

Brownian motion conditioned to stay in a cone Rodolphe Garbit

▶ To cite this version:

Rodolphe Garbit. Brownian motion conditioned to stay in a cone. 2008. hal-00341032v2

HAL Id: hal-00341032 https://hal.science/hal-00341032v2

Preprint submitted on 28 Nov 2008 (v2), last revised 6 Jun 2009 (v3)

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

BROWNIAN MOTION CONDITIONED TO STAY IN A CONE

RODOLPHE GARBIT

ABSTRACT. A result of R. Durrett, D. Iglehart and D. Miller states that Brownian meander is Brownian motion conditioned to stay positive for a unit of time, in the sense that it is the weak limit, as x goes to 0, of Brownian motion started at x>0 and conditioned to stay positive for a unit of time. We extend this limit theorem to the case of multi-dimensional Brownian motion conditioned to stay in a smooth convex cone. Properties of the limit process are obtained and applications to random walks are given.

1. Introduction

The purpose of this paper is to prove the existence of a process which is, in some sense, multidimensional Brownian motion started at the vertex of a smooth convex cone and conditioned to stay in it for a unit of time.

Let \mathcal{C}_{∞} be the space of continuous functions $w:[0,+\infty)\to\mathbb{R}^d$, $d\geq 1$, endowed with the topology of uniform convergence on compact subsets induced by the distance

$$d(w, w') = \sum_{n>0} 2^{-n} (\max_{t \in [0, n]} ||w(t) - w'(t)|| \wedge 1) ,$$

and let \mathcal{F} be the corresponding Borel σ -algebra. We shall use \mathcal{C}_{∞} as a concise notation for $(\mathcal{C}_{\infty}, \mathcal{F})$. If (μ_n) , μ are probability measures on \mathcal{C}_{∞} , the sequence (μ_n) is said to converge weakly to μ if

$$\int f \, d\mu_n \to \int f \, d\mu$$

for all bounded and continuous function $f: \mathcal{C}_{\infty} \to \mathbb{R}$. Weak convergence will be denoted by the symbol \Rightarrow .

Let $\{X_t, t \geq 0\}$ be the canonical process on \mathcal{C}_{∞} for which $X_t(w) = w(t)$ for any $w \in \mathcal{C}_{\infty}$. Consider an open cone C with vertex at the origin 0 and let $\tau_C = \inf\{t > 0 : X_t \notin C\}$ be the first exit time of the canonical process from C. For any $x \in C$ we define the law $\widetilde{W}_{x,1}^C$ of the Brownian motion started at x and conditioned to stay in C for a unit of time by the formula

$$\widetilde{W}_{x,1}^{C}(*) = W_x(* | \tau_C > 1)$$
,

²⁰⁰⁰ Mathematics Subject Classification. Primary 60J65; Secondary 60B10.

Key words and phrases. Brownian motion, Brownian meander, conditioned Brownian motion, cone, weak convergence, limit theorems.

where W_x is the distribution on \mathcal{C}_{∞} of the standard d-dimensional Brownian motion started at x.

Our main result is the following theorem which states that the Brownian motion started at x and conditioned to stay in C for unit of time converges in law to a limit process as $x \in C$ tends to 0, when C is a *nice* cone. The precise definition of a nice cone is given in Section 5; for example, any circular or ellipsoidal cone is nice.

Theorem 1.1. Suppose C is a nice cone. As $x \in C$ goes to 0, the law $\widetilde{W}_{x,1}^C$ converges weakly on \mathcal{C}_{∞} to a limit $\widetilde{W}_{0,1}^C$.

For any $t \in (0,1]$, the entrance law $\widetilde{W}_{0,1}^C(X_t \in dy)$ has the density e(t,y) (w.r.t. Lebesque measure) given by the formula (23).

