lipschitz domain. Suppose V is open, M is compact and $M \subset V$. Then there exists a constant c = c(M, V, D) such that if u and v are two positive and harmonic functions

on D that both vanish continuously on $V \cap \partial D$, then

$$\frac{u(x)}{v(x)} \le c \, \frac{u(y)}{v(y)}, \quad x, y \in M \cap D.$$

The next theorem is a simplified version of Theorem 1 of [1].

Theorem 2.3 Assume that D is a Lipschitz domain. Then there exist constants $r_0 = r_0(D) > 0$, $c = c(D) < \infty$ and a = a(D) > 1 such that if $z \in \partial D$ and $0 < r \le r_0$ then for all functions u and v that are bounded, positive and harmonic on $D \cap B(z, ar)$, and vanishing continuously on $\partial D \cap B(z, ar)$, we have

$$\frac{u(x)}{v(x)} \le c \frac{u(y)}{v(y)}, \quad x, y \in D \cap B(z, r).$$

Remark 2.4 Theorem 2.3 can be used to estimate the constant c(M, V, D) in Theorem 2.2 as follows. Suppose that r_0 and a are as in Theorem 2.3 and we can find balls $B_i(x_i, r_i), i = 1, ..., n$, and $B'_j(y_j, \rho), j = 1, ..., m, \rho > 0, r_i \le r_0, x_i \in \partial D, y_j \in D, M \subset \bigcup_i B_i(x_i, r_i) \cup \bigcup_j B'_j(y_j, \rho)$, and $\bigcup_i B_i(x_i, ar_i) \subset V$ and $\bigcup_j B'_j(y_j, 2\rho) \subset D$. A simple chaining argument based on Theorems 2.1 and 2.3 then shows that the constant c(M, V, D) in Theorem 2.2 depends only on n, m and D.

Next we recall some notation and results from [11]. Fix $d \ge 2$ and p > 0. Let

$$h(\theta) = h_{p,d}(\theta) = F(-p, p+d-2; (d-1)/2; (1-\cos\theta)/2), \qquad (2.1)$$

where

$$F(a,b;c;x) = \sum_{k=0}^{\infty} \frac{(a)_k(b)_k}{(c)_k k!} x^k, \quad |x| < 1,$$

denotes the hypergeometric function and $(a)_k = a(a + 1) \dots (a + k - 1), (a)_0 = 1$. The function *h* has at least one zero in $(0, \pi)$; let $\theta_{p,d}$ denote the smallest one. The quantity $\theta_{p,d}$ is strictly decreasing in *p* for any fixed $d \ge 2$, and strictly increasing to $\pi/2$ in *d* for any fixed p > 1. In particular, if p = 2, then

$$h_{2,d}(\theta) = 1 - \frac{d}{d-1}\sin^2\theta,$$

 $\theta_{2,d} = \arccos \frac{1}{\sqrt{d}} \operatorname{and} \cot \theta_{2,d} = \frac{1}{\sqrt{d-1}}$. Therefore $\theta_{2,2} = \pi/4$ and p < 2 is equivalent to $\cot \theta_{p,d} < \frac{1}{\sqrt{d-1}}$.

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For $d \ge 2$ and p > 0 we let θ be the angle between y and $(0, \dots, 0, 1)$,

$$K_{p,d} = \left\{ y \in \mathbb{R}^d : y \neq 0, \ 0 \le \theta < \theta_{p,d} \right\},\$$

and let *O* denote the axis of $K_{p,d}$. Obviously p < p' implies $K_{p',d} \subset K_{p,d}$. We will drop the subscripts *p* and *d* and write *K* instead of $K_{p,d}$ whenever there is no danger of confusion.

The function $v(x) = |x|^p h(\theta)$, where *h* is given by (2.1), is positive and harmonic inside *K* and continuous on \overline{K} with v(x) = 0 for $x \in \partial K$.

Let (\mathbb{P}^x, X_t) be *d*-dimensional Brownian motion and for a Borel set $A \subset \mathbb{R}^d$ define

$$T_A = \inf \{ t > 0 : X_t \in A \}.$$
(2.2)

Lemma 2.5 Let *F* denote the intersection of $K = K_{p,d}$ and a hyperplane orthogonal to *O*. Let z_0 be the point of intersection of *O* with *F* and assume that $z_0 \in K$. There exists c = c(p, d) such that for all $z_1, z_2 \in O$ with $|z_0| < |z_1| < |z_2|$, we have

$$\frac{\mathbb{P}^{z_2}\left(T_F < T_{\partial K}\right)}{\mathbb{P}^{z_1}\left(T_F < T_{\partial K}\right)} \ge c \left(\frac{|z_2|}{|z_1|}\right)^{2-d-p}.$$
(2.3)

Proof Let K_* be the unbounded component of $K \setminus F$ and

$$u(z) = \mathbb{P}^{z} \left(T_F < T_{\partial K} \right), \quad z \in K_*.$$

Then *u* is positive and harmonic in K_* and continuous on $\overline{K}_* \setminus (F \cap \partial K)$, with u(z) = 0 for $z \in \partial K \setminus F$. It is easy to see that $u(x) \to 0$ as $|x| \to \infty$.

If $I(x) = x/|x|^2$, then the function $\tilde{u}(x) = |x|^{2-d}u(I(x))$ is positive and harmonic in $\tilde{K} = I(K_*)$ (see Lemma 1.18 of [5]). The function \tilde{u} vanishes continuously on $\partial \tilde{K} \setminus I(F)$. Let $K' = (1/2)\tilde{K}$. Recall that $v(x) = |x|^p h(\theta)$ is positive and harmonic inside K and continuous on \overline{K} with v(x) = 0 for $x \in \partial K$. By the boundary Harnack principle,

$$\frac{\widetilde{u}(z)}{\widetilde{u}(z')} \ge c \frac{v(z)}{v(z')},\tag{2.4}$$

for $z, z' \in K'$, where *c* depends on \widetilde{K} and K' and does not depend on *z* and *z'*. Note that $u(x) = |x|^{2-d} \widetilde{u}(I(x))$. Hence, for $z_1, z_2 \in O \cap I(K')$,

$$\frac{u(z_2)}{u(z_1)} = \frac{|z_2|^{2-d}\widetilde{u}(I(z_2))}{|z_1|^{2-d}\widetilde{u}(I(z_1))} \ge c\frac{|z_2|^{2-d}v(I(z_2))}{|z_1|^{2-d}v(I(z_1))} = c\frac{|z_2|^{2-d}|z_2|^{-p}h(0)}{|z_1|^{2-d}|z_1|^{-p}h(0)} = c\left(\frac{|z_2|}{|z_1|}\right)^{2-p-d}.$$

The inequality holds for all $z_1, z_2 \in O \cap I(K_*)$ (possibly with a different value of *c*) because the function *u* is bounded below and above on $O \setminus I(K')$ by strictly positive and finite constants. This completes the proof of (2.3).

We will use the following estimate in the proof of Lemma 4.1.

Lemma 2.6 There exists a cone $K' \subset K = K_{p,d}$ and a constant c = c(K, K') such that for $x \in K'$ and $t \ge |x|^2$,

$$c^{-1} \left(\frac{t}{|x|^2}\right)^{-\frac{p}{2}} \le \mathbb{P}^x \left(T_{\partial K} > t\right) \le c \left(\frac{t}{|x|^2}\right)^{-\frac{p}{2}}.$$
(2.5)

Proof See [4,6,14] or [22].

3 A boundary Harnack principle

Let $D \subset \mathbb{R}^d$, $d \ge 2$, be a bounded Lipschitz domain and let $A \subset D$ be a compact set with Int $A \neq \emptyset$. For $x \in D$, define

$$f(x) = \mathbb{P}^{x}(T_A < T_{\partial D}),$$

$$g(x) = \mathbb{E}^{x}T_{\partial D}.$$

Theorem 3.1 Assume that the Lipschitz constant L of D satisfies $L < \frac{1}{\sqrt{d-1}}$. Then there exists a constant c = c(A, D) such that for all $x \in D$,

$$\frac{1}{c} \le \frac{f(x)}{g(x)} \le c. \tag{3.1}$$

Remark 3.2 The condition $L < \frac{1}{\sqrt{d-1}}$ is sharp. See Example 3.3 below.

Proof of RHS of (3.1) Since A is compact, $\inf_{x \in A} \operatorname{dist}(x, D^c) = c_1 > 0$. Therefore,

$$\inf_{x \in A} \mathbb{E}^{x} T_{\partial D} \ge \inf_{x \in A} \mathbb{E}^{x} T_{\partial B(x,c_{1})} = c_{2} > 0.$$

By the strong Markov property applied at T_A , we have for $x \in D$,

$$\mathbb{E}^{x} T_{\partial D} \ge c_2 \mathbb{P}^{x} (T_A < T_{\partial D}),$$

which implies the RHS of (3.1).

Proof of LHS of (3.1) Since *D* is a bounded Lipschitz domain with Lipschitz constant $L < \frac{1}{\sqrt{d-1}}$, it is easy to see that there exist $p \in (0, 2)$ and $\rho > 0$ with the following properties.

- (i) dist $(A, \partial D) > 2\rho$.
- (ii) Consider any $x \in D$ with dist $(x, \partial D) < \rho 2^{-5}$. Then there exists $x_0 \in \partial D$ and an orthonormal coordinate system $CS = CS_{x_0}$ with the following properties.

The origin of *CS* is $x_0, K_{p,d} \cap B(x_0, 2\rho) \subset D \cap B(x_0, 2\rho)$, and $x \in O$ (that is, *x* belongs to the axis of $K_{p,d}$). For r > 0 and integer *k*, let

$$E_r^* = \left\{ y \in \mathbb{R}^d \text{ in } CS : |\widetilde{y} - \widetilde{x}_0| \le r \tan(\theta_{p,d}), |y^d - x_0^d| \le r \right\},\$$

$$\widetilde{E}_k = E_{2^{-k}}^*.$$

We can choose x_0 and *CS* so that for some Lipschitz function $F = F_{x_0}$: $\mathbb{R}^{d-1} \to \mathbb{R}$ with Lipschitz constant *L*, and all *k* such that $2^{-k} \le \rho$,

$$D \cap \widetilde{E}_k = \left\{ (y^1, \dots, y^d) \text{ in } CS : y^d > F(\widetilde{y}) \right\} \cap \widetilde{E}_k.$$

We fix $x \in D$ with dist $(x, \partial D) < \rho 2^{-5}$ and the corresponding coordinate system *CS* for the rest of the proof.

Let $E_k = \widetilde{E}_k \setminus \widetilde{E}_{k+1}$ and $C_k = \text{Int}(D \cap E_k)$ for $k = N_0, \dots, N_1$, where

$$N_0 = \min\{k : 2^{-k} \le \rho\}, \qquad N_1 = \max\{k : |x| = x^d \le 2^{-k-3}\}.$$

Also let $C_{N_0-1} = \operatorname{Int} \left(D \setminus \widetilde{E}_{N_0} \right)$ and $C_{N_1+1} = \operatorname{Int} \left(D \cap \widetilde{E}_{N_1+1} \right)$.

Note that $C_i \cap C_j = \emptyset$ if $i \neq j$, and $D = \overline{C}_{N_0-1} \cup \cdots \cup \overline{C}_{N_1+1}$.

Let G(x, y) denote the Green function for Brownian motion killed on exiting *D*. Then

$$g(x) = \mathbb{E}^{x} T_{\partial D} = \int_{D} G(x, y) \, dy = \sum_{k=N_{0}-1}^{N_{1}+1} \int_{C_{k}} G(x, y) \, dy.$$
(3.2)

For $k = N_0, ..., N_1$ denote by y_k the midpoint of the line segment being the intersection of C_k with x^d -axis in *CS*. In other words, $\{y_k\} = \partial E^*_{(3/4)2^{-k}} \cap O$. Fix k and j such that $j \ge 1, k \ge N_0, j + k \le N_1$ and consider the points y_k and y_{k+j} . Let

$$F_k = \overline{C}_k \cap \overline{C}_{k+1} \cap K_{p,d},$$

and

$$u(z) = \mathbb{P}^{z} \Big(T_{F_{k+j}} < T_{\partial K_{p,d}} \Big).$$

By Lemma 2.5,

$$u(y_k) \ge c_1 u(y_{k+j}) \left(\frac{2^{-k}}{2^{-k-j}}\right)^{2-p-d} = c_1 u(y_{k+j}) 2^{j(2-p-d)},$$

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where $c_1 = c_1(p, d)$. By scaling properties of Brownian motion, $u(y_{k+j}) = c_2 = c_2(p, d)$, that is, $u(y_{k+j})$ depends only on p and d. We obtain

$$\mathbb{P}^{z}\left(T_{F_{k+j}} < T_{\partial K_{p,d}}\right) \ge c_{3} 2^{-j(p+d-2)},\tag{3.3}$$

where $c_3 = c_3(p, d)$.

Let

$$v(z) = \mathbb{P}^{z} \left(T_{F_{k+j}} < T_{\partial D} \right).$$

Note that $v(y_{k+j}) \le 1$ and $v(y_k) \ge u(y_k) \ge c_3 2^{-j(p+d-2)}$, by (3.3).

We will apply Theorem 2.2 with $M = \partial E^*_{(3/4)2^{-k-j}}$ and $V = E_{k+j}$. It follows from Remark 2.4 that the constant $c_5 = c(M, V, D)$ may be chosen independent of k and j. The boundary Harnack principle implies that

$$\frac{G(x,z)}{G(x,y_{k+j})} \ge c_5 \frac{v(z)}{v(y_{k+j})},$$
(3.4)

for $z \in D \cap M$. The harmonic functions $G(x, \cdot)$ and v have zero boundary values on $\partial D \setminus \overline{E}^*_{(3/4)2^{-k-j}}$, so the inequality (3.4) extends to all $z \in D \setminus E^*_{(3/4)2^{-k-j}}$, in particular, it applies to $z = y_k$. Hence,

$$\frac{G(x, y_k)}{G(x, y_{k+j})} \ge c_5 \frac{v(y_k)}{v(y_{k+j})} \ge c_5 c_3 2^{-j(p+d-2)} = c_6 2^{-j(p+d-2)}.$$
(3.5)

Now consider the function

$$h_m(z) = \mathbb{P}^z \left(T_{\widetilde{E}_{m+2}} < T_{\partial D} \right).$$

By the scaling properties of Brownian motion, $h_m(y_m) \ge c_7 > 0$ for all $m = N_0, \ldots, N_1$. By the boundary Harnack principle (Theorem 2.2) applied to $u(z) = G(x, z), v(z) = h_m(z), M = \overline{C}_m$ and $V = \text{Int}(\widetilde{E}_{m-1} \setminus E^*_{(3/4)2^{-m-1}})$, we have

$$\frac{G(x, y)}{h_m(y)} \le c_8 \frac{G(x, y_m)}{h_m(y_m)}$$

for $y \in C_m$, where c_8 depends only on D, by Remark 2.4. Therefore, for $y \in C_m$,

$$G(x, y) \le c_8 G(x, y_m) \frac{h_m(y)}{h_m(y_m)} \le c_8 \frac{1}{c_7} G(x, y_m) = c_9 G(x, y_m).$$
(3.6)

This implies

$$\int_{C_{k+j}} G(x, y) \, dy \le c_9 G(x, y_{k+j}) \operatorname{vol}(C_{k+j}) \le c_{10} \, 2^{-d(k+j)} G(x, y_{k+j}), \quad (3.7)$$

where c_{10} depends only on *D*.

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On the other hand, by the usual Harnack inequality,

$$G(x, y) \ge c_{11}G(x, y_k)$$

for $y \in B_k = B(y_k, 2^{-k-1})$, because $B(y_k, 2^{-k-1}) \subset D \setminus \{x\}$. This implies that

$$\int_{C_k} G(x, y) \, dy \ge c_{11} G(x, y_k) \operatorname{vol}(B_k) = c_{12} 2^{-kd} G(x, y_k), \tag{3.8}$$

where c_{12} does not depend on k.

Combining (3.5), (3.7) and (3.8) we have

$$\int_{C_{k+j}} G(x, y) \, dy \le c_{13} 2^{j(p-2)} \int_{C_k} G(x, y) \, dy,$$

where $c_{13} = c_{13}(D)$. Fix q < 1. Since $p \in (0, 2)$, we may choose j so large that $c_{13}2^{j(p-2)} \le q < 1$. Let $a_k = \int_{C_k} G(x, y) \, dy$, then

$$a_{k+j} \le q a_k, \quad k = N_0, \dots, N_1 - j.$$

Let $N_2 = \min(N_1, N_0 + j - 1)$. The last inequality implies that

$$\sum_{k=N_0}^{N_1} a_k = \sum_{k=N_0}^{N_2} \sum_{m=0}^{\infty} a_{k+mj} \mathbf{1}_{\{k+mj \le N_1\}} \le \sum_{k=N_0}^{N_2} \sum_{m=0}^{\infty} a_k q^m = c_{14} \sum_{k=N_0}^{N_2} a_k.$$
(3.9)

Recall that $G(x, \cdot)$ has zero boundary values on ∂D , so it is bounded by $\sup_{z \in C_{N_0}} G(x, z)$ on the set $D \setminus \widetilde{E}_{N_0}$. This and (3.6) imply that $\sup_{z \in D \setminus \widetilde{E}_{N_0}} G(x, z) \leq c_{15}G(x, y_{N_0})$. We use (3.8) to see that

$$a_{N_0-1} = \int_{C_{N_0-1}} G(x, y) \, dy \le c_{15} G(x, y_{N_0}) \operatorname{vol}(C_{N_0-1})$$

$$\le c_{15} \operatorname{vol}(C_{N_0-1}) c_{12}^{-1} 2^{N_0 d} \int_{C_{N_0}} G(x, y) \, dy = c_{16} a_{N_0}.$$
(3.10)

Recall the definition of N_0 to see that c_{16} depends only on D.