Theorem 1.1 is the multidimensional analog of Durrett, Iglehart and Miller result ([5], Theorem 2.1) in which they consider Brownian motion conditioned to stay positive for a unit of time and identify the limit as the Brownian meander. In the case of two-dimensional Brownian motion, Theorem 1.1 is due to Shimura ([8], Theorem 2).

For geometric reasons, the extension of Shimura's result to higher dimensions is not straightforward. Shimura calculates the finite-dimensional distributions and proves tightness. In order to prove the convergence of the finite-dimensional distributions, he obtains first a formula for the transitions of a two-dimensional Ornstein-Uhlenbeck process killed on the boundary of C and then derives by scaling the expression of the finite-dimensional distributions. Our proof rather uses directly an explicit formula for the heat kernel of a cone given by Bañuelos and Smits in [1]. But the main difference is to be found in the proof of tightness of the laws $\widetilde{W}_{x,1}^C$ as $x \in C \to 0$. We use a principle already present in Shimura's article: If $\widetilde{W}_{x,1}^C$ converges weakly as x tends to any point $x_0 \in \partial C \setminus \{0\}$, then the tightness as $x \in C \to 0$ follows. For a two-dimensional cone, proving weak convergence of $\widetilde{W}_{x,1}^C$ as $x \to x_0 \in \partial C \setminus \{0\}$ is quite easy because ∂C is locally linear at x_0 , so the proof is nearly the same as in the one-dimensional case. But in higher dimensions the geometry of the boundary of a cone is not so simple and we are led to a quite more general problem: Given an open set U and a point $x_0 \in \partial U$, does the law $\widetilde{W}_{x,1}^U$ of Brownian motion started at $x \in U$ and conditioned to stay in U for unit of time converge weakly as $x \in U \to x_0$? The major part of this paper is in fact concerned with the study of this question.

In Section 2, we consider the general problem of Brownian motion conditioned to stay in an open set U and give some useful properties of the conditioned laws $\widetilde{W}_{x,1}^U$, such as the Markov property and a form of continuity with respect to the variable x. In Section 3, we recall Durret-Iglehart-Miller result on Brownian motion conditioned to stay positive which extends immediately to the case of Brownian motion conditioned to stay in a half-plane. From the half-space case, we then derive in Section 4 a convergence theorem

for $\widetilde{W}_{x,1}^U$ as $x \to x_0 \in \partial U$ when U is nice at x_0 . This new result is based on the *ball estimate* (Lemma 4.6) which constitute the heart of this paper. Finally, in Section 5 we present a complete proof of Theorem 1.1, and we give some properties of the limit process, such as the distribution of its first exit time from the cone after time 1.

Our interest in such a limit theorem is related to random walks. Indeed, Bolthausen proved in [3] that a one-dimensional square integrable centered random walk conditioned to stay positive until time n converges to a Brownian meander, that is, a Brownian motion conditioned to stay positive for a unit of time. By analogy, the limit process given in Theorem 1.1 is very likely to be the limit in law of a multidimensional random walk conditioned to stay in a cone. Under the assumption that such a functional Central Limit Theorem holds, we derive in Section 5.3 information on the tail probability of the first exit time of the random walk from a cone. The functional CLT in itself will be the subject of a forthcoming article.

Unusual notations. If μ is a probability measure on a space (X, \mathcal{A}) , we will denote by $\mu(f)$ the expectation of a measurable function with respect to μ . For a set $A \in \mathcal{A}$ and a measurable function f, the notation $\mu(A; f)$ stands for $\mu(\mathbb{1}_A \times f)$, where $\mathbb{1}_A$ is the characteristic function of the set A. For consistency, $\mu(A; B)$ will often be preferred to $\mu(A \cap B)$.

2. Basic facts about the conditioned laws

2.1. Markov property. Let U be an open subset of \mathbb{R}^d and let τ_U be the first exit time from U. For any $x \in U$ and t > 0, we set

$$\widetilde{W}_{x,t}^{U}(*) = W_x(* | \tau_U > t) = \frac{W_x(*; \tau_U > t)}{W_x(\tau_U > t)}.$$

For convenience, we will also use the notation $\widetilde{W}_{x,t}^U := W_x$ for any $t \leq 0$ and $x \in \mathbb{R}^d$.

Let \mathcal{F}_t be the σ -algebra generated by the random variables $\{X_s, s \leq t\}$ and let $\mathcal{F}_{t+} = \bigcap_{s>t} \mathcal{F}_s$. The shift operator θ_t on \mathcal{C}_{∞} is defined by $\theta_t(w)(s) = w(t+s)$.

The laws \widetilde{W}^U inherit a Markov property from Brownian motion:

Proposition 2.1 (Markov property). Let $x \in U$ and t > 0. For all $s \ge 0$, $A \in \mathcal{F}_{s^+}$ and $B \in \mathcal{F}$,

$$\widetilde{W}^{U}_{x,t}\left(A;\theta_{s}^{-1}B\right)=\widetilde{W}^{U}_{x,t}\left(A;\widetilde{W}^{U}_{X(s),t-s}(B)\right)\;.$$

Proof. Suppose first that $s \in [0, t)$. By the Markov property of Brownian motion, we get

$$W_{x} (A; \theta_{s}^{-1}B; \tau_{U} > t)$$

$$= W_{x} (A; \tau_{U} > s; W_{X(s)}(B; \tau_{U} > t - s))$$

$$= W_{x} (A; \tau_{U} > s; \widetilde{W}_{X(s),t-s}^{U}(B)W_{X(s)}(\tau_{U} > t - s)).$$

Moreover, $\widetilde{W}^U_{X(s),t-s}(B)$ is \mathcal{F}_s -measurable, and the Markov property also gives

$$\begin{split} W_x\left(A;\widetilde{W}_{X(s),t-s}^U(B);\tau_U>t\right)\\ &=W_x\left(A;\widetilde{W}_{X(s),t-s}^U(B);\tau_U>s;W_{X(s)}(\tau_U>t-s)\right)\;. \end{split}$$

Thus we have

$$W_x\left(A;\theta_s^{-1}B;\tau_U>t\right)=W_x\left(A;\widetilde{W}_{X(s),t-s}^U(B);\tau_U>t\right)\;,$$

and the result follows.

Now suppose $s \geq t$. Since $\{\tau_U > t\}$ is \mathcal{F}_s -measurable, the Markov property gives

$$W_x\left(A;\theta_s^{-1}B;\tau_U>t\right)=W_x\left(A;W_{X(s)}(B);\tau_U>t\right)\ ,$$

which is the desired result.

We shall also need the related strong Markov property; the proof is quite similar and is omitted here.

Proposition 2.2 (Strong Markov property). Let $x \in U$ and t > 0. For any optional time τ , any $A \in \mathcal{F}_{\tau^+}$ and any positive measurable function $f(s, w) : [0, +\infty) \times \mathcal{C}_{\infty} \to \mathbb{R}$, we have

$$\widetilde{W}_{x,t}^{U}\left(A; \tau < t; f(\tau, \theta_{\tau})\right) = \widetilde{W}_{x,t}^{U}\left(A; \tau < t; \widetilde{W}_{X(\tau), t-\tau}^{U}(f(s, \cdot))_{|s=\tau}\right) .$$

2.2. Continuity. Let U be an open subset of \mathbb{R}^d . We will say that U is co-regular if $W_x(\tau_{\overline{U}} > 0) = 0$ for every $x \in \partial U$; that is, a Brownian motion started at any point of the boundary of U visits instantaneously the complement of \overline{U} . For such a set, τ_U and $\tau_{\overline{U}}$ are almost surely equal.

Proposition 2.3. Suppose U is co-regular. For every bounded continuous function f on C_{∞} , the mapping $(x,t) \mapsto \widetilde{W}_{x,t}^U(f)$ is continuous on $U \times (0,+\infty)$.