The following calculation is presented in the case $d \ge 3$ only. The case d = 2 requires minor modifications and is left to the reader.

Let $\widetilde{G}(x, y)$ denote the Green function for Brownian motion in \mathbb{R}^d , and let $\overline{G}(x, y)$ be the Green function for Brownian motion in $B(x, 2^{-N_1-4})$. It is well known that for $d \ge 3$, $\widetilde{G}(x, y) = c_{17}|x - y|^{2-d}$, where c_{17} depends on d, and $\overline{G}(x, y) =$

 $\widetilde{G}(x, y) - \widetilde{G}(x, z)$, for $y \in B(x, 2^{-N_1-4})$ and $z \in \partial B(x, 2^{-N_1-4})$. It follows that for $|y - x| \le 2^{-N_1-5}$,

$$\overline{G}(x, y) \ge c_{18}\widetilde{G}(x, y). \tag{3.11}$$

We have $G(x, y) \leq \widetilde{G}(x, y)$ for $y \in D$, and $\int_{B(x,r)} \widetilde{G}(x, y) dy = c_{19}r^2$. Therefore,

$$a_{N_{1}+1} = \int_{C_{N_{1}+1}} G(x, y) dy \leq \int_{C_{N_{1}+1}} \widetilde{G}(x, y) dy$$
$$\leq \int_{B(x, \operatorname{diam}(\widetilde{E}_{N_{1}+1}))} \widetilde{G}(x, y) dy = c_{20} 2^{-2N_{1}}.$$
(3.12)

Since $B(x, 2^{-N_1-4}) \subset D$,

$$G(x, y) \ge \overline{G}(x, y). \tag{3.13}$$

Put $y_{N_1+1} = (\tilde{x}, x_d + 2^{-N_1-5})$. Then by (3.11) and (3.13),

$$G(x, y_{N_1+1}) \ge c_{18}\widetilde{G}(x, y) = c_{21}(2^{-N_1})^{2-d}.$$

Moreover, by the usual Harnack inequality,

$$G(x, y) \ge c_{22}G(x, y_{N_1+1}),$$

for $y \in B(y_{N_1}, 2^{-N_1-2})$. Therefore,

$$a_{N_1} = \int_{C_{N_1}} G(x, y) dy \ge \int_{B(y_{N_1}, 2^{-N_1 - 2})} G(x, y) dy$$

$$\ge c_{22} G(x, y_{N_1 + 1}) \operatorname{vol}(B(y_{N_1}, 2^{-N_1 - 2})) \ge c_{23} (2^{-N_1})^{2 - d} \cdot 2^{-N_1 d} = c_{24} 2^{-2N_1}.$$
(3.14)

Combining (3.12) and (3.14), we obtain

$$a_{N_1+1} \le c_{25}a_{N_1}.\tag{3.15}$$

Let $C_* = C_{N_0-1} \cup \cdots \cup C_{N_2}$ and note that $A \subset C_*$. Let $\sigma_{C_*} = \int_0^{T_{\partial D}} 1_{\{X_s \in C_*\}} ds$. Then (3.9), (3.10) and (3.15) imply that

$$\mathbb{E}^{x} T_{\partial D} \le c_{26} \mathbb{E}^{x} \sigma_{C_{*}}.$$
(3.16)

Since *D* is bounded, $\sup_{z \in D} \mathbb{E}^z T_{\partial D} = c_{27} < \infty$. By the strong Markov property applied at the hitting time of C_* , for $z \in D$,

$$\mathbb{E}^{\mathbb{Z}}\sigma_{C_*} \leq c_{27} \mathbb{P}^{\mathbb{Z}}(T_{C_*} < T_{\partial D}).$$

This and (3.16) yield

$$\mathbb{E}^{x} T_{\partial D} \le c_{28} \mathbb{P}^{x} (T_{C_{*}} < T_{\partial D}). \tag{3.17}$$

Consider functions

$$\xi_1(z) = \mathbb{P}^z(T_A < T_{\partial D}),$$

$$\xi_2(z) = \mathbb{P}^z(T_{C_*} < T_{\partial D}).$$

Both functions are positive and harmonic in $D \setminus \overline{C}_*$, and continuous on $\overline{D} \setminus \overline{C}_*$ with u(z) = v(z) = 0 for $z \in \partial D \setminus \overline{C}_*$. We apply the boundary Harnack principle with $V = D \setminus \overline{C}_*$ and $M = \widetilde{E}_{N_2+1}$ to see that

$$\frac{\xi_1(x)}{\xi_2(x)} \ge c_{29} \frac{\xi_1(y_{N_2+1})}{\xi_2(y_{N_2+1})}.$$
(3.18)

We use Remark 2.4 to see that c_{29} may be chosen so that it depends only on D. It follows from the definitions of N_0 , N_2 and j that for some constant c_{30} , we have dist $(y_{N_2+1}, \partial D) > c_{30}$. This implies that $\xi_1(y_{N_2+1}) = \mathbb{P}^{y_{N_2+1}}(T_A < T_{\partial D}) \ge c_{31}$, for some c_{31} depending only on D. We obtain from (3.18) that $\xi_1(x)/\xi_2(x) \ge c_{29}c_{31}$, and this combined with (3.17) gives

$$\mathbb{E}^{x} T_{\partial D} \leq (c_{28}/c_{29}c_{31})\mathbb{P}^{z}(T_{A} < T_{\partial D}).$$

We have proved the LHS of (3.1) for x satisfying dist(x, ∂D) $\leq \rho 2^{-5}$.

It is easy to check that $\inf\{f(x) : \operatorname{dist}(x, \partial D) \ge \rho 2^{-5}\} > 0$ and $\sup\{g(x) : x \in D\}$ < ∞ , so the LHS of (3.1) holds for all $x \in D$.

Example 3.3 The condition $L < \frac{1}{\sqrt{d-1}}$ in Theorem 3.1 is sharp. To see this, note that for any $L > \frac{1}{\sqrt{d-1}}$ there is a p > 2, such that the cone $K = K_{p,d}$ is a Lipschitz domain with the Lipschitz constant L. Let r > 0 be such that for every $x \in O$, $B(x, r|x|) \subset K$. Then $g(x) = \mathbb{E}^x T_{\partial K} \ge \mathbb{E}^x T_{\partial B(x, r|x|)} \ge c_1 r^2 |x|^2$. Recall that $f(x) = \mathbb{P}^x (T_A < T_{\partial K})$ and let $u(x) = |x|^p h_{p,d}(\theta)$. By the boundary Harnack principle applied to f and uin a neighborhood of 0, $f(x) \le c_2 |x|^p$ for $x \in O$, |x| < 1. Since p > 2, we cannot have $f(x) \ge c_3 g(x)$ in a neighborhood of 0. The domain K is unbounded but it is easy to extend the argument to $K \cap B(0, 1)$.

4 Construction of an auxiliary process from Brownian excursions

Let Ω denote the family of all functions $\omega : [0, \infty) \to \mathbb{R}^d \cup \{\delta\}$ continuous up to their lifetime $R(\omega) = \inf \{t \ge 0 : \omega(t) = \delta\}$ and constantly equal to δ for $t \ge R$,

where δ denotes the coffin state outside \mathbb{R}^d . Let *X* be the canonical process on Ω , i.e., $X_t(\omega) = \omega(t)$ and let \mathbb{P}^x denote the distribution of Brownian motion starting from $x \in \mathbb{R}^d$. As in (2.2), for a Borel set $A \subset \mathbb{R}^d$ let $T_A = \inf \{t > 0 : X_t \in A\}$. Let $K = K_{p,d}$ for some p > 0, and let X' denote the process

$$X'_t = \begin{cases} X_t, & \text{for } t < T_{\partial K}, \\ \delta, & \text{otherwise,} \end{cases}$$

i.e., X' is the process X killed on exiting K. If X has the distribution \mathbb{P}^x , then X' is called Brownian motion in K and its distribution is denoted by \mathbb{P}^x_K .

Let *U* denote the family of all functions $\omega : [0, \infty) \to K \cup \{\delta\}$ such that $\omega(0) = 0$, continuous up to their lifetime *R*. Let H^0 denote a standard excursion law of Brownian motion in $K_{p,d}$ starting from 0. Namely, H^0 is a nonnegative and σ -finite measure on Ω such that *X* is strong Markov under H^0 with the \mathbb{P}_K transition probabilities and $H^0(\lim_{t\to 0} X_t \neq 0) = 0$. We have $H^0(\Omega \setminus U) = 0$. The existence of H^0 follows from results of [8] and [20].

Lemma 4.1 There exists $c \in (0, \infty)$ such that

$$H^{0}(R > t) = c t^{-\frac{p}{2}}, \quad t > 0.$$
(4.1)

Proof Let $y_{\varepsilon} = (0, ..., 0, \varepsilon) \in \mathbb{R}^d$ and let $G_K(x, y)$ denote the Green function for *K*. By Theorem 4.1 of [8],

$$H^{0}(R > t) = c_{1} \lim_{\substack{z \to 0 \\ z \in K}} \frac{\mathbb{P}^{z}(T_{\partial K} > t)}{G_{K}(z, y_{1})}.$$
(4.2)

By Theorem 2.2 of [8], which is an improvement of the boundary Harnack principle, there exists $c(K, \varepsilon)$ such that for all functions h_1 and h_2 which are positive and harmonic in K and vanish continuously on ∂K , we have

$$c(K,\varepsilon)^{-1}\frac{h_1(y)}{h_2(y)} \le \frac{h_1(x)}{h_2(x)} \le c(K,\varepsilon)\frac{h_1(y)}{h_2(y)}$$

for all $x, y \in K \cap B(0, \varepsilon)$, and $\lim_{\varepsilon \to 0} c(K, \varepsilon) = 1$. Therefore, the limit

$$\lim_{\substack{z \to 0 \\ z \in K}} \frac{h_1(z)}{h_2(z)}$$

exists and belongs to $(0, \infty)$ for all functions h_1, h_2 satisfying the above assumptions. We apply this claim to $h_1(z) = G_K(z, y_1)$ and $h_2(z) = |z|^p h(\theta)$, to conclude that

$$\lim_{\varepsilon \to 0} \frac{G_K(y_\varepsilon, y_1)}{\varepsilon^p} = c \in (0, \infty),$$

and

$$H^{0}(R > t) = c \lim_{\varepsilon \to 0} \frac{\mathbb{P}^{y_{\varepsilon}}(T_{\partial K} > t)}{\varepsilon^{p}}.$$
(4.3)

By Lemma 2.6,

$$c^{-1}t^{-\frac{p}{2}} \leq \frac{\mathbb{P}^{y_{\varepsilon}}(T_{\partial K} > t)}{\varepsilon^{p}} \leq ct^{-\frac{p}{2}},$$

for $t \ge \varepsilon^2$ which implies $c^{-1}t^{-\frac{p}{2}} \le H^0(R > t) \le ct^{-\frac{p}{2}}$, for $t \ge 0$. Therefore $H^0(R > 1)$ is a positive and finite number.

Now mimicking the proof of Proposition 5.1 of [8], using (4.3) instead of (4.2), we easily see that if $\{X(t), t \ge 0\}$ has the distribution H^0 , then for every a > 0 the scaled process $\{\sqrt{a}X(t/a), t \ge 0\}$ has the distribution $a^{p/2}H^0$. In particular, for every a > 0

$$H^0(R > t) = a^{p/2} H^0(R > at), \quad t \ge 0,$$

and putting a = 1/t we obtain (4.1) with $c = H^0(R > 1)$.

Let λ denote the Lebesgue measure on $\mathbb{R}_+ = [0, \infty)$ and let \mathcal{P} be a Poisson point process on $\mathbb{R}_+ \times U$ with characteristic measure $\lambda \times H^0$, i.e., \mathcal{P} is a random subset of $\mathbb{R}_+ \times U$ such that for every pair A_1, A_2 of disjoint nonrandom subsets of $\mathbb{R}_+ \times U$, card($\mathcal{P} \cap A_1$) and card($\mathcal{P} \cap A_2$) are independent random variables with Poisson distributions with means ($\lambda \times H^0$)(A_1) and ($\lambda \times H^0$)(A_2), respectively [18]. With probability 1, there are no two points with the same first coordinate, and therefore the elements of \mathcal{P} may be unambiguously denoted by (t, e_t) . Let

$$R_t = \inf \{s > 0 : e_t(s) = \delta\}.$$

By abuse of notation, for a generic element e of U we will write

$$R(e) = \inf \{s > 0 : e(s) = \delta\}.$$

Lemma 4.2 *If* $p \in (0, 2)$ *, then for every* s > 0*,*

$$\sum_{t\leq s}R_t<\infty, \quad \text{a.s.}$$

Proof We use Theorem 4.6 of [18]: if $\varphi : \mathbb{R}_+ \times U \to \mathbb{R}_+$ is a measurable function, then

$$\sum_t \varphi(t, e_t) < \infty, \quad \text{a.s.}$$

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iff

$$\iint_{\mathbb{R}_+\times U} (\varphi(t,e)\wedge 1) dt H^0(de) <\infty.$$

In particular, if $\varphi(t, e) = R(e)\mathbf{1}_{[0,s]}(t)$, then

$$\sum_{t\leq s}R_t<\infty,\quad \text{a.s.}$$

iff

$$\iint_{[0,s]\times U} (R(e)\wedge 1)dt H^0(de) < \infty.$$

If we let $U^- = \{e \in U : R(e) \le 1\}$ and $U^+ = \{e \in U : R(e) > 1\}$ then

$$\iint_{[0,s]\times U} (R(e) \wedge 1) dt H^0(de) = s \int_{U^-} R(e) H^0(de) + s H^0(U^+).$$

By Lemma 4.1,

$$H^{0}(U^{+}) = \int_{1}^{\infty} H^{0}(R \in dt) = c \int_{1}^{\infty} t^{-p/2 - 1} dt < \infty,$$

because p > 0, and

$$\int_{U^{-}} RdH^{0} = \int_{0}^{1} tH^{0}(R \in dt) = c \int_{0}^{1} t \cdot t^{-p/2-1}dt < \infty,$$

because p < 2.

Let $\sigma_v = \sum_{s \le v} R_s$ and $\sigma_{v-} = \sum_{u < v} R_u$ for $v \ge 0$. By Lemma 4.2, if $p \in (0, 2)$ then $\sigma_v < \infty$ for all $v < \infty$, a.s.

Lemma 4.3 The process σ is a stable subordinator with index p/2.

Proof The process σ is increasing and has values in $[0, \infty)$. Its paths are rightcontinuous with left limits. Note that $\{(t, R(e_t))\}_{e \in \mathcal{P}}$ is a Poisson point process on $\mathbb{R}_+ \times \mathbb{R}_+$ with characteristic measure $\lambda \times \Pi$, where Π is given by

$$\Pi(dx) = H^0(R \in dx) = c \, x^{-p/2 - 1} dx,$$

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where the last formula follows from Lemma 4.1. This implies that σ is a process with independent and stationary increments, so σ is a Lévy process. Moreover σ is a sub-ordinator, since it has values in $[0, \infty)$ only. We use calculations that can be found in Section 0.5 and on page 73 of [7] to see that the Laplace transform of σ is

$$E \exp(-\lambda \sigma_t) = \exp\left\{-t \int_0^\infty (1 - e^{-\lambda x}) \Pi(dx)\right\}$$
$$= \exp\left\{-ct \int_0^\infty (1 - e^{-\lambda x}) x^{-p/2 - 1} dx\right\}$$
$$= \exp(-ct \lambda^{p/2}).$$

Therefore σ is stable with index p/2.

It is well known that for a stable subordinator σ we have $\lim_{v\to\infty} \sigma_v = \infty$, a.s. So with probability 1, for every $t \ge 0$, the formula $r = \inf \{v \ge 0 : \sigma_v \ge t\}$ defines a unique $r \ge 0$. For $t \ge 0$ let

$$Z_t = \begin{cases} e_r(t - \sigma_{r-}), & \text{if } \sigma_{r-} < \sigma_r \text{ and } t \in (\sigma_{r-}, \sigma_r), \\ 0, & \text{otherwise.} \end{cases}$$
(4.4)

The process Z takes values in $K \cup \{0\}$.

Remark 4.4 The above construction is similar to the classical Itô representation of Brownian motion using the Poisson point process of excursions, see [24, Chap. XII]. The construction of a Markov process from excursions is presented in [25]. The history of the idea, related papers and results are discussed in that article. The process *Z* is strong Markov by [25, Thm. 4.1]—it is straightforward to check that the assumptions of that theorem are satisfied in our case.