Proof. Since $W_x(\tau_U > t) > 0$ for any $(x, t) \in U \times (0, +\infty)$, it suffices to prove that the mapping $(x, t) \mapsto W_x(f; \tau_U > t)$ is continuous on $U \times (0, +\infty)$. Suppose $x_n \to x \in U$ and $t_n \to t > 0$, and set

$$\phi_n(w) = f(x_n - x + w) \mathbb{1}_{\{\tau_U > t_n\}} (x_n - x + w)$$

and

$$\phi(w) = f(w) \, 1\!\!1_{\{\tau_U > t\}}(w) \; .$$

Since $W_{x_n}(f; \tau_U > t_n) = W_x(\phi_n)$ and $W_x(f; \tau_U > t) = W_x(\phi)$, it is enough to show that $\phi_n(w) \to \phi(w)$ for W_x -almost every w.

Set $\Omega = \{X_0 = x; \tau_U = \tau_{\overline{U}} \neq t\}$ and choose a path $w \in \Omega$ and a sequence (w_n) that converges to w on \mathcal{C}_{∞} .

If $\tau_U(w) > t$, then there exists $\epsilon > 0$ such that $w(s) \in U$ for all $s \in [0, t + \epsilon]$. Since U is open and $w_n \to w$ uniformly on $[0, t + \epsilon]$, for n large enough we have $w_n(s) \in U$ for all $s \in [0, t + \epsilon]$ and we have $t_n < t + \epsilon$, thus $\tau_U(w_n) > t_n$. On the other hand, if $\tau_U(w) = \tau_{\overline{U}}(w) < t$, then there exists s < t such that $w(s) \notin \overline{U}$. If n is sufficiently large, then $w_n(s) \notin \overline{U}$ and $s < t_n$, so that $\tau_U(w_n) < t_n$. Hence, in each case,

$$1_{\{\tau_U > t_n\}}(w_n) \to 1_{\{\tau_U > t\}}(w)$$
.

Together with $W_x(\Omega) = 1$, this proves that $\phi_n(w) \to \phi(w)$ for W_x -almost every w.

2.3. Finite-dimensional distributions. Let $U \subset \mathbb{R}^d$ be a co-regular open set and x_0 a boundary point of U. Thanks to the Markov property and the continuity, it is now easy to find a sufficient condition for the weak convergence of the finite-dimensional distributions of $\widetilde{W}_{x,1}^U$, as x tends to x_0 , that only involves the first transitions.

Proposition 2.4. Suppose that for any $t \in (0,1)$, the first transition law $\widetilde{W}_{x,1}^U(X_t \in dy)$ converges weakly as $x \to x_0$ to a probability measure for which ∂U is a null set. Then, the finite-dimensional distributions of $\widetilde{W}_{x,1}^U$ converge weakly as $x \to x_0$ to some probability measures.

Proof. For all $t \in (0,1)$, let us denote by μ_t the limit of $\widetilde{W}_{x,1}^U(X_t \in dy)$ as $x \to x_0$. By Portemanteau Theorem,

$$\mu_t(\overline{U}) \ge \limsup_{x \to x_0} \widetilde{W}_{x,1}^U(X_t \in U) = 1$$
.

Since $\mu_t(\partial U) = 0$, we have $\mu_t(U) = 1$.

Now, consider times $0 < t_1 < t_2 < \cdots < t_n$ and a bounded continuous function $f: \mathbb{R}^n \to \mathbb{R}$, and set $F = f(X_{t_1}, X_{t_2}, \dots, X_{t_n})$. Fix $0 < t < \min(t_1, 1)$ and observe that for all $x \in U$ we have

$$\widetilde{W}_{x,1}^{U}(F) = \widetilde{W}_{x,1}^{U}\left(\widetilde{W}_{X_t,1-t}^{U}(f(X_{t_1-t}, X_{t_2-t}, \dots, X_{t_n-t}))\right)$$