Corollary 4.5 Let Z_t^1, \ldots, Z_t^N be jointly independent copies of Z_t defined in (4.4). If $p < 2 - \frac{2}{N}$ and $T \in (1, \infty)$ then $\inf_{1/T \le t \le T} \max_{1 \le i \le N} |Z_t^i| > 0$, a.s.

Proof For each *i*, let σ_t^i be a stable subordinator associated with the process Z_t^i as in Lemma 4.3 and let $A_i = \{t \in [\frac{1}{T}, T] : Z_t^i = 0\}$. In other words, A_i is the range of σ_t^i over $[\frac{1}{T}, T]$. We use the following result of Hawkes [17]: The ranges of two independent stable subordinators with indices α and β intersect if and only if $\alpha + \beta > 1$, in which case the intersection is stochastically equivalent to the range of a stable subordinator of index $\alpha + \beta - 1$. Therefore, by induction, $A_1 \cap \cdots \cap A_N = \emptyset$, a.s., if and only if $\frac{Np}{2} - N + 1 < 0$. This condition holds since $p < 2 - \frac{2}{N}$.

It is easy to see that $t \to |Z_t^i|$ is lower semicontinuous. Hence, $t \to \max_{1 \le i \le N} |Z_t^i|$ is also lower semicontinuous and, therefore, it attains its infimum on $[\frac{1}{T}, T]$. It follows that $\{\inf_{1/T \le t \le T} \max_{1 \le i \le N} |Z_t^i| > 0\} = \{A_1 \cap \cdots \cap A_N = \emptyset\}$. We have shown that the last event has probability one if $p < 2 - \frac{2}{N}$.

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5 Construction of a Fleming-Viot process

We recall the following description of a Fleming-Viot-type particle system from [10]. Consider an open set $D \subset \mathbb{R}^d$ and an integer $N \geq 2$. Let $\mathbf{X}_t = (X_t^1, \dots, X_t^N)$ be a process with values in D^N defined as follows. Let $\mathbf{X}_0 = (x^1, \dots, x^N) \in D^N$. Then the processes X_t^1, \ldots, X_t^N evolve as independent Brownian motions until the time τ_1 when one of them, say, X^j hits the boundary of D. At this time one of the remaining particles is chosen uniformly, say, X^k , and the process X^j jumps at time τ_1 to $X_{\tau_1}^k$. The processes X_t^1, \ldots, X_t^N continue evolving as independent Brownian motions after time τ_1 until the first time $\tau_2 > \tau_1$ when one of them hits the boundary of D. Again at the time τ_2 the particle which approaches the boundary jumps to the current location of a particle chosen uniformly at random from amongst the ones strictly inside D. The subsequent evolution of X proceeds in the same way. The total number of jumps may be finite or infinite. The above recipe defines the process X_t only for $t < \tau_{\infty} = \lim_{k \to \infty} \tau_k$. There is no natural way to define the process \mathbf{X}_t for $t \geq \tau_{\infty}$. Hence, we add a cemetery state δ to the state space and we let $\mathbf{X}_t = \delta$ for all $t \geq \tau_{\infty}$. We define \mathbf{X}_t so that it is right-continuous with left limits on the interval $[0, \tau_{\infty})$. We do not make any *a priori* claims about existence or non-existence of the left limit $\mathbf{X}_{\tau_{\infty}-}$.

Remark 5.1 The proof of the main theorem in this section, Theorem 5.4, involves an inductive construction of the Fleming-Viot process. In particular the proof relies on a special construction of a Brownian motion started in D and stopped on hitting ∂D . In preparation we introduce a sequence of stopped processes which may be used as an alternative construction of a Fleming-Viot process. Let $\mathbf{X}_t = (X_t^1, \dots, X_t^N)$ be a Fleming-Viot process in D. For each particle X_t^i and $n \in \mathbb{N}$ we name a sequence of stopping times $s_n^i = \inf\{t \ge \tau_n : X_{t^-}^i \in \partial D\}$. Notice there is exactly one $j \in \{1, \dots, N\}$ for which $s_n^j = \tau_n$. For every $i \neq j$, s_n^i is the first time after τ_n that X_t^i hits ∂D . Notice that $s_{n+1}^i = s_n^i$ for every $i \neq j$. Now let \mathcal{Q}^n be the distribution of the stopped process $(X_{t \wedge s_n^1}^1, \dots, X_{t \wedge s_n^N}^N)$. Informally we allow the process to evolve until τ_n , at which point all but one of the particles are in the interior of D. After τ_n the remaining particles continue as independent Brownian motions until they are stopped on exiting D. It is easy to construct a process \mathbf{X}_t^1 distributed as \mathcal{Q}^1 because trivially each particle evolves independently as a Brownian motion stopped on exiting D. Now suppose that we have constructed $\mathbf{X}_{t}^{\ell} = (X_{t}^{\ell,1}, \dots, X_{t}^{\ell,N})$ distributed as \mathcal{Q}^{ℓ} . There is exactly one particle, say $X^{\ell,j}$ such that $s_{\ell}^j = \tau_{\ell}$ and $X_{\tau_{\ell}^-}^{\ell,j} \in \partial D$. For every $i \neq j$ we have $s_{\ell+1}^i = s_{\ell}^i$. Therefore to construct $\mathcal{Q}^{\ell+1}$ from \mathcal{Q}^{ℓ} we need only to extend the lifetime of X_t^j until $s_{\ell+1}^{J}$ in accordance with the rules of Fleming-Viot. So choose $\lambda_{\ell} \in \{1, \ldots, \mathbb{N}\} \setminus \{j\}$ uniformly and independently of \mathbf{X}^{ℓ} . Set $x_{\ell} = X_{\tau_{\ell}}^{\ell,\lambda_{\ell}}$ and let \widetilde{X}_{t}^{ℓ} be a Brownian motion independent of everything else, started at x_{ℓ} and stopped on exiting D. We may construct $\mathbf{X}^{\ell+1}$ distributed as $\mathcal{Q}^{\ell+1}$ by setting $X_t^{\ell+1,i} = X_t^{\ell,i}$ whenever $t \leq s_\ell^i, i \neq j$, and $X_t^{\ell+1,j} = \widetilde{X}_{t-\tau_\ell}^{\ell}$ for $t \in [\tau_\ell, s_{\ell+1}^j)$. By an application of the Kolmogorov extension theorem there is a unique process \mathbf{X}_t with $\mathbf{X}_t = \mathbf{X}_t^{\ell}$ whenever $t < \tau_{\ell}$ and \mathbf{X}_t agrees in distribution with the construction at the start of this section.

The following lemma shows that if $\tau_{\infty} < \infty$ then all processes X_t^1, \ldots, X_t^N approach ∂D at time τ_{∞} . This result does not require any assumptions on the smoothness or regularity of ∂D , unlike our main results, so it may have independent interest.

Lemma 5.2 Let $R_t = \max_{1 \le i \le N} \operatorname{dist}(X_t^i, \partial D)$. If $D \subset \mathbb{R}^d$ is an open set and $N \ge 2$ then

$$\mathbb{P}\left(\{\tau_{\infty} = \inf\{t > 0 : R_{t-} = 0\} < \infty\} \cup \{\tau_{\infty} = \infty\}\right) = 1.$$
(5.1)

Proof Let Λ_j be the closure of the set $\{t \ge 0 : X_{t-}^j \in \partial D\}$. Suppose that $\tau_{\infty} < \infty$ with positive probability. Then at least one of the processes X_t^j must have an infinite number of jumps before τ_{∞} . For every *j* with this property we have $\tau_{\infty} \in \Lambda_j$. We will show that there are no processes X_t^j with only a finite number of jumps before τ_{∞} , a.s.

Let τ_k^j denote the time of the *k*-th jump of X_t^j for $j \in \{1, ..., N\}$. Let $\widehat{X}_t^j = X_t^j$ for $t \in [0, \tau_1^j)$. Then we define inductively $\widehat{X}_t^j = X_t^j + \widehat{X}_{\tau_k^j}^j - X_{\tau_k^j}^j$ for $t \in [\tau_k^j, \tau_{k+1}^j), k \ge 1$. It is easy to see that $\{\widehat{X}_t^j, 0 \le t \le \tau_k^j\}$ is a *d*-dimensional Brownian motion for every *k*. Hence, $\{\widehat{X}_t^j, 0 \le t < \tau_\infty\}$ is also a Brownian motion defined on a random time interval. Let m_j be the number of jumps of X_t^j before τ_∞ . Suppose that $\tau_\infty < \infty$ and $m_n = \infty$ for some *n*. Assume that $\limsup_{k\to\infty} \text{dist}(X_{\tau_k^n}^n, \partial D) > r$ for some r > 0. Then \widehat{X}_t^n has an infinite number of oscillations of size greater than or equal to *r* on every time interval of the form $(\tau_\infty - \varepsilon, \tau_\infty)$, for every $\varepsilon > 0$. This implies that \widehat{X}_t^n does not have a left limit at τ_∞ . The last event has zero probability, so we conclude that, for all rational r > 0 and all $n \in \{1, ..., N\}$,

$$\mathbb{P}\left(\{\tau_{\infty} < \infty\} \cap \{m_n = \infty\} \cap \left\{\limsup_{k \to \infty} \operatorname{dist}(X^n_{\tau^n_k}, \partial D) \ge r\right\}\right) = 0$$

Hence, for all $n \in \{1, \ldots, N\}$,

$$\mathbb{P}\left(\{\tau_{\infty} < \infty\} \cap \{m_n = \infty\} \cap \left\{\lim_{k \to \infty} \operatorname{dist}(X^n_{\tau^n_k}, \partial D) \neq 0\right\}\right) = 0.$$
(5.2)

Next suppose that $m_j < \infty$ for some $j \in \{1, ..., N\}$. By continuity of Brownian motion, the left limit $X_{\tau_{\infty}-}^j$ exists on the event $\{\tau_{\infty} < \infty, m_j < \infty\}$. We will prove that

$$\mathbb{P}\left(\{\tau_{\infty} < \infty\} \cap \{m_j < \infty\} \cap \left\{X_{\tau_{\infty}^{-}}^j \notin \partial D\right\}\right) = 0$$
(5.3)

for every j. Suppose to the contrary that for some j,

$$\mathbb{P}\left(\{\tau_{\infty} < \infty\} \cap \{m_j < \infty\} \cap \left\{X_{\tau_{\infty}^{-}}^j \notin \partial D\right\}\right) > 0.$$
(5.4)

Then there exists a rational r > 0 such that

$$\mathbb{P}\left(\{\tau_{\infty} < \infty\} \cap \{m_j < \infty\} \cap \left\{X_{\tau_{\infty}-}^j \in D_r\right\}\right) > 0,$$

where $D_r = \{x \in D : \operatorname{dist}(x, \partial D) > r\}$. Recall that each process X_t^n jumps at times τ_k^n to the location of another process $X^i, i \neq n$, chosen in a uniform way. Let *n* be such that $m_n = \infty$ (there exists at least one such *n*, a.s.). There exists an infinite subsequence $\{s_k\}_{k\geq 1}$ of $\{\tau_k^n\}_{k\geq 1}$, such that $X_{s_k}^n = X_{s_k}^j$. It follows that, on the event $\{X_{\tau_{\infty}-}^j \in D_r\}$, it is not true that $\lim_{k\to\infty} \operatorname{dist}(X_{\tau_k}^n, \partial D) = 0$. This contradicts (5.2), so we conclude that (5.4) is false. This completes the proof of (5.3). The lemma follows from (5.2) and (5.3).

Remark 5.3 The proof of Theorem 5.4, the main result of this paper, is quite complicated so we will outline the proof of a similar result in the 1-dimensional case to help the reader follow the main argument. The structure of the proof in the 1-dimensional case is the same as in the higher dimensional case but there are fewer technical details to deal with. See [26] for a more general argument based on a similar idea.

Let *Y* be one dimensional Brownian motion starting from $Y_0 = 0$ and $v_0 \ge 0$. It is well known ([19, Sect. 3.6 C]) that, a.s., there exist unique continuous processes *V* and *L* such that $V_0 = v_0$ and

$$dV_t = dY_t + dL_s, \quad \text{for } t \ge 0. \tag{5.5}$$

Here *L* is the local time of *V* at 0. In other words, *L* is a non-decreasing continuous process which does not increase when *V* is 0, i.e., $\int_0^\infty \mathbf{1}_{\{0\}}(V_t)dL_t = 0$, a.s. The process *V* is called reflected Brownian motion driven by *Y*. The construction of *V* is based on deterministic Skorokhod lemma [19, Lemma 3.6.14] so we have strong existence and uniqueness for (5.5).

Let $D = (0, \infty), N \ge 2$, let $Y^k, k = 1, ..., N$, be independent 1-dimensional Brownian motions and $x_1, ..., x_N \in D$. Let (V^k, L^k) be the solution to (5.5) driven by Y^k , with $V_0^k = x_k$, for k = 1, ..., N.

Let τ_1 be the first time when one of the processes V^k hits 0. Suppose that $V_{\tau_1}^j = 0$. We let $X_t^k = V_t^k$ for $t \in [0, \tau_1]$ and $k \neq j$, and $X_t^j = V_t^j$ for $t \in [0, \tau_1)$. We choose uniformly an integer M_1 in the set $\{1, \ldots, N\} \setminus \{j\}$ and let $X_{\tau_1}^j = X_{\tau_1}^{M_1}$. Note that $X_t^k \ge V_t^k$ for $t \in [0, \tau_1]$ and all k. We proceed by induction. Suppose that τ_1, \ldots, τ_n have been defined and $\{X_t^k, t \in [0, \tau_n]\}$ have also been defined. Assume that $X_t^k \ge V_t^k$ for $t \in [0, \tau_n]$ and all k. Let $\{(V_t^{k,n}, L_t^{k,n}), t \ge \tau_n\}$ be the solution to (5.5) driven by Y^k , with $V_{\tau_n}^{k,n} = X_{\tau_n}^k$ for all k. Note that $V_{\tau_n}^{k,n} \ge V_t^k$ for all k. Then, by the strong uniqueness of solutions to (5.5), we have $V_t^{k,n} \ge V_t^k$ for all $t \ge \tau_n$ and k. Let τ_{n+1} be the first time $t \ge \tau_n$ when one of the processes $\{V_t^{k,n}, t \ge \tau_n\}$ hits 0. Suppose that $V_{\tau_{n+1}}^m = 0$. We let $X_t^k = V_t^{k,n}$ for $t \in (\tau_n, \tau_{n+1}]$ and $k \ne m$, and $X_t^m = V_t^{m,n}$ for $t \in (\tau_n, \tau_{n+1})$. We choose uniformly an integer M_{n+1} in the set $\{1, \ldots, N\} \setminus \{m\}$ and let $X_{\tau_{n+1}}^m = X_{\tau_{n+1}}^{M_{n+1}}$. Let $\tau_{\infty} = \lim_{n \to \infty} \tau_n$. The process (X_t^1, \ldots, X_t^N) is well defined on the interval $(0, \tau_{\infty})$. The process $R_t = ((V_t^1)^2 + \cdots + (V_t^N)^2)^{1/2}$ is *N*-dimensional Bessel process and, therefore, it never hits 0. On the event $\{\tau_{\infty} < \infty\}$ we have $R_{\tau_{\infty}} > 0$ and, therefore, $\lim_{t \to \tau_{\infty}} X_t^k > 0$ for at least one *k*. In view of Lemma 5.2, we conclude that the probability of the event $\{\tau_{\infty} < \infty\}$ is 0.

The above argument is more complicated in higher dimensions because it is much harder to construct a process which always lies "closer to the boundary" of D than X^k and has a structure that can be easily analyzed. The construction of such a process uses the process Z defined in (4.4).

Theorem 5.4 There exists a constant c = c(N, d) such that if $D \subset \mathbb{R}^d$ is a bounded Lipschitz domain with the Lipschitz constant L < c(N, d), then $\tau_{\infty} = \infty$, a.s. Moreover, c(N, d) increases in N, decreases in d and

$$\lim_{N \to \infty} c(N, d) = c(d) = \frac{1}{\sqrt{d-1}}.$$
(5.6)

Proof Part 1. We start by defining c(N, d) and some other constants used in the proof. Recall the definition of $\theta_{p,d}$ and $K_{p,d}$. Let p' = 2 - 2/N and $c(N, d) = \cot \theta_{p',d}$, and fix a p such that $L < \cot \theta_{p,d} < c(N, d)$. Recall that $D_r = \{x \in D : \operatorname{dist}(x, \partial D) > r\}$. Since D is bounded and Lipschitz, there exists a small r > 0 for which the following is true. For every $x \in D \setminus D_r$ there exist an orthonormal coordinate system CS_x , $\mathcal{O}_x \in \partial D$, a Lipschitz function $F_x : \mathbb{R}^{d-1} \to \mathbb{R}$ and a cone K_x , such that \mathcal{O}_x is the origin of CS_x , K_x has vertex \mathcal{O}_x and axis passing through x, K_x can be described in CS_x as $K_{p,d}$, and

$$D \cap B(\mathcal{O}_x, r) \subset \left\{ y \text{ in } CS_x : y^d > F_x(\widetilde{y}) \right\} \cap B(\mathcal{O}_x, r),$$

$$K_{p,d} \cap B(\mathcal{O}_x, r) \subset D \cap B(\mathcal{O}_x, r).$$

As we have chosen $\cot \theta_{p,d} > L$, there exists $c_1 = c_1(p, D) > 0$ such that $\operatorname{dist}(y, \partial D) > c_1 | y - \mathcal{O}_x |$ for every $y \in K_x$ and $x \in D \setminus D_{2r}$. We set \mathcal{T}_x to be an isometry that maps $K_{p,d}$ onto K_x .

Part 2. Recall that $R_t = \max_{1 \le i \le N} \operatorname{dist}(X_t^i, \partial D)$. The process R_t is continuous on $[0, \tau_{\infty})$ because Brownian motion is continuous and R_t does not jump at any stopping time τ_k .

We will split the lifetime of the process \mathbf{X}_t into two phases, a 'safe' phase where R_t is large and an 'unsafe' phase where R_t is small. The main part of the proof will involve a special construction of \mathbf{X}_t in the unsafe phase that ensures the process cannot terminate. The two phases are defined using two sequences of stopping times V_i^r , \hat{V}_i^r , $i \ge 0$. Fix r > 0 as in Part 1 of the proof. Set $V_0^r = 0$ and for $i \ge 0$ let

$$\hat{V}_{i}^{r} = \inf\{t > V_{i}^{r} : R_{t} \le \frac{r}{2}\}, \\
V_{i+1}^{r} = \inf\{t > \hat{V}_{i}^{r} : R_{t} \ge r\},$$

with the convention that $\inf \emptyset = \infty$.

If for some *i* we have $\widehat{V}_i^r = \infty$ then $\widehat{V}_j^r = \infty$ and $V_j^r = \infty$ for all j > i. Hence, $\lim_{i\to\infty} V_i^r = \infty$. Similarly, if $V_i^r = \infty$ for some *i* then $\widehat{V}_j^r = \infty$ and $V_{j+1}^r = \infty$ for all $j \ge i$. In this case we also have $\lim_{i\to\infty} V_i^r = \infty$.

Next suppose that $\widehat{V}_i^r < \infty$ and $V_i^r < \infty$ for all *i*.

Recall Brownian motions \widehat{X}_t^j from the proof of Lemma 5.2. On any interval $[V_i^r, \widehat{V}_i^r), i \ge 1$, at least one of the processes X_t^j must travel a distance of $\frac{r}{2}$. Thus, at least one of the processes \widehat{X}_t^j must travel a distance of $\frac{r}{2}$ on this interval. Since Brownian motions \widehat{X}_t^j cannot make infinitely many such oscillations on a finite time interval and the number N of processes \widehat{X}_t^j is finite, we have $\lim_{i\to\infty} V_i^r = \infty$, a.s.

In view of (5.1), the probability that $\tau_{\infty} < \infty$ and $\tau_{\infty} \in [V_i^r, \widehat{V}_i^r]$ for some *i* is zero. So if $\tau_{\infty} < \infty$ with positive probability then there exists $i \ge 0$ such that $\mathbb{P}(\tau_{\infty} \in [\widehat{V}_i^r, V_{i+1}^r)) > 0$.

Note that $\widehat{V}_0^r = 0$ if $\mathbf{X}_0 \in (D \setminus D_{r/2})^N$. Suppose that we can show that

$$\mathbb{P}(\tau_{\infty} \in [\widehat{V}_{i}^{r}, V_{i+1}^{r})) = 0$$
(5.7)

for i = 0 and arbitrary $\mathbf{X}_0 \in (D \setminus D_{r/2})^N$. Then, by the strong Markov property applied at \widehat{V}_i^r 's, (5.7) holds for all $i \ge 0$, a.s., and this implies the theorem.

Part 3. We present an informal overview of the remaining part of the proof.

Our aim is to construct a coupling of a process \mathbf{X}_t with a vector of independent copies of the excursion process Z_t constructed in such a way that $\operatorname{dist}(X_t^j, \partial D) \ge c_2 |Z_t^j|$ for some fixed constant $c_2 = c_2(N, D)$, at least up until the stopping time V_r^1 .

The construction consists of three consecutive inductive constructions. The first inductive construction generates a coupling of a Brownian motion Y_t in the cone $K_{p,d}$ and a copy of the process Z_t . The key point is to couple Y and Z in such a way that Z is at the vertex of $K_{p,d}$ when Y hits the boundary of $K_{p,d}$. In addition the details of the construction give the bound $|Y_t| \ge |Z_t| \cos \theta_{p,d}$.

In the second level of the inductive construction we map processes Y_t from the first level of the construction into $D \setminus D_r$ using the maps \mathcal{T}_x described in Part 1. As $|Z_t| = 0$ when Y_t exits the cone we may concatenate a sequence of such processes to define a Brownian motion X_t stopped on exiting $D \setminus D_r$, coupled with an excursion process Z_t . Using estimates in Parts 1 and 4 we may bound dist $(X_t, \partial D)$ away from zero by $c_2|Z_t|$. If X_t enters D_r then a Fleming Viot process constructed using X_t has survived until V_r^1 and is safe. The coupling with Z_t is no longer needed and X_t is allowed to continue independently of Z_t until exiting D.

The second level of the construction produces a Brownian motion started anywhere in $D \setminus D_r$ coupled appropriately with a copy of Z_t . For the third and last level we follow the construction in Remark 5.1 using the products of the second level as building blocks. If a particle X_t^j exits D at time $\tau_n < V_r^1$ its sister process Z_t^i must be at the origin at time τ_i and we may extend the lifetime of the particle using the construction in Part 5. The resulting construction gives us the desired coupling of \mathbf{X}_t with Nindependent copies of Z_t until at least V_r^1 . After which the processes are allowed to decouple to give us the full process \mathbf{X}_t as required. Combining this coupling with Corollary 4.5, we may bound R_t away from zero up until time V_r^1 . Therefore the process cannot terminate before V_r^1 and we must have $\tau_{\infty} = \infty$ by observations in Part 2.

Part 4. We now present a detailed description of the first level of construction. For $y^d > 0$ we construct a coupling of the excursion process Z_t with a Brownian motion Y_t , started at $Y_0 = y = (0, 0, ..., y^d)$ and stopped on exiting the cone $K_{p,d}$. Let Y_t be a Brownian motion started at $Y_0 = y$ and consider the moving cone $C_t = Y_t - K_{p,d}$. Then $0 \in C_t$ as long as $Y_t \in K_{p,d}$ so we may let Y_t and Z_t evolve independently until the first time Z_t hits ∂C_t . At this point Z_t 'sticks' to ∂C_t and Y and Z evolve together with $dZ_t = dY_t$ until the next time Z_t exits $K_{p,d}$, after which Z_t starts again from zero evolving independently from Y_t . The process is repeated until Y_t exits $K_{p,d}$.

More formally, let W_t be a *d*-dimensional Brownian motion started at the origin and independent of Z_t and define Y_t inductively through two sequences of stopping times ξ_i, ξ'_i with $\xi_0 = 0$ and

$$\xi_{i}' = \inf\{t > \xi_{i} : Y_{t} - Z_{t} \in \partial K_{p,d}\},\\ \xi_{i+1} = \inf\{t > \xi_{i} : Z_{t} = 0\}.$$

Let $Y_0 = y = (0, 0, \dots, y^d)$ and

$$Y_t = Y_{\xi_i} + (W_t - W_{\xi_i}) \text{ for } t \in [\xi_i, \xi_i'), Y_t = Y_{\xi_i'} + (Z_t - Z_{\xi_i'}) \text{ for } t \in [\xi_i', \xi_{i+1}).$$

We stop the process Y_t at $\zeta = \inf\{t : Y_t \in \partial K_{p,d}\}$. An induction on *i* shows that Y_t is well defined and adapted to (Z_t, W_t) up until any stopping time of the form $\zeta \land \xi_i$.

By the strong Markov property of (W, Z), Y_t is a Brownian motion on all intervals $[\xi_i, \xi'_i)$ and $[\xi'_i, \xi_{i+1})$. Notice that each time ξ'_i corresponds to a separate excursion of Z_t , and as $Y_{\xi'_i} - Z_{\xi'_i} \in \partial K_{p,d}$ we must have $|Z_{\xi'_i}| \ge \operatorname{dist}(Y_{\xi'_i}, \partial K_{p,d})$. The probability that Brownian motion hits the vertex of $K_{p,d}$ is zero. For any a > 0 the excursion law H^0 is finite on the set of excursions that exit B(0, a), so as Y_t is continuous there may be only finitely many times $\xi'_i < \zeta - \varepsilon$, for any fixed $\varepsilon > 0$. All these observations imply that Y_t is well defined and a Brownian motion on $[0, \zeta]$. Furthermore, by construction, we have $Y_t \in Z_t + \overline{K_{p,d}}$ for all $t \in [0, \zeta]$ and so we must have $Z_{\zeta} = 0$. Furthermore $Z_t \in \overline{K_{p,d}}$. Hence by considering the projections of Y_t and Z_t onto the axis of the cone it is easy to see that $|Y_t| \ge |Z_t| \cos \theta_{p,d}$.

With probability 1, $\{Y_t, 0 \le t \le \zeta\}$ and ζ are unique functions of y and $\{(W_t, Z_t), 0 \le t \le \zeta\}$. We will denote these functions $Y_t = \mathcal{U}_y(W_t, Z_t)$ and $\zeta = \zeta_y(W_t, Z_t)$ respectively.

Part 5. We construct a Brownian motion X_t , stopped on exiting D by concatenating processes constructed in Part 4. To guide the construction we name some stopping times of X_t .

First set $\zeta_0 = 0$ and choose $x_0 \in D \setminus D_r$ arbitrarily. Then if X_t is a Brownian motion started at x_0 we may set

$$\tau = \inf\{t : X_{t^-} \in \partial D\}, \quad v = \inf\{t : X_{t^-} \in \partial (D \setminus D_r)\}.$$

Next recall the definitions of K_x and \mathcal{O}_x from Part 1. Define a sequence of stopping times $\zeta_n \leq v$ inductively by setting $\zeta_{n+1} = v \wedge \inf\{t > \zeta_n : X_t \in \partial K_{x_n}\}$ where $x_n = X_{\zeta_n}$. Consider the random sequence of cones K_{x_n} . At each time $\zeta_n < v$ the particle X_t is on the axis of the cone K_{x_n} and on the boundary of the cone $K_{x_{n-1}}$. Hence we have split the lifetime of the process into a sequence of Brownian motions in cones isomorphic to $K_{p,d}$. We construct the Brownian motion X_t by mapping processes Y_t^n constructed as in Part 4 into the random cones K_{x_n} to form a continuous process in $D \setminus D_r$. First we must argue that this construction partitions the entire lifetime of a Brownian motion in $D \setminus D_r$. That is the sequence of times ζ_n converges to $\tau \wedge v$, the first exit time of $D \setminus D_r$.

By definition $\zeta_n \leq v$, notice also that if $\zeta_n < \tau \wedge v$, then $X_t \in K_{x_n} \cap D \setminus D_r$ for $t \in [\zeta_n, \zeta_{n+1})$. As the probability that X_t hits the vertex \mathcal{O}_{x_n} is zero and from Part 1 dist $(X_t, \partial D) \geq c_1 | X_t - \mathcal{O}_{x_n} |$ in that interval we must have $\zeta_{n+1} < \tau$. So as τ is finite we may set $\zeta_{\infty} = \lim_{n \to \infty} \zeta_n$. By continuity of Brownian motion we have $x_n \to X_{\zeta_{\infty}}$ as $n \to \infty$. We have $|x_n - x_{n+1}| \geq |x_n - \mathcal{O}_{x_n}| \sin \theta_{p,d}$ whenever $\zeta_{n+1} < v$. As in addition $|x_n - \mathcal{O}_{x_n}| \geq \operatorname{dist}(x_n, \partial D)$, the sequence x_n cannot converge to any point in D and we have $\zeta_{\infty} = \tau \wedge v$ with probability one.

Let (W_t, Z_t) be as in Part 4. Our aim is to construct a Brownian motion X_t in such a way that X_t is adapted to (W_t, Z_t) and there exists a constant $c_2 = c_2(N, D)$ such that $dist(X_t, \partial D) \ge c_2|Z_t|$.

Suppose we have constructed X_t satisfying the above on the interval $[0, \zeta_n)$ with $\zeta_n < v$ and $Z_{\zeta_n} = 0$. This is trivial for n = 0. As X_t is adapted to (W_t, Z_t) and $Z_{\zeta_n} = 0$ we may set $W_t^n = W_{\zeta_n+t} - W_{\zeta_n}$ and $Z_t^n = Z_{\zeta_n+t}$. The pair (W_t^n, Z_t^n) agrees in distribution with (W_t, Z_t) and by the strong Markov property is independent of $X_{t \wedge \zeta_n}$.

Recall the maps \mathcal{T}_x from Part 1. As $\zeta_n < v$ we have $x_n \in D \setminus D_r$ and we may set $y_n = \mathcal{T}_{x_n}^{-1}$. Now using the construction in Part 4 set $Y_t^n = \mathcal{U}_{y_n}(W_t^n, Z_t^n), \zeta_n^\star = \zeta_{y_n}(W_t^n, Z_t^n)$ and map the process into D by setting $\widetilde{X}_t^n = \mathcal{T}_{x_n}(Y_t^n)$. If \widetilde{X}_t^n hits D_r at some time $v' < \zeta_n^\star$ during its lifetime we set $\zeta_{n+1} = v = \zeta_n + v'$, if not we set $\zeta_{n+1} = \zeta_n + \zeta_n^\star$.

By construction X_t^n is a Brownian motion started at x_n and stopped on exiting $K_{x_n} \cap D \setminus D_r$ and so dist $(X_t^n, \partial D) > c_1 |Y_t^n|$. Recall from Part 4 that $|Y_t^n| \ge |Z_t^n| \cos \theta_{p,d}$ for $t < \zeta_n^*$. So, let $c_2 = c_1 \cos \theta_{p,d}$ and set $X_t = X_{t-\zeta_n}^n$ for all $t \in [\zeta_n, \zeta_{n+1})$. Then X_t is a concatenation of two independent Brownian motions and is a Brownian motion. Furthermore we have dist $(\tilde{X}_t^n, \partial D) \ge c_2 |Z_t|$ on the interval $[\zeta_n, \zeta_{n+1})$ as required.

Now if $\zeta_{n+1} < v$ then $Z_{\zeta_{n+1}} = Z_{\zeta_n}^n = 0$ by construction and we may repeat the inductive step. If $\zeta_{n+1} = v$ we stop the construction. In this case we have constructed a Brownian motion until time v with dist $(X_t, \partial D) \ge c_2 |Z_t|$, so we need only continue X_t until time τ in such a way that it is adapted to (W_t, Z_t) . We achieve this by setting $X_t = X_v + W_t - W_v$ on the interval $[v, \tau]$.

If the construction does not terminate then X_t is a Brownian motion satisfying the above conditions until any stopping time ζ_n . Therefore we may take the limit as $n \to \infty$ and X_t is a Brownian motion up until ζ_∞ . Arguing as above we have $\zeta_\infty = \tau$.

So we have constructed a process X_t stopped on exiting D at time τ . As before, with probability 1, we may express X and τ as functions $\{X_t, t \ge 0\} = \mathcal{V}_{x_0}(\{W_t, t \ge 0\})$

 $t \ge 0$ }, $\{Z_t, t \ge 0\}$) and $\tau = \tau_{x_0}(\{W_t, t \ge 0\}, \{Z_t, t \ge 0\})$. We will use the following abbreviations, $X = \mathcal{V}_{x_0}(W, Z)$ and $\tau = \tau_{x_0}(W, Z)$.

Part 6. In the final stage we construct a coupling of a Fleming-Viot process X_t with a vector of independent excursion processes. Let Z_t^1, \ldots, Z_t^N be independent copies of the excursion process Z_t and let W_t^1, \ldots, W_t^N be independent Brownian motions. We use each pair (W_t^i, Z_t^i) as driving noise for a particle X_t^i using the constructions in Parts 4 and 5.

Recall the stopping times s_n^i and the distributions Q^n of stopped processes from Remark 5.1. Our strategy is to follow the method outlined in Remark 5.1 and construct a sequence of processes $\mathbf{X}_t^n = (X_t^{n,1}, \dots, X_t^{n,N})$ with the following properties.

- \mathbf{X}_{t}^{n} is distributed as \mathcal{Q}^{n} and adapted to \mathcal{F}_{t} .
- dist $(X_t^{n,j}, \partial D) \ge c_2 |Z_t^j|$ for every $j \in 1, \dots, \mathbb{N}$ and $t < V_r^1$.
- The sequence of processes is coherent in the sense that with probability 1 for every m > n and $t < \tau_n$ we have $\mathbf{X}_t^n = \mathbf{X}_t^m$.