$$= \widetilde{W}_{x,1}^{U}(H_t(X_t)),$$

with $H_t(y) = \widetilde{W}_{y,1-t}^U(f(X_{t_1-t}, X_{t_2-t}, \dots, X_{t_n-t}))$. By Proposition 2.3, the function $y \mapsto H_t(y)$ is continuous on U. Since $\widetilde{W}_{x,1}^U(X_t \in dy)$ converges weakly to $\mu_t(dy)$ as $x \to x_0$, it follows from the continuous mapping theorem that

$$\lim_{x \to x_0} \widetilde{W}_{x,1}^U(H_t(X_t)) = \mu_t(H_t) .$$

That is

$$\lim_{x \to x_0} \widetilde{W}_{x,1}^U(F) = \int \widetilde{W}_{y,1-t}^U(f(X_{t_1-t}, X_{t_2-t}, \dots, X_{t_n-t})) \, \mu_t(dy) \; .$$

The expression on the right side of this equation clearly defines a probability measure on \mathbb{R}^n , thus Proposition 2.4 is proved.

Let $x \in U$ and t > 0. The Markov property of Brownian motion gives

$$\begin{split} \widetilde{W}_{x,1}^{U}(X_{t} \in dy) &= \frac{W_{x}(X_{t} \in dy; \tau_{U} > 1)}{W_{x}(\tau_{U} > 1)} \\ &= \frac{W_{x}(X_{t} \in dy; \tau_{U} > t; W_{X_{t}}(\tau_{U} > 1 - t))}{W_{x}(\tau_{U} > 1)} \,. \end{split}$$

Since the transitions $W_x(X_t \in dy; \tau_U > t)$ of Brownian motion killed on the boundary of U have densities $p^U(t, x, y)$ with respect to Lebesgue measure dy, we get

$$\widetilde{W}_{x,1}^{U}(X_t \in dy) = \frac{p^{U}(t, x, y)}{W_x(\tau_U > 1)} W_y(\tau_U > 1 - t) dy$$
.

Hence, proving convergence of the finite-dimensional distributions of $\widetilde{W}_{x,1}^U$ consists essentially in finding an equivalent of the heat kernel $p^U(t,x,y)$ as $x \to x_0$.

2.4. Neat convergence. Let $U \subset \mathbb{R}^d$ be a co-regular open set and let x_0 be a boundary point of U. We shall now take the opposite view to the initial problem and exhibit conditions ensuring that some properties of the conditional distributions are preserved by weak limit.

Suppose that there exists a law $\widetilde{W}_{x_0,1}^U$ on \mathcal{C}_{∞} such that $\widetilde{W}_{x,1}^U \Rightarrow \widetilde{W}_{x_0,1}^U$ as $x \in U$ tends to x_0 . We will say that the convergence is neat (or that $\widetilde{W}_{x,1}^U$ converges neatly to $\widetilde{W}_{x_0,1}^U$ as $x \in U \to x_0$) if the limit process does not leave U before time 1, *i.e.* $\widetilde{W}_{x_0,1}^U(\tau_U > 1) = 1$. The next proposition gives a sufficient condition for neat convergence and states that the Markov property then holds for the limit process.

Proposition 2.5. Suppose $\widetilde{W}_{x,1}^U \Rightarrow \widetilde{W}_{x_0,1}^U$ as $x \in U \to x_0$. If $\widetilde{W}_{x_0,1}^U(X_t \in \partial U) = 0$ for all 0 < t < 1, then the convergence is neat, and the limit process $\widetilde{W}_{x_0,1}^U$ has the Markov property: For all t > 0, $A \in \mathcal{F}_{t^+}$ and $B \in \mathcal{F}$,

(1)
$$\widetilde{W}_{x_0,1}^U(A;\theta_t^{-1}B) = \widetilde{W}_{x_0,1}^U\left(A;\widetilde{W}_{X_t,1-t}^U(B)\right)$$
.