It is easy to construct a process \mathbf{X}_t^1 satisfying the above as \mathcal{Q}^1 is just the distribution of N independent stopped Brownian motions. For any starting vector $\mathbf{X}_0 = (x_0^1, \ldots, x_0^N) \in (D \setminus D_r)^N$ set $X^{1,i} = \mathcal{V}_{x_0^i}(W^i, Z^i)$ for each $i \in \{1, \ldots, N\}$ and $t < s_1^i = \tau_{x_0^i}(W^i, Z^i)$. As the pairs (W_t^i, Z_t^i) are independent and distributed as (W_t, Z_t) in Part 5 the Brownian motions $X_t^{1,i}$ are independent and satisfy the required bound on dist $(X_t^i, \partial D)$.

Now suppose after ℓ inductive steps we have constructed processes $\mathbf{X}_{t}^{1}, \ldots, \mathbf{X}_{t}^{\ell}$ satisfying the three conditions above. Suppose further that for each $j \in \{0, \ldots, N\}$ the shifted processes $Z_{t}^{j,\ell} = Z_{t+s_{\ell}^{j}}^{j}$ and $W_{t}^{j,\ell} = W_{t+s_{\ell}^{j}}^{j} - W_{s_{\ell}^{j}}^{j}$ are independent of \mathbf{X}^{ℓ} . This fact is easy to check for $\ell = 1$.

Recall from Remark 5.1 that there is exactly one particle, say $X_t^{\ell,j}$ such that $s_\ell^j = \tau_\ell$. To construct $\mathbf{X}^{\ell+1}$ we must extend the lifetime of $X_t^{\ell,j}$ by adding an independent Brownian motion with an appropriately chosen starting position.

So as in Remark 5.1, set $X_t^{\ell+1,i} = X_t^{\ell,i}$ whenever $t < s_\ell^i$. Next choose, $\lambda_{\ell+1}$ uniformly from $\{1, \ldots, N\} \setminus \{j\}$ and independent of every other random variable. Then the particle $X^{\ell+1,j}$ will jump to $x_\ell = X_{\tau_\ell}^{\ell,\lambda_\ell}$ at time τ_ℓ .

If $\tau_{\ell} < V_r^1$ then we must have $Z_{\tau_{\ell}}^j = 0$ and $x_{\ell} \in D \setminus D_r$ so we may set $s_{\ell+1}^j = s_{\ell}^j + \tau_{x_{\ell}}(W^{\ell,j}, Z^{\ell,j})$ and $X^{\ell+1,j} = \mathcal{V}_{x_{\ell}}(W^{\ell,j}, Z^{\ell,j})$ for $t \in [s_{\ell}^j, s_{\ell+1}^j)$. By construction we have dist $(X_t^{\ell+1,j}, \partial D) \ge c_2 |Z_t^j|$ on the interval $[\tau_j, s_{\ell+1}^j \land v)$.

Alternatively if $\tau_{\ell} > V_r^1$ we do not need to couple $X_t^{\ell+1,j}$ with Z_t^j after time τ_{ℓ} so we may set $X_t^{\ell+1,j} = x_{\ell} + W_t^j - W_{s_{\ell}^j}^j$ until $s_{\ell+1}^j$, the next time $X_t^{\ell+1,j}$ exits D.

By assumption the pair (W_t^{ℓ}, Z_t^{ℓ}) is independent of \mathbf{X}_t^{ℓ} . In both cases above $X_{t-s_\ell^j}^{\ell+1,j}$ is a Brownian motion adapted to (W_t^{ℓ}, Z_t^{ℓ}) . Arguing as in Remark 5.1, $\mathbf{X}_t^{\ell+1}$ is distributed as $\mathcal{Q}^{\ell+1}$ and is adapted to \mathcal{F}_t . The independence assumption for $\mathbf{X}^{\ell+1}$ follows from the strong Markov property.

Now extend the sequence \mathbf{X}_t^n to a Fleming-Viot process by setting $\mathbf{X}_t = \mathbf{X}_t^n$ whenever $t < \tau_n$.

Suppose there exists a finite deterministic time T with $\mathbb{P}(1/T < \tau_{\infty} < T) > 0$. Let $A = \{1/T < \tau_{\infty} < T\}$. Then from Lemma 5.2 we must have $R_{\tau_n} \to 0$ as $n \to \infty$, on A. But from Corollary 4.5 the maximum process $\max_{i \in \{1, \dots, N\}} |Z_i^i|$ is bounded below on [1/T, T] by a strictly positive random variable, a.s. Therefore we must have $\operatorname{dist}(X_{\tau_n}^i, \partial D) < c_2 |Z_i^i|$ for all but finitely many n, for all i, assuming A holds. According to our construction, this is impossible unless $\tau_{\infty} > V_r^1$ on A. Letting $T \to \infty$, we obtain $\mathbb{P}(\tau_{\infty} > V_r^1) = 1$. Part 2 now implies that we must have $\mathbb{P}(\tau_{\infty} < \infty) = 0$.

Remark 5.5 Let $D \subset \mathbb{R}^d$ be a bounded Lipschitz domain with the Lipschitz constant $L < \frac{1}{\sqrt{d-1}}$. Then by (5.6) we see that there exists N_0 so large that L < c(N, d) for all $N \ge N_0$. In consequence, the Fleming-Viot-type particle process \mathbf{X}_t in D is well defined for all $t \ge 0$ provided it consists of N particles with $N \ge N_0$.

Remark 5.6 We will argue that Theorems 1.3 and 1.4 in [10] hold true even though we do not know whether Theorem 1.1 in that paper is true.

Theorem 1.3 is concerned with a fixed time t > 0. A sequence of processes \mathbf{X}^N is considered and it is assumed that the initial empirical distributions $(1/N) \sum_{k=1}^{N} \delta_{X_{k}^{k}}$ converge weakly to a probability measure μ_0 in D. Let A be a compact subset of D such that $\mu_0(A) > 0$. Since the distance r from A to ∂D is strictly positive and Brownian motion can stay in the ball of radius r/2 for time t with a strictly positive probability p_1 , it follows that if *j* Brownian motions start from points in A then with probability equal to or greater than $1 - (1 - p_1)^j$ at least one of these Brownian motions never comes closer to ∂D than r/2 units on the interval [0, t]. Fix an arbitrarily small p_2 and let j be so large that $(1 - p_1)^j < p_2/2$. Let N_0 be so large that for $N \ge N_0$, the probability that there are at least j processes X_0^k in A is equal to or greater than $1-p_2/2$. Then with probability equal to or greater than $1-p_2$ there exists a process X^k which never approaches ∂D closer than r/2 units on the interval [0, t]. By Lemma 5.2, this implies that $\tau_{\infty} \ge t$ with probability equal to or greater than $1 - p_2$. Hence, the empirical distribution $(1/N) \sum_{k=1}^{N} \delta_{X_{k}^{k}}$ is well defined with probability equal to or greater than $1 - p_2$. When the empirical distribution of \mathbf{X}_t is not well defined at time *t*, we can define it arbitrarily to be the atom at (x_0, x_0, \ldots, x_0) for some $x_0 \in D$. Theorem 1.3 of [10] makes an assertion about convergence of $(1/N) \sum_{k=1}^{N} \delta_{X_{k}^{k}}$ in probability. Since $p_2 > 0$ is arbitrarily small, the proof given in [10] and the remarks given above show that Theorem 1.3 is true.

We will now discuss Theorems 1.4 in [10]. That theorem is concerned with domains which satisfy the internal ball condition with radius r > 0. The family of such domains is not contained in the class of Lipschitz domains and neither does it contain all Lipschitz domains. To see this, consider a square which is a Lipschitz domain but does not satisfy the internal ball condition. A two-dimensional example illustrating the opposite claim is $D = B((10, 0), 10) \cup B((0, 1), 1) \cup B((0, -1), 1)$. This example shows that Theorem 5.4 of the present paper cannot be applied to some domains satisfying the internal ball condition. On page 698 of [10] it is shown that $dist(X_t^k, D^c) \ge r - R_t^k$, where R^k 's are independent *d*-dimensional Bessel processes reflected at *r*. A claim is made in [10] that this relation holds for all finite *t*, based on Theorem 1.1. Although Theorem 1.1 has incorrect proof, the argument given in the proof of Theorem 1.4 does show that $\operatorname{dist}(X_t^k, D^c) \ge r - R_t^k$ holds for all $t < \tau_{\infty}$. The process $\Gamma_t = ((r - R_t^1)^2 + \dots + (r - R_t^N)^2)^{1/2}$ has the distribution absolutely continuous with respect to the distribution of *N*-dimensional Bessel process on every finite time interval, by the Girsanov theorem. Hence Γ_t does not hit 0 at any finite time. We now reason as in Remark 5.3. On the event $\{\tau_{\infty} < \infty\}$ we have $\Gamma_{\tau_{\infty}} > 0$ and, therefore, $\limsup_{t \uparrow \tau_{\infty}} \operatorname{dist}(X_t^k, D^c) > 0$ for at least one *k*. In view of Lemma 5.2, we conclude that the probability of the event $\{\tau_{\infty} < \infty\}$ is 0. This shows that the process \mathbf{X}_t is well defined for all *t* under assumptions of Theorem 1.4 in [10] and, therefore, Theorem 1.4 is true.

Example 5.7 The proof of Theorem 1.1 in [10] contains an error. Formula (2.1) in [10] does not follow "by induction" from the previous statement. We will show that the error is irreparable in the following sense. The proof of Theorem 1.1 in [10] is based only on two properties of Brownian motion—the strong Markov property and the fact the the hitting time distribution of a compact set has no atoms (assuming that the starting point lies outside the set). Hence, if some version of that argument were true, it would apply to almost all non-trivial examples of Markov processes with continuous time, and in particular to all diffusions. However we may find a diffusion for which the analogue of Theorem 1.1 in [10] is false. Let X_t be the diffusion on $[0, \infty)$, started at $X_0 = 1$ and satisfying the SDE

$$dX_t = dW_t - \frac{5}{2X_t}dt.$$

We make 0 absorbing so that it can play the role of the boundary for the domain $D = (0, \infty)$. Notice that although X_t is not a Bessel process, as we have reversed the drift term, it scales in the same way. That is, for $\alpha > 0$, $\alpha X_{t\alpha^{-2}}$ is a diffusion satisfying the same SDE, but started at α . Let $\mathbf{Y}_t^i = (Y_t^{i,1}, Y_t^{i,2})$, $i = 1 \dots \infty$, be a double sequence of independent copies of X_t , and set

$$\sigma_{i} = \inf\{t > 0 : Y_{t}^{i,1} \land Y_{t}^{i,2} = 0\},\ \alpha_{i} = Y_{\sigma_{i}}^{i,1} \lor Y_{\sigma_{i}}^{i,2}.$$

Now, construct a two-particle Fleming-Viot type process $\mathbf{X}_t = (X_t^1, X_t^2)$ as follows. First let $\tau_1 = \sigma_1$ and set $\mathbf{X}_t = \mathbf{Y}_t^1$ for $t \in [0, \tau_1)$. At τ_1 one of the particles hits the boundary and jumps to $\xi_1 = \alpha_1$. To continue the process we use the scaling property of \mathbf{Y}_t and set $\mathbf{X}_t = \xi_1 \mathbf{Y}_{(t-\tau_1)\xi_1^{-2}}^2$ for $t \in [\tau_1, \tau_2)$ where $\tau_2 = \tau_1 + \xi_1^2 \sigma_2$. At τ_2 a second particle hits the boundary and jumps, this time to $\xi_2 = \alpha_2 \xi_1$, and we continue the process in the same way by setting

$$\xi_{i} = \prod_{j=1}^{i} \alpha_{j}, \quad \tau_{i} = \sum_{j=1}^{i} \xi_{j-1}^{2} \sigma_{j},$$
$$\mathbf{X}_{t} = \xi_{i} \mathbf{Y}_{(t-\tau_{i})\xi_{i}^{-2}}^{i}, \quad \text{for } t \in [\tau_{i}, \tau_{i+1}).$$

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Then \mathbf{X}_t evolves as two independent copies of X_t with Fleming-Viot type jumps when a particle hits the boundary. The process \mathbf{X}_t is well defined up until τ_{∞} and if the analogue of [10, Theorem 1.1] were to hold for this process we would have $\tau_{\infty} = \infty$ almost surely. In fact the opposite is true. We will show now that $\mathbb{E}\tau_{\infty} < \infty$ and hence $\tau_{\infty} < \infty$ almost surely. To do this it will be sufficient to show $\mathbb{E}(\alpha_1^{-2}) < 1$ and $\mathbb{E}\sigma_1 < \infty$. Let $f(x, y) = x^4 + y^4 - x^2y^2$ and notice $f(x, x) = f(x, 0) = f(0, x) = x^4$. We may check using Ito's formula that $f(Y_{t,\Lambda\sigma_i}^{i,1}, Y_{t,\Lambda\sigma_i}^{i,2})$ is a positive local martingale and hence a supermartingale. By the optional stopping theorem

$$\mathbb{E}\left(\alpha_{1}^{4}\right) = \mathbb{E}f\left(Y_{\sigma_{1}}^{1,1}, Y_{\sigma_{1}}^{1,2}\right) \leq \mathbb{E}f\left(Y_{0}^{1,1}, Y_{0}^{1,2}\right) = 1.$$

Furthermore, α_1 is not almost surely constant and so by Jensen's inequality

$$\mathbb{E}(\alpha_1^2) < \sqrt{\mathbb{E}(\alpha_1^4)} = 1.$$

We may use Ito's formula again to show that $X_t^2 + 4t$ is a local martingale and so by the optional stopping theorem again we have that $\mathbb{E}(\sigma_1) \leq \frac{1}{4}$.

By independence of the **Y**^{*i*} processes we have that $\mathbb{E}(\xi_i^2) = \mathbb{E}(\alpha_1^2)^i$ and so

$$\mathbb{E}\tau_{\infty} = \sum_{j=1}^{\infty} \mathbb{E}\left(\xi_{j-1}^{2}\sigma_{j}\right) \leq \frac{1}{4} \sum_{j=0}^{\infty} \mathbb{E}\left(\alpha_{1}^{2}\right)^{j} < \infty.$$

6 Hitting probabilities of compact sets

This section is devoted to a technical estimate needed in the proof of Theorem 7.1. Recall definitions of D_r and $\mathbf{X}_t = (X_t^1, \dots, X_t^N)$.

Lemma 6.1 Fix $N \ge 2$ and let $D \subset \mathbb{R}^d$ be a bounded Lipschitz domain with the Lipschitz constant L < c(N, d).

(i) For any fixed $k \in \{1, ..., N\}$, and for every r > 0 such that $\operatorname{Int} D_r \neq \emptyset$, there exist c > 0 and t > 0 such that for all $\mathbf{x} \in D^N$,

$$\mathbb{P}^{\mathbf{x}}\left(X_t^k \in D_r\right) \ge c.$$

(ii) For every r > 0 such that $\operatorname{Int} D_r \neq \emptyset$, there exist c > 0 and t > 0 such that for all $\mathbf{x} \in D^N$,

$$\mathbb{P}^{\mathbf{X}}\left(\mathbf{X}_{t}\in D_{r}^{N}\right)\geq c.$$

Proof (i) Fix r > 0 such that $Int D_r \neq \emptyset$. Recall that notation such as T_{D_r} , $T_{\partial D}$, etc. refers to hitting times by Brownian motion. By Theorem 3.1 there exists $c_0 = c_0(r)$ such that for all $x \in D$,

$$\mathbb{P}^{x}\left(T_{D_{r}} < T_{\partial D}\right) \ge c_{0}\mathbb{E}^{x}T_{\partial D}.$$
(6.1)

Fix k and let $T_{D_r}^{X^k} = \inf \{t \ge 0 : X_t^k \in D_r\}$, and

$$Y_t = X^k (t \wedge T_{D_r}^{X^k}).$$

Define $T_0 = 0$ and

$$T_{n+1} = \inf \left\{ t > T_n : \lim_{s \to t^-} Y_s \in \partial D \right\} \wedge T_{D_r}^{X^k}.$$

Let $M_0 = 0$ and

$$M_n = \frac{1}{c_0} \mathbf{1}_{\{Y(T_n) \in D_r\}} - T_n, \quad n \ge 1,$$

and

$$\mathcal{F}_n = \sigma(\mathbf{X}_t, t \leq T_n).$$

It is easy to see that $ET_n < \infty$ so $E|M_n| < \infty$. For $\mathbf{x} = (x_1, x_2, \dots, x_N) \in D^N$ with $x_k \notin D_r$,

$$\mathbb{E}^{\mathbf{x}} (M_{n+1} - M_n \mid \mathcal{F}_n) = \frac{1}{c_0} \mathbb{E}^{\mathbf{x}} \left(\mathbf{1}_{\{Y(T_{n+1}) \in D_r\}} (\mathbf{1}_{\{Y(T_n) \notin D_r\}} + \mathbf{1}_{\{Y(T_n) \in D_r\}}) - \mathbf{1}_{\{Y(T_n) \in D_r\}} \mid \mathcal{F}_n \right) - \mathbb{E}^{\mathbf{x}} (T_{n+1} - T_n \mid \mathcal{F}_n)$$

$$= \frac{1}{c_0} \mathbb{E}^{\mathbf{x}} \left(\mathbf{1}_{\{Y(T_{n+1}) \in D_r\}} \mathbf{1}_{\{Y(T_n) \notin D_r\}} + \mathbf{1}_{\{Y(T_n) \in D_r\}} - \mathbf{1}_{\{Y(T_n) \in D_r\}} \mid \mathcal{F}_n \right) - \mathbb{E}^{\mathbf{x}} (T_{n+1} - T_n \mid \mathcal{F}_n)$$

$$= \frac{1}{c_0} \mathbf{1}_{\{Y(T_n) \notin D_r\}} \mathbb{P}^{\mathbf{x}} (Y(T_{n+1}) \in D_r \mid \mathcal{F}_n)$$

$$-\mathbb{E}^{\mathbf{x}} (T_{n+1} - T_n \mid \mathcal{F}_n) .$$

We have on the event $\{Y(T_n) \notin D_r\}$,

$$\mathbb{E}^{\mathbf{x}}\left(M_{n+1}-M_n\mid \mathcal{F}_n\right) \geq \frac{1}{c_0} \mathbb{P}^{X^k(T_n)}\left(T\left(D_r\right) < T_{\partial D}\right) - \mathbb{E}^{X^k(T_n)}T_{\partial D} \geq 0,$$

by (6.1). On the event $\{Y(T_n) \in D_r\}$, we have $T_{n+1} = T_n, Y_{T_{n+1}} \in D_r$, and so

$$\mathbb{E}^{\mathbf{x}}(M_{n+1}-M_n\mid \mathcal{F}_n)=0.$$

Combining the last two formulas, we conclude that $\{M_n\}$ is a submartingale with respect to $\{\mathcal{F}_n\}$.

Define

$$S = \inf \left\{ j : T_j \ge 1 \right\} \wedge \inf \left\{ j : Y_{T_j} \in D_r \right\}.$$

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Fix an $\mathbf{x} \in D^N$ and consider two cases. First, we may have

$$\mathbb{P}^{\mathbf{X}}\left(S = \inf\left\{j : Y_{T_j} \in D_r\right\}\right) \ge 1/2.$$

In this case,

$$\mathbb{P}^{\mathbf{X}}(T_{D_r}^{X^k} \le 1) \ge 1/2.$$
(6.2)

The second case is when

$$\mathbb{P}^{\mathbf{X}}\left(S=\inf\left\{j:Y_{T_{j}}\in D_{r}\right\}\right)<1/2.$$

In this case, $\mathbb{P}^{\mathbf{x}}(S \ge 1) \ge 1/2$, so $\mathbb{E}^{\mathbf{x}}T_S \ge 1/2$. The submartingale M_n is bounded above by $1/c_0$ so we can apply the optional stopping theorem to obtain

$$\mathbb{E}^{\mathbf{X}}M_S \geq \mathbb{E}^{\mathbf{X}}M_0 = 0.$$

Hence

$$\mathbb{P}^{\mathbf{x}}\left(Y_{T_{S}} \in D_{r}\right) \ge c_{0}\mathbb{E}^{\mathbf{x}}T_{S} \ge c_{0}/2.$$

$$(6.3)$$

We will show that for some t_0 ,

$$\mathbb{P}^{\mathbf{X}}\left(T_{D_r}^{\mathbf{X}^k} \le t_0\right) \ge c_0/4.$$
(6.4)

If $T_S > s_0$ for some $s_0 > 1$ then X_t^k must not hit $D_r \cup \partial D$ for $t \in (1, s_0)$. The probability of this event is bounded above by the probability of the event that Brownian motion starting from X_1^k will not leave the ball $B(X_1^k, 2 \operatorname{diam}(D))$ for $s_0 - 1$ units of time. The last probability is $c_1 < 1$, depending on $s_0 > 1$, but not depending on X_1^k . By the Markov property,

$$\sup_{\mathbf{x}\in D^N} \mathbb{P}^{\mathbf{x}} \left(T_S > s_0 \right) \le c_1 < 1.$$

Applying the Markov property repeatedly at times $s_0, 2s_0, \ldots$, we obtain for any $\mathbf{x} \in D^N$,

$$\mathbb{P}^{\mathbf{X}}(T_S > ns_0) \leq c_1^n$$

We choose *n* so large that $c_1^n \le c_0/4$ and let $t_0 = ns_0$. Then for $\mathbf{x} \in D^N$,

$$\mathbb{P}^{\mathbf{x}}(T_S > t_0) \le c_0/4. \tag{6.5}$$

We use (6.3) and (6.5) to see that

$$c_0/2 \leq \mathbb{P}^{\mathbf{x}} \left(Y_{T_S} \in D_r \right)$$

= $\mathbb{P}^{\mathbf{x}} \left(Y_{T_S} \in D_r, T_S > t_0 \right) + \mathbb{P}^{\mathbf{x}} \left(Y_{T_S} \in D_r, T_S \leq t_0 \right)$
 $\leq \mathbb{P}^{\mathbf{x}} (T_S > t_0) + \mathbb{P}^{\mathbf{x}} \left(T_{D_r}^{X^k} \leq t_0 \right)$
 $\leq c_0/4 + \mathbb{P}^{\mathbf{x}} \left(T_{D_r}^{X^k} \leq t_0 \right).$

This implies (6.4). We combine the two cases, that is, (6.2) and (6.4), to see that for some $t_1 < \infty$ and c_2 , for all $\mathbf{x} \in D^N$,

$$\mathbb{P}^{\mathbf{x}}\left(T_{D_r}^{X^k} \le t_1\right) \ge c_2. \tag{6.6}$$

Let r_1 be such that $0 < r < r_1$ and $Int D_{r_1} \neq \emptyset$. Let t_2 and c_3 be such that (6.6) holds with r_1, t_2 and c_3 in place of r, t_1 and c_2 , i.e.,

$$\mathbb{P}^{\mathbf{x}}\left(T_{D_{r_1}}^{X^k} \le t_2\right) \ge c_3. \tag{6.7}$$

Let $r_2 = (r_1 - r)/2$ and $p_1 = \mathbb{P}^0(T_{\partial B(0,r_2)} \ge t_2) > 0$. By translation invariance of Brownian motion, $p_1 = \mathbb{P}^y(T_{\partial B(y,r_2)} \ge t_2)$ for every y. If the process X^k hits D_{r_1} before time t_2 and then stays in the ball $B(X^k(T_{D_{r_1}}^{X^k}), r_2)$ for at least t_2 units of time then X^k will be inside D_r at time t_2 . By the strong Markov property applied at the stopping time $T_{D_{r_1}}^{X^k}$, we obtain, using (6.7), for all $\mathbf{x} \in D^N$,

$$\mathbb{P}^{\mathbf{x}}(X_{t_2}^k \in D_r) \ge p_1 \mathbb{P}^{\mathbf{x}} \left(T_{D_{r_1}}^{X^k} \le t_2 \right) \ge p_1 c_3 > 0.$$
(6.8)

This completes the proof of part (i) of the lemma.

(ii) Recall that r > 0 is fixed and such that $\operatorname{Int} D_r \neq \emptyset$. Let r_3 and r_4 be such that $0 < r < r_3 < r_4$ and $\operatorname{Int} D_{r_4} \neq \emptyset$. Let $r_5 = \min(r_3 - r, r_4 - r_3)/2$. We choose t_3 and c_4 so that (6.8) can be applied with r_4 in place of r,

$$\mathbb{P}^{\mathbf{X}}(X_{t_3}^k \in D_{r_4}) \ge c_4 > 0.$$

Let $p_2 = \inf_{y \in D} \mathbb{P}^y(T_{\partial D} \le t_3)$ and note that $p_2 > 0$. Let $p_3 = \mathbb{P}^y(T_{\partial B(y,r_5)} \ge 2t_3) > 0$ and note that p_3 does not depend on y.

Let *A* be the intersection of the following events.

- (a) The process X^1 is in D_{r_4} at time t_3 , and it stays in $B(X^1_{t_3}, r_5)$ for all $t \in [t_3, 3t_3]$.
- (b) For every j = 2, ..., N, the process X^j jumps at a time $s_j \in [t_3, 2t_3]$ to $X^1_{s_j}$, and then stays in the ball $B(X^j_{s_j}, r_5) = B(X^1_{s_j}, r_5)$ for all $t \in [s_j, s_j + 2t_3]$.

By the strong Markov property and the definition of the process **X**, the probability of *A* is bounded below by $c_5 = c_4 p_3 (p_2(1/(N-1))p_3)^{N-1}$. If *A* occurs then $\mathbf{X}_{3t_3} \in D_r^N$.

$$\mathbb{P}^{\mathbf{X}}\left(\mathbf{X}_{3t_{3}}\in D_{r}^{N}\right)\geq c_{5}>0.$$

This proves part (ii) of the lemma.

Hence, for every $\mathbf{x} \in D^N$,

7 Stationary distribution for the particle system

The two theorems proved in this section generalize the analogous results in [10], where the proofs were given only for domains satisfying the internal ball condition.

Theorem 7.1 Suppose that $D \subset \mathbb{R}^d$ is a bounded Lipschitz domain with the Lipschitz constant L < c(N, d), where c(N, d) is as in Theorem 5.4. Then there exists a unique stationary probability distribution \mathcal{M}^N for \mathbf{X}_t . The process \mathbf{X}_t converges to its stationary distribution exponentially fast, i.e., there exists $\lambda > 0$ such that for every $A \subset D^N$,

$$\lim_{t \to \infty} e^{\lambda t} \sup_{\mathbf{x} \in D^N} \left| \mathbb{P}^{\mathbf{x}} \left(\mathbf{X}_t \in A \right) - \mathcal{M}^N(A) \right| = 0.$$
(7.1)

Proof We have shown in Lemma 6.1 (ii) that for any r > 0, with probability higher than $p_0 = p_0(r) > 0$, the process \mathbf{X}_t can reach the compact set D_r^N within $t_0 > 0$ units of time. This and the strong Markov property applied at times $2t_0, 4t_0, 6t_0, \ldots$ show that the hitting time of D_r^N is stochastically bounded by an exponential random variable with the expectation independent of the starting point of \mathbf{X}_t . Since the transition densities $p_t^{\mathbf{X}}(\mathbf{x}, \mathbf{y})$ for \mathbf{X}_t are bounded below by the densities for the Brownian motion killed at the exit time from D^N , we see that $p_t^{\mathbf{X}}(\mathbf{x}, \mathbf{y}) > c_1 > 0$ for $\mathbf{x}, \mathbf{y} \in D_r^N$. Fix arbitrarily small s > 0 and consider the "skeleton" $\{\mathbf{X}_{ns}\}_{n\geq 0}$. The properties listed in this paragraph imply that the skeleton has a stationary probability distribution and that it converges to that distribution exponentially fast, i.e., (7.1) holds for the skeleton, by Theorem 2.1 in [15] or Theorem 16.0.2 (ii) and (vi) of [21]. See the proof of Proposition 1.2 in [9] for an argument showing how to pass from the the statement of uniform ergodicity for the skeleton to the analogous statement for the continuous process $t \to \mathbf{X}_t$. We sketch this argument here. Take any $\varepsilon > 0$ and find $t_1 = n_1 s$ such that

$$\left| e^{\lambda t} \sup_{\mathbf{x} \in D^{N}} \left| \mathbb{P}^{\mathbf{x}} \left(\mathbf{X}_{t} \in A \right) - \mathcal{M}^{N}(A) \right| \leq \varepsilon$$
(7.2)

holds for $t \ge t_1$ of the form t = ns. Consider an arbitrary $t_2 > t_1$, not necessarily of the form *ns*. Let *m* be the integer part of t_2/s and let $u = t_2 - ms$. Note that $m \ge n_1$. Since (7.2) holds for t = ms, the semigroup property applied at time *u* shows that (7.2) holds also at time t_2 .

Theorem 7.2 Suppose that D is a bounded Lipschitz domain with the Lipschitz constant $L < \frac{1}{\sqrt{d-1}}$. For $N \ge N_0$ (see Remark 5.5) let $\mathcal{X}_{\mathcal{M}}^N$ be the stationary empirical measure. Let φ be the first eigenfunction for Laplacian in D with the Dirichlet boundary conditions, normalized so that $\int_D \varphi = 1$. Then the sequence of random measures $\mathcal{X}_{\mathcal{M}}^N$, $N \ge N_0$, converges as $N \to \infty$ to the (non-random) measure with the density φ , in the sense of weak convergence of random measures.

Proof Recall processes Y^j defined in the proof of Theorem 5.4. By construction, we have dist $(Y_t^j, \partial D) \leq \text{dist}(X_t^j, \partial D)$, for all *j* and *t*.

It is elementary to see that the process Z constructed in Sect. 4 has the property that

$$\lim_{r \downarrow 0} \limsup_{t \to \infty} \frac{1}{t} \int_{0}^{t} \mathbf{1}_{\{\operatorname{dist}(Z_{s}, \partial D) \leq r\}} ds = 0, \text{ a.s.}$$

In view of the construction of Y^{j} from independent copies of Z, we also have, for every j,

$$\lim_{r \downarrow 0} \limsup_{t \to \infty} \frac{1}{t} \int_{0}^{t} \mathbf{1}_{\{\operatorname{dist}(Y_{s}^{j}, \partial D) \leq r\}} ds = 0, \text{ a.s.}$$

Hence, for every j,

$$\lim_{r \downarrow 0} \limsup_{t \to \infty} \frac{1}{t} \int_{0}^{t} \mathbf{1}_{\{\operatorname{dist}(X_{s}^{j}, \partial D) \le r\}} ds = 0, \text{ a.s.}$$

This implies that for every $p_1 > 0$, one can find r > 0 so small that if **X** has the stationary measure \mathcal{M}^N then for every t, $\mathbb{P}(X_t^j \notin D_r) \leq p_1$. It follows that for any N, the mean measure $\mathcal{EX}_{\mathcal{M}}^N$ of the compact set D_r is not less than $1 - p_1$. Hence, the mean measures $\mathcal{EX}_{\mathcal{M}}^N$ are tight in D. Lemma 3.2.7, p. 32, of [13] implies that the sequence of random measures $\mathcal{X}_{\mathcal{M}}^N$ is tight and so it contains a convergent subsequence.

One can complete the proof of the claim that the random measures $\mathcal{X}_{\mathcal{M}}^N$ converge as $N \to \infty$ to the measure with the density φ exactly as in the proof of Theorem 1.4 in [10], starting on line 9 of page 699.

8 Polyhedral domains

In this section we show that the Lipschitz constant c(N, d) in Theorem 5.4 is not sharp, that is, $\tau_{\infty} = \infty$, a.s., in some Lipschtz domains with arbitrarily large Lipschitz constant. Specifically, we will demonstrate the existence of the two particle process for all times in arbitrary polyhedral domains. Unfortunately, our method cannot be easily adapted to the multiparticle case, so we leave this generalization as an open problem.

Definition 8.1 We say an open set $D \subset \mathbb{R}^d$ is a *polyhedral domain* if there exist simplicial complexes $\mathcal{K} \supset \partial \mathcal{K}$ such that $\overline{D} = |\mathcal{K}|$ and $\partial D = |\partial \mathcal{K}|$.

For the remainder of this section we will assume that $D = \text{Int}|\mathcal{K}|$ is a polyhedral domain. Let $\mathbf{X}_t = (X_t^1, X_t^2)$ be a Fleming-Viot process in D and define jump times τ_i as before. We will show:

Theorem 8.2 If *D* is a polyhedral domain and $\mathbf{X}_t = (X_t^1, X_t^2)$ is a Fleming-Viot process with jump times τ_i then $\tau_i \to \infty$ as $i \to \infty$ almost surely.

As \mathbf{X}_t is a CADLAG process we have $X_{\tau_i}^1 = X_{\tau_i}^2$ for each $i \in \mathbb{N}$, so we may define a sequence of *jump points* $\xi_i = X_{\tau_i}^1 = X_{\tau_i}^2$. Since \overline{D} is compact, ξ_i has at least one limit point in \overline{D} . To prove Theorem 8.2 we will examine the behavior of \mathbf{X}_t when both particles are close to a limit point of ξ_i and, assuming that $\tau_{\infty} < \infty$, arrive at a contradiction.

First we will show that if $t \in [\tau_i, \tau_{i+1})$ then \mathbf{X}_t cannot stray too far from (ξ_i, ξ_i) .

Lemma 8.3 Set $V_t^1 = ||X_t^1 - \xi_i||, V_t^2 = ||X_t^2 - \xi_i||$ for $t \in [\tau_i, \tau_{i+1})$. If $\tau_{\infty} < \infty$ then $V_t^1 \to 0$ and $V_t^2 \to 0$ as $t \to \tau_{\infty}$.

Proof It suffices to consider only V_t^1 . Notice that V_t^1 is a *d*-dimensional Bessel process (Bes(*d*), for short), reset to 0 at each τ_i . So setting $\Delta V_i^1 = V_{\tau_i^-}^1$ we may extract a Brownian motion

$$W_t = V_t^1 + \sum_{\{i \in \mathbb{N} : \tau_i \le t\}} \Delta V_i^1 - \int_0^t \frac{d-1}{2V_t^1} dt.$$

Consider $\varepsilon > 0$. We will count the number of upcrossings of the interval $[\frac{\varepsilon}{2}, \varepsilon]$ within a short time interval $[t, t + \delta]$, where $\delta = \varepsilon^2/(4(d - 1))$. Consider times $t < s' < s < t + \delta$ where $V_s^1 \ge \varepsilon$ and $s' = \sup\{\tilde{s} < s : V_{\tilde{s}}^1 = \frac{\varepsilon}{2}\}$. Notice as V^1 only jumps downwards there is no $i \in \mathbb{N}$ such that $s' < \tau_i \le s$. We have

$$W_{s} - W_{s'} = V_{s}^{1} - V_{s'}^{1} - \int_{s'}^{s} \frac{d-1}{2V_{t}^{1}} dt$$
$$\geq \frac{\varepsilon}{2} - (s-s')\frac{d-1}{\varepsilon}$$
$$\geq \frac{\varepsilon}{2} - \frac{\varepsilon^{2}}{4(d-1)}\frac{d-1}{\varepsilon} = \frac{\varepsilon}{4}$$

So on a short time interval, each upcrossing of $[\frac{\varepsilon}{2}, \varepsilon]$ by V^1 corresponds to an oscillation of $\frac{\varepsilon}{4}$ by W. As W is a Brownian motion, with probability 1, V^1 makes only finitely many upcrossings of $[\frac{\varepsilon}{2}, \varepsilon]$ in a given time interval $[t, t + \delta]$. If $\tau_{\infty} < \infty$, we may find $n \in \mathbb{N}$ with $\tau_n \ge \tau_{\infty} - \delta$. So if $V_t^1 > \varepsilon$ for some $\tau_n < \tau_i < t < \tau_{i+1}$ then as V^1 is reset to 0 at τ_i there must be an upcrossing of $[\frac{\varepsilon}{2}, \varepsilon]$ in the interval $[\tau_i, \tau_{i+1}) \subset [\tau_n, \tau_n + \delta]$.

So $V_t^1 > \varepsilon$ in only finitely many intervals $[\tau_i, \tau_{i+1})$ and, as ε is arbitrary, $V_t^1 \to 0$ as $t \to \tau_{\infty}$.

Corollary 8.4 If $\tau_{\infty} < \infty$ then the sequence ξ_i has no limit point $\xi_{\infty} \in D$.

Proof Fix $x \in D$. As D is open, there exists some ε with $B(x, 2\varepsilon) \subset D$. If $\xi_i \in B(x, \varepsilon)$ then, as both particles follow continuous paths until one exits D, we must have $V_t^1 \vee V_t^2 > \varepsilon$ for some $t \in [\tau_i, \tau_{i+1})$. So if $V_t^1, V_t^2 \to 0$ as $t \to \tau_\infty$ then $\xi_i \in B(x, \varepsilon)$ for only finitely many i. As x is arbitrary we see that so long as $V_t^1, V_t^2 \to 0$ as $t \to \tau_\infty, \xi_i$ can have no limit point in D.

It is convenient at this point to introduce some notation that will allow us to consider the behavior of \mathbf{X}_t when it is close to the boundary of a simplicial complex. Let σ be a *k*-simplex with vertices $\{v_0, \ldots, v_k\}$, that is

$$\sigma = \left\{ \sum_{i=0}^k \lambda_i v_i \lambda_0, \dots, \lambda_k \ge 0, \quad \sum_{i=0}^k \lambda_i = 1 \right\}.$$

Then define the *interior* of σ

$$\overset{\circ}{\sigma} = \left\{ \sum_{i=0}^{k} \lambda_i v_i \lambda_0, \dots, \lambda_k > 0, \quad \sum_{i=0}^{k} \lambda_i = 1 \right\}$$

and the *span* of σ to be the subspace

$$S_{\sigma} = \left\{ \sum_{i=0}^{k} \lambda_i v_i \sum_{i=0}^{k} \lambda_i = 0 \right\}.$$

For two simplices $\sigma_1, \sigma_2 \in \mathcal{K}$ we write $\sigma_1 \leq \sigma_2$ if σ_1 is a face of σ_2 and $\sigma_1 < \sigma_2$ if σ_1 is a proper face of σ_2 . We name the *star* of a simplex σ to be the set

$$St(\sigma) = \{\sigma_1 \in \mathcal{K} : \sigma_1 \ge \sigma\}$$

and define the *neighborhood* of σ as

$$\mathcal{N}(\sigma) = \{x \in \overline{D} : x \in \overset{\circ}{\sigma_1} \text{ for some } \sigma_1 \ge \sigma\}.$$

Given simplices $\sigma \leq \sigma_1$ name the vertices of σ and $\sigma_1, \{v_0, \ldots, v_k\}$ and $\{v_0, \ldots, v_n\}$ respectively. Define the *wedges*

$$\mathcal{W}(\sigma, \sigma_1) = \left\{ \sum_{i=0}^n \lambda_i v_i \lambda_{k+1}, \dots, \lambda_n > 0, \quad \sum_{i=0}^n \lambda_i = 1 \right\},$$
$$\mathcal{W}(\sigma) = \bigcup_{\sigma_1 \in St(\sigma)} \mathcal{W}(\sigma, \sigma_1).$$

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Notice that $\mathcal{N}(\sigma) \subset \overline{\mathcal{W}(\sigma)}$ and that $\mathcal{N}(\sigma)$ is open with respect to the subspace topology of \overline{D} . Notice also that $\mathcal{W}(\sigma)$ is a product space

$$\mathcal{W}(\sigma) = \mathcal{C}(\sigma) \times \mathcal{S}_{\sigma},$$

where the cone $\mathcal{C}(\sigma)$ is the projection of $\mathcal{W}(\sigma)$ onto $\mathcal{S}_{\sigma}^{\perp}$.

Now, consider $\sigma \in \partial \mathcal{K}$ and suppose there exists a subsequence $\xi_{i_n} \to \xi_{\infty} \in \overset{\circ}{\sigma}$. As $\xi_{\infty} \in \overset{\circ}{\sigma} \subset \mathcal{N}(\sigma)$ and $\mathcal{N}(\sigma)$ is open in \overline{D} we may assume without loss of generality that $\xi_{i_n} \in \mathcal{N}(\sigma)$ for each *n*. So consider \mathbf{X}_t started at (ξ_{i_n}, ξ_{i_n}) at time τ_{i_n} and stopped at the first time $T > \tau_{i_n}$ where one of X_t^1, X_t^2 exits $\mathcal{N}(\sigma)$. Of course, as $\mathcal{N}(\sigma) \subset \mathcal{W}(\sigma) \cap \overline{D}$, this has the same distribution as a Fleming-Viot process in $\mathcal{W}(\sigma)$ started and stopped in the same way.

So, let \mathbb{P}^x_{σ} and \mathbb{E}^x_{σ} be the probability measure and expectation operator associated with a Fleming-Viot process in $\mathcal{W}(\sigma)$ started at $\mathbf{X}_0 = (x, x)$. The \mathcal{S}_{σ} and $\mathcal{S}_{\sigma}^{\perp}$ components are not quite independent as they have the same jumps, but \mathbb{P}_{σ} allows a partial factorization as follows.

Lemma 8.5 If \mathbf{X}_t is a Fleming-Viot process in $\mathcal{W}(\sigma)$ then there is a well defined decomposition $\mathbf{X}_t = \mathbf{Y}_t + \mathbf{Z}_t$ with $\mathbf{Y}_t = (Y_t^1, Y_t^2) \in \mathcal{C}(\sigma)^2$, $\mathbf{Z}_t = (Z_t^1, Z_t^2) \in \mathcal{S}_{\sigma}^2$ with the following properties

- \mathbf{Y}_t is a Fleming-Viot process in $\mathcal{C}(\sigma)$;
- there exists a Brownian motion \tilde{Z}_t in S_{σ} (not adapted to the filtration of \mathbf{X}_t), independent of \mathbf{Y}_t , such that for each $i \in \mathbb{N}$ we have $\tilde{Z}_{\tau_i} = \zeta_i$ with $\zeta_i = Z_{\tau_i}^1 = Z_{\tau_i}^2$.

Proof Obviously, as $C(\sigma) \subset S_{\sigma}^{\perp}$, the factorization $\mathbf{X}_t = \mathbf{Y}_t + \mathbf{Z}_t$ is unique. Further, on each interval $[\tau_i, \tau_{i+1})$, the processes Y_t^1, Y_t^2, Z_t^1 and Z_t^2 evolve as independent Brownian motions on S_{σ}^{\perp} and S_{σ} respectively. So as S_{σ} is a subspace and has no boundary, X_t^j jumps when and only when Y_t^j hits $\partial C(\sigma)$, and so \mathbf{Y}_t is indeed a Fleming-Viot process on $C(\sigma)$.

Now for each $i \in \mathbb{N}$ only one of X_t^1, X_t^2 has a discontinuity at τ_{i+1} , so there is a well defined sequence of random variables $J_i \in \{1, 2\}$ such that $X_t^{J_i}$ is continuous on the closed interval $[\tau_i, \tau_{i+1}]$ and we may define a continuous process

$$\tilde{Z}_t = Z_t^{J_i}, \quad t \in \left[\tau_i, \tau_{i+1}\right].$$

Then $\tilde{Z}_{\tau_i} = \zeta_i$ for every *i* and it remains to show that \tilde{Z}_t is a Brownian motion independent of \mathbf{Y}_t . Of course \tilde{Z}_t is only defined up to τ_{∞} . But we may continue \tilde{Z}_t after τ_{∞} with an independent Brownian motion if necessary.

Now as \tilde{Z}_t follows either Z_t^1 or Z_t^2 then the quadratic variation $\langle \tilde{Z} \rangle_t = t$ and, by Lévy's characterization, we need only check that \tilde{Z}_t is a martingale with respect to its own natural filtration and is independent of \mathbf{Y}_t . Furthermore, although \tilde{Z}_t is not adapted to \mathbf{X}_t , for each τ_i , the path $\tilde{Z}|_{[0,\tau_i]}$ is measurable with respect to $\mathbf{X}|_{[0,\tau_i]}$. Therefore, by the strong Markov property, it is sufficient to consider only intervals $[\tau_i, \tau_{i+1})$. In fact it suffices to consider only the first time interval $[0, \tau_1)$. Let X_t be a Fleming-Viot process started at $\xi_0 \in \mathcal{W}(\sigma)$ and stopped at τ_1 . Then the left limit process is a pair of independent Brownian motions stopped at $\tau = \tau_1^-$. Set $J = J_0$ and we have $\xi_1 = X_\tau^J \in \mathcal{W}(\sigma)$ and $X_\tau^{3-J} \in \partial \mathcal{W}(\sigma)$.

So set $\mathbf{X}_t = \mathbf{Y}_t + \mathbf{Z}_t$ as in the statement of the lemma and let $\mathcal{F}_t^{\mathbf{Y}}, \mathcal{F}_t^{\mathbf{Z}}$ and $\mathcal{F}_t^{\tilde{Z}}$ be the natural filtrations of \mathbf{Y}, \mathbf{Z} and \tilde{Z} respectively. Set $\zeta_0 = Z_0^1, \zeta_1 = Z_\tau^J$ to be the $\mathcal{F}_{\tau}^{\mathbf{X}}$ -measurable \mathbf{Z} -components of ξ_0 and ξ_1 , respectively. Thus, τ is a stopping time of $\mathcal{F}_t^{\mathbf{Y}}$ and J is measurable with respect to $\mathcal{F}_{\tau}^{\mathbf{Y}}$. Now crucially \mathbf{Y} and \mathbf{Z} are independent processes so for $t < \tau$ we have

$$\mathbb{E}^{\xi_0}_{\sigma}\left(\zeta_1 \mid \mathcal{F}^{\mathbf{Y}}_{\tau} \lor \mathcal{F}^{\mathbf{Z}}_{t}\right) = \mathbb{E}^{\xi_0}_{\sigma}\left(Z^J_{\tau} \mid \mathcal{F}^{\mathbf{Y}}_{\tau} \lor \mathcal{F}^{\mathbf{Z}}_{t}\right) = \tilde{Z}_t.$$

Thus \tilde{Z} is a martingale, and hence a Brownian motion, with respect to the filtration $\mathcal{G}_t = \mathcal{F}_{\tau}^{\mathbf{Y}} \vee \mathcal{F}_t^{\mathbf{Z}}$. Therefore \tilde{Z} is independent of $\mathcal{F}_{\tau}^{\mathbf{Y}} \subset \mathcal{G}_0$ and is a Brownian motion with respect to its own natural filtration $\mathcal{F}_t^{\tilde{Z}} \subset \mathcal{G}_t$.

Now \mathbf{Y}_t is a process in a cone and if ξ_i converges to some point in $\overset{\circ}{\sigma}$ then \mathbf{Y}_t must converge to the apex of $\mathcal{C}(\sigma)$. Our next step is to show that this cannot be the case.

Lemma 8.6 If \mathbf{Y}_t is a Fleming-Viot process in a cone $C \subset \mathbb{R}^d$ then, with probability one, \mathbf{Y}_t does not converge to $(\underline{0}, \underline{0})$.

To prove this we will need to consider the angular components, $\Phi_t^j = \frac{Y_t^j}{\|Y_t^j\|}$, of **Y**. We will recall briefly some facts about spherical Brownian motion. We will omit details, which can be found in [23, Chapter 8], particularly Example 8.5.8.

Let B_t be a Brownian motion on \mathbb{R}^d , let the unit sphere be denoted

$$\mathbb{S}^{d-1} = \{ x \in \mathbb{R}^d : ||x|| = 1 \},\$$

and define the map $\phi : \mathbb{R}^d \setminus \{\underline{0}\} \to \mathbb{S}^{d-1}$ by $\phi(x) = \frac{x}{\|x\|}$. Now let $\Phi_t = \phi(B_t)$. Applying Ito's formula,

$$d\Phi_t = \frac{1}{\|B_t\|} \left(I - \Phi_t \Phi_t^\top \right) dB_t - \frac{d-1}{2\|B_t\|^2} \Phi_t dt.$$

Note we are interpreting Φ_t as a column vector so $\Phi_t \Phi_t^{\top}$ is a square matrix. Now define a differential operator $A : C^2(\mathbb{S}^{d-1}, \mathbb{R}) \to C^0(\mathbb{S}^{d-1}, \mathbb{R})$ by

$$Af(x) = \frac{1}{2} \left(\Delta f(x) - \sum_{i,j} x_i x_j \frac{\partial^2 f}{\partial x_i \partial x_j} \right) - \frac{d-1}{2} \sum_i x_i \frac{\partial f}{\partial x_i}.$$

Applying Ito's formula again, we see that $f(\Phi_t) - \int_0^t \frac{Af(\Phi_t)}{\|B_t\|^2} dt$ is a local martingale for each $f \in C^2(\mathbb{S}^{d-1}, \mathbb{R})$.

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We may extend this to functions of two Brownian motions by defining $\mathcal{A}^1, \mathcal{A}^2$ by

$$\mathcal{A}^1 f(x, y) = A(f(\cdot, y))(x),$$

$$\mathcal{A}^2 f(x, y) = A(f(x, \cdot))(y).$$

Then by a similar application of Ito's formula, if B_t^1 and B_t^2 are independent Brownian motions and $\Phi_t^1 = \phi(B_t^1)$, $\Phi_t^2 = \phi(B_t^2)$, $\Phi_t = (\Phi_t^1, \Phi_t^2)$, then

$$N_t^f = f\left(\Phi_t^1, \Phi_t^2\right) - \int_0^t \left(\frac{\mathcal{A}^1 f(\Phi_t)}{\|B_t^1\|^2} + \frac{\mathcal{A}^2 f(\Phi_t)}{\|B_t^2\|^2}\right) dt$$
(8.1)

is a local martingale.

Now apply a time change to Φ_t as follows. If $\alpha(t) = \inf\{s \in \mathbb{R}^+ : \int_0^s \|B_{\tilde{s}}\|^{-2} d\tilde{s} \ge t\}$, then $\Theta_t = \Phi_{\alpha(t)}$ is a Markov diffusion on \mathbb{S}^{d-1} with generator A. Let $\mathbb{P}_{\mathbb{S}}^{\theta_1,\theta_2}$ and $\mathbb{E}_{\mathbb{S}}^{\theta_1,\theta_2}$ be the probability measure and expectation operator associated with two independent copies of Θ_t started at θ_1 and $\theta_2 \in \mathbb{S}^{d-1}$ respectively.

Lemma 8.7 Let U be an open subset of \mathbb{S}^{d-1} and set

$$T_1^U = \inf\{t \in \mathbb{R} : \Theta_t^1 \in \partial U\},\$$

$$T_2^U = \inf\{t \in \mathbb{R} : \Theta_t^2 \in \partial U\},\$$

$$h^U(\theta_1, \theta_2) = \mathbb{P}_{\mathbb{S}}^{\theta_1, \theta_2}[T_1^U < T_2^U].$$

Then $h^U \in \mathcal{C}^2(U^2, \mathbb{R})$ and $\mathcal{A}^1 h^U = -\mathcal{A}^2 h^U \ge 0$.

Proof The process (Θ_t^1, Θ_t^2) is a Markov diffusion with generator $\mathcal{A}^1 + \mathcal{A}^2$, so by Dynkin's formula $\mathcal{A}^1 h^U + \mathcal{A}^2 h^U = 0$ and it remains to show that $\mathcal{A}^1 h^U \ge 0$.

By definition of the Markov generator

$$\begin{aligned} \mathcal{A}^{1}h^{U}(\theta_{1},\theta_{2}) &= \lim_{t \to 0} \frac{1}{t} \left(\mathbb{E}_{\mathbb{S}}^{\theta_{1},\theta_{2}} \left(h^{U} \left(\Theta_{t}^{1},\theta_{2} \right) \right) - h^{U} \left(\theta_{1},\theta_{2} \right) \right) \\ &= \lim_{t \to 0} \frac{1}{t} \left(\mathbb{E}_{\mathbb{S}}^{\theta_{1},\theta_{2}} \left(\mathbb{P}_{\mathbb{S}}^{\Theta_{t}^{1},\theta_{2}} \left[T_{1}^{U} < T_{2}^{U} \right] \right) - \mathbb{P}_{\mathbb{S}}^{\theta_{1},\theta_{2}} \left[T_{1}^{U} < T_{2}^{U} \right] \right). \end{aligned}$$

But $\mathbb{E}_{\mathbb{S}}^{\theta_1,\theta_2}(\mathbb{P}_{\mathbb{S}}^{\Theta_t^1,\theta_2}(\cdot))$ is the probability measure associated with the process $(\Theta_{s+t}^1, \Theta_s^2), s > 0$, obtained by giving Θ_1 a headstart. So we have

$$\mathbb{E}_{\mathbb{S}}^{\theta_{1},\theta_{2}}\left(\mathbb{P}_{\mathbb{S}}^{\Theta_{t}^{1},\theta_{2}}\left[T_{1}^{U} < T_{2}^{U}\right]\right) \geq \mathbb{P}_{\mathbb{S}}^{\theta_{1},\theta_{2}}\left[T_{1}^{U} - t < T_{2}^{U}\right] - \mathbb{P}_{\mathbb{S}}^{\theta_{1},\theta_{2}}\left[T_{1}^{U} < t\right]$$
$$\geq \mathbb{P}_{\mathbb{S}}^{\theta_{1},\theta_{2}}\left[T_{1}^{U} < T_{2}^{U}\right] - \mathbb{P}_{\mathbb{S}}^{\theta_{1},\theta_{2}}\left[T_{1}^{U} < t\right]$$

and, since $\frac{1}{t} \mathbb{P}^{\theta_1, \theta_2}_{\mathbb{S}}(T_1^U < t) \to 0$ as $t \to 0$, we may pass to the limit, and we see that $\mathcal{A}^1 h^U(\theta^1, \theta^2) \ge 0.$

We are ready to prove Lemma 8.6.

Proof of Lemma 8.6 Set $C = \{\lambda u : \lambda \in \mathbb{R}^+, u \in U\}$ for some open subset $U \subset \mathbb{S}^{d-1}$ and let \mathbf{Y}_t be a Fleming-Viot process in C.

We deal first with the special case when d = 1, in which case either $C = \mathbb{R}$ and there is nothing to prove or $C = \mathbb{R}^+$. If $C = \mathbb{R}^+$ then \mathbf{Y}_t is a 2-dimensional Brownian motion in the quarter plane with jumps $(y, 0) \mapsto (y, y)$ or $(0, y) \mapsto (y, y)$ whenever the process exits the first quadrant. As these jumps only increase $\|\mathbf{Y}_t\|$ then $\|\mathbf{Y}_t\|$ dominates a Bes(2) process and \mathbf{Y}_t does not converge to 0.

For $d \ge 2$ define a function

$$\mu(x) = \begin{cases} \log \|x\|, & \text{if } d = 2, \\ \frac{\|x\|^{2-d}}{2-d} & \text{if } d \ge 3, \end{cases}$$

and define processes

$$\begin{aligned} \Phi^1_t &= \phi(Y^1_t), & M^1_t &= \mu(Y^1_t), \\ \Phi^2_t &= \phi(Y^2_t), & M^2_t &= \mu(Y^2_t), \\ H_t &= h^U(\Phi^1_t, \Phi^2_t), & S_t &= M^1_t + (M^2_t - M^1_t)H_t. \end{aligned}$$

Now, μ is harmonic on \mathbb{R}^d and it will be key to our argument that M_t^1 and M_t^2 are both local martingales except when \mathbf{Y}_t jumps. We say a \mathbf{Y}_t -adatapted process R_t is a martingale between jumps if $R_t - \sum_{\{i \in \mathbb{N}: \tau_i \leq t\}} (R_{\tau_i} - R_{\tau_i})$ is a continuous local martingale. The process S_t is a convex combination of M_t^1 and M_t^2 , so if both Y_t^1 and Y_t^2 converge to the origin, then S_t converges to $-\infty$. Notice also that if Y^1 approaches ∂C then $H_t \to 1$ and so $S_t \to M_t^2$. Similarly, if Y_t^2 approaches the boundary then $S_t \to M_t^1$. So S_t is continuous.

Set

$$N_{s} = H_{s} - \int_{0}^{s} \left(\frac{\mathcal{A}^{1} h^{U}(\Phi_{t})}{\|B_{t}^{1}\|^{2}} + \frac{\mathcal{A}^{2} h^{U}(\Phi_{t})}{\|B_{t}^{2}\|^{2}} \right) dt$$

By (8.1) N_t is a martingale between jumps. We may check that the cross variation terms $\langle M^1, \Phi^1 \rangle_t = \langle M_t^2, \Phi_t^2 \rangle = 0$ and so, as H_t is a C^2 function of Φ_t^1 and Φ_t^2 , we have $\langle M^1, H \rangle_t = \langle M_t^2, H \rangle_t = 0$ and for $s \in [\tau_i, \tau_{i+1})$ we may calculate

$$S_{s} = S_{\tau_{i}} + \int_{\tau_{i}}^{s} (1 - H_{t}) dM_{t}^{1} + \int_{\tau_{i}}^{s} H_{t} dM_{t}^{2} + \int_{\tau_{i}}^{s} (M_{t}^{2} - M_{t}^{1}) dH_{t}$$

$$= S_{\tau_{i}} + \int_{\tau_{i}}^{s} (1 - H_{t}) dM_{t}^{1} + \int_{\tau_{i}}^{s} H_{t} dM_{t}^{2} + \int_{\tau_{i}}^{s} (M_{t}^{2} - M_{t}^{1}) dN_{t}$$

$$+ \int_{\tau_{i}}^{s} \left(M_{t}^{2} - M_{t}^{1}\right) \left(\frac{\mathcal{A}^{1}h^{U}(\Phi_{t})}{\|B_{t}^{1}\|^{2}} + \frac{\mathcal{A}^{2}h^{U}(\Phi_{t})}{\|B_{t}^{2}\|^{2}}\right) dt.$$

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Therefore $S_s - \int_{\tau_i}^s (M_t^2 - M_t^1) \left(\frac{\mathcal{A}^{1}h^U(\Phi_t)}{\|B_t^1\|^2} + \frac{\mathcal{A}^{2}h^U(\Phi_t)}{\|B_t^2\|^2} \right) dt$ is a martingale between jumps. Now from Lemma 8.7 we have $\mathcal{A}^{1}h^U = -\mathcal{A}^2h^U \ge 0$ and so

$$\frac{\mathcal{A}^{1}h^{U}(\Phi_{t})}{\|B_{t}^{1}\|^{2}} + \frac{\mathcal{A}^{2}h^{U}(\Phi_{t})}{\|B_{t}^{2}\|^{2}} = \mathcal{A}^{1}h^{U}(\Phi_{t})\left(\|B_{t}^{1}\|^{-2} - \|B_{t}^{2}\|^{-2}\right).$$

But μ is an increasing function of the norm $\|\cdot\|$, so for $\tau_i \leq s_1 \leq s_2 < \tau_{i+1}$,

$$\int_{s_1}^{s_2} \left(M_t^2 - M_t^1 \right) \left(\frac{\mathcal{A}^1 h^U(\Phi_t)}{\|B_t^1\|^2} + \frac{\mathcal{A}^2 h^U(\Phi_t)}{\|B_t^2\|^2} \right) dt \ge 0.$$

Therefore S_t is a continuous local submartingale and it cannot converge to $-\infty$. Thus \mathbf{Y}_t does not converge to (0, 0).

Corollary 8.8 If \mathbf{X}_t is a Fleming-Viot process in a polyhedral domain D then with probability one the sequence of jump points ξ_i does not converge to any $\xi_{\infty} \in \partial D$ as $i \to \infty$.

Proof First, for $\sigma \in \partial \mathcal{K}$, let F^{σ} be the event that $\xi_i \to \xi_{\infty}$ for some $\xi_{\infty} \in \overset{\circ}{\sigma}$ and assume without loss of generality that $\underline{0} \in \sigma$. Set

$$F_i^{\sigma} = F^{\sigma} \cap \left[X_t^j \in \mathcal{N}(\sigma); \ t \ge \tau_i, \quad j = 1, 2 \right]$$

Then, as $\mathcal{N}(\sigma)$ is open in \overline{D} , from Lemma 8.3, F_i^{σ} increases to F^{σ} up to an event of probability 0. By the strong Markov property and Lemma 8.6,

$$\mathbb{P}(F_i^{\sigma}) = \mathbb{P}_{\sigma}^{\xi_i} \left(\mathbf{Y}_t \to (\underline{0}, \underline{0}) \cap \left[X_t^j \in \mathcal{N}(\sigma); \ t \ge \tau_i, \quad j = 1, 2 \right] \right) = 0.$$

So as $\partial \mathcal{K}$ is a finite set of simplices we have $\mathbb{P}[\exists \xi_{\infty} \in \partial D \text{ s.t. } \xi_i \to \xi_{\infty} \text{ as } i \to \infty] = 0.$

To complete the proof of Theorem 8.2 we consider the set

$$L = \{ \sigma \in \mathcal{K} : \text{ there exists a subsequence } \xi_{i_n} \to \xi \in \overset{\circ}{\sigma} \text{ as } n \to \infty \}$$

It is easy to check that the event { $\sigma \in L$ } is **X**-measurable. We say σ is a *local maximum* of *L* if $L \cap St(\sigma) = \{\sigma\}$. Of course any non-empty subset of a finite lattice contains at least one local maximum, and *L* is non empty by compactness of \overline{D} . We will prove Theorem 8.2 by showing that for each $\sigma \in \mathcal{K}$ the event that $\tau_{\infty} < \infty$ and σ is a local maximum of *L* has probability 0.

Proof of Theorem 8.2 Fix $\sigma \in \partial \mathcal{K}$, and note that $\mathcal{N}(\sigma) \setminus \sigma$ is non empty. We show first that if ξ_i has a limit point in $\overset{\circ}{\sigma}$ and $\tau_{\infty} < \infty$, then ξ_i has a second limit point in $\mathcal{N}(\sigma) \setminus \sigma$.

First suppose that $\sigma = \{v\}$ is a vertex of \mathcal{K} and v is a limit point of ξ_i . By Corollary 8.8, the sequence ξ_i does not converge to v as $i \to \infty$, so we may choose $\varepsilon > 0$ such that $B(v, \varepsilon) \cap \overline{D} \subset \mathcal{N}(\sigma)$ and that $\|\xi_i - v\| > \varepsilon$ infinitely often. If this is the case then there are infinitely many pairs $(\xi_{i_n}, \xi_{i_{n+1}})$ such that $\xi_{i_n} \in B(v, \varepsilon)$ and $\xi_{i_{n+1}} \notin B(v, \varepsilon)$. But from Lemma 8.3 we have $\|\xi_i - \xi_{i+1}\| \to 0$ as $i \to \infty$ hence $\|\xi_{i_n} - v\| \to \varepsilon$ as $i \to \infty$. Therefore, as $\partial B(v, \varepsilon)$ is compact, ξ_i must have some limit point in $\partial B(v, \varepsilon) \cap \overline{D} \subset \mathcal{N}(\sigma) \setminus \{v\}$.

If σ is a *k*-simplex for 0 < k < d then for each $x \in \overset{\circ}{\sigma}$ there exists $\varepsilon > 0$ such that $B(x, 2\varepsilon) \cap \overline{D} \subset \mathcal{N}(\sigma)$. We will consider upcrossings of the interval $[\varepsilon, 2\varepsilon]$ by $||\xi_i - x||$. Define sequences $i_n, j_n \in \mathbb{N} \cup \{\infty\}$ and $T_n, \eta_n \in \mathbb{R} \cup \{\infty\}$ by: $j_0 = 0$,

$$i_{n+1} = \inf\{i > j_n : \xi_i \in B(x, \varepsilon)\},$$

$$j_n = \inf\{j > i_n : \xi_j \notin B(x, 2\varepsilon)\},$$

$$T_n = \inf\{t > \tau_{i_n} : X_t^1 \notin B(x, 2\varepsilon) \text{ or } X_t^2 \notin B(x, 2\varepsilon)\},$$

$$\eta_n = \sup\{\tau_i : \tau_i < T_n\}.$$

Then we put $N = \sup\{n \in \mathbb{N} : j_n < \infty\}$ to be the number of upcrossings.

Note that $B(x, 2\varepsilon) \cap \overline{D} \subset \mathcal{N}(\sigma)$ and so $\mathbf{X}_{(t+\tau_{i_n})\wedge T_n}$ is a Fleming-Viot process in $\mathcal{W}(\sigma)$ started at (ξ_{i_n}, ξ_{i_n}) and stopped on exiting $B(x, 2\varepsilon)$. So we may consider $\mathbb{P}_{\sigma}^{\xi_{i_n}}$ and factorize $\mathbf{X}_t = \mathbf{Y}_t + \mathbf{Z}_t$ as in Lemma 8.5. For $t \in [\tau_{i_n}, \eta_n]$, the process \tilde{Z}_t is measurable with respect to $\mathbf{X}|_{[\tau_{i_n}, T_n]}$ which is distributed according to $\mathbb{P}_{\sigma}^{\xi_{i_n}}$. Hence $\tilde{Z}|_{[\tau_{i_n}, \eta_n]}$ is a Brownian motion in \mathcal{S}_{σ} with respect to its own natural filtration.

Recall $\mathbf{Z}_{\tau_i} = (\zeta_i, \zeta_i)$ and set

$$\tilde{V}_t = \begin{cases} \|\tilde{Z}_t - \zeta_{i_n}\|, & \text{if } t \in [\tau_{i_n}, \eta_n], \\ 0, & \text{otherwise.} \end{cases}$$

Then \tilde{V}_t is dominated by a Bes(*d*) process reset to zero at times τ_{i_n} . So arguing as in the proof of Lemma 8.3, if $\tau_{\infty} < \infty$ and the number of upcrossings $N = \infty$, then $\tau_{i_n} < \tau_{\infty} < \infty$ for each $n \in \mathbb{N}$, and $\tilde{V}_t \to 0$ as $\tau_i \to \infty$. But $\eta_n = \sup_i \{\tau_i : \tau_i < T_n\}$, hence $\mathbf{X}_{\eta_n} = (\xi_{k_n}, \xi_{k_n})$ for some $k_n \in \mathbb{N}$ and either $||x - X_{T_n}^1|| = 2\varepsilon$ or $||x - X_{T_n}^2|| = 2\varepsilon$. So if $\tau_{\infty} < \infty$ and $N = \infty$, we must have $||\xi_{k_n} - x|| \to 2\varepsilon$ as $n \to \infty$ and so ξ_{k_n} has a limit point $\xi_{\infty} \in \partial B(x, 2\varepsilon) \cap \overline{D} \subset \mathcal{N}(\sigma)$. But $\tilde{V}_t \to 0$ as $t \to \tau_{\infty}$ with probability one, so we cannot have $\xi_{\infty} \in \sigma$ and we must have $\xi_{\infty} \in \mathcal{N}(\sigma) \setminus \sigma$.

Now let Q^{σ} be a countably dense subset of $\overset{\circ}{\sigma}$ and suppose ξ_i has some limit point $x \in \overset{\circ}{\sigma}$. By Corollary 8.8, ξ_i does not converge to x as $i \to \infty$ and we may choose some rational $\varepsilon > 0$ such that $B(x, 3\varepsilon) \cap \overline{D} \subset \mathcal{N}(\sigma)$ and $\|\xi_i - x\| > 3\varepsilon$ infinitely often. Now choose $q \in Q^{\sigma} \cap B(x, \varepsilon)$ and notice that $\|\xi_i - q\|$ makes infinitely many upcrossings of the interval $[\varepsilon, 2\varepsilon]$. If $\tau_{\infty} < \infty$ then as Q^{σ} is countable, with probability one we may find some limit point $\xi_{\infty} \in \mathcal{N}(\sigma) \setminus \sigma$.

Recall the definition of the set *L*. As *L* is nonempty there must exist some local maximum σ . However if $\tau_{\infty} < \infty$ then, by Corollary 8.4, we have $L \subseteq \partial \mathcal{K}$. We have just shown that if $\tau_{\infty} < \infty$ then *L* has no local maximum in $\partial \mathcal{K}$. Hence we must have $\tau_{\infty} = \infty$.

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