Proof. Let $t \in (0,1)$. By Portemanteau Theorem, we get

$$\widetilde{W}_{x_0,1}^U(X_t \in \overline{U}) \ge \limsup_{x \to x_0} \widetilde{W}_{x,1}^U(X_t \in U) = 1$$
.

Since $\widetilde{W}_{x_0,1}^U(X_t \in \partial U) = 0$, we see that $\widetilde{W}_{x_0,1}^U(X_t \in U) = 1$. Note that it is enough for the right hand side of equation (1) to be correctly defined. We will now prove that the Markov property holds for $\widetilde{W}_{x_0,1}^U$ and then use it to show that the convergence is neat.

Consider a finite sequence $0 < s_1 < s_2 < \cdots < s_n$ and a function $G = g(X_{s_1}, X_{s_2}, \dots, X_{s_n})$, where $g : \mathbb{R}^n \to \mathbb{R}$ is bounded and continuous. Choose $\epsilon \in (0, s_1)$, and consider another finite sequence $0 \le t_1 < t_2 < \cdots < t_m \le t_n$

 $t+\epsilon$ and a function $F = f(X_{t_1}, X_{t_2}, \dots, X_{t_m})$, where $f : \mathbb{R}^m \to \mathbb{R}$ is bounded and continuous. For every $x \in U$, we have

$$\widetilde{W}_{x,1}^{U}(F;G \circ \theta_{t}) = \widetilde{W}_{x,1}^{U}(F;g(X_{t+s_{1}},\ldots,X_{t+s_{n}}))$$

$$= \widetilde{W}_{x,1}^{U}\left(F;\widetilde{W}_{X(t+\epsilon),1-t-\epsilon}^{U}(g(X_{s_{1}-\epsilon},\ldots,X_{s_{n}-\epsilon}))\right)$$

$$= \widetilde{W}_{x,1}^{U}(F;H_{\epsilon}(X_{t+\epsilon})),$$

where $H_{\epsilon}(y) = \widetilde{W}_{y,1-t-\epsilon}^{U}(g(X_{s_1-\epsilon}, \dots, X_{s_n-\epsilon})).$ If $t + \epsilon \geq 1$, then H_{ϵ} is defined on \mathbb{R}^d by the formula

$$H_{\epsilon}(y) = W_{y}(g(X_{s_{1}-\epsilon}, \dots, X_{s_{n}-\epsilon}))$$

and is therefore continuous everywhere; thus the function $H_{\epsilon}(X_{t+\epsilon})$ is continuous on \mathcal{C}_{∞} .

If $t+\epsilon < 1$, then H_{ϵ} is defined and continuous on U by Proposition 2.3. Since the coordinate mapping $X_{t+\epsilon}$ is continuous on \mathcal{C}_{∞} and $\widetilde{W}_{x_0,1}^U$ -almost surely takes its values in U, we see that the mapping $H_{\epsilon}(X_{t+\epsilon})$ is almost-surely continuous with respect to $\widetilde{W}_{x_0,1}^U$.

Thus, in each case, the continuous mapping theorem gives

$$\lim_{x \to x_0} \widetilde{W}_{x,1}^U(F; H_{\epsilon}(X_{t+\epsilon})) = \widetilde{W}_{x_0,1}^U(F; H_{\epsilon}(X_{t+\epsilon})) .$$

We have also, by hypothesis

$$\lim_{x \to x_0} \widetilde{W}^U_{x,1}(F; G \circ \theta_t) = \widetilde{W}^U_{x_0,1}(F; G \circ \theta_t) \ .$$

Hence

(2)
$$\widetilde{W}_{x_0,1}^U(F;G\circ\theta_t) = \widetilde{W}_{x_0,1}^U(F;H_{\epsilon}(X_{t+\epsilon})).$$

By density, relation (2) holds for all $\mathcal{F}_{s+\epsilon}$ -measurable function F, and a fortiori for all $F \in \mathcal{F}_{s+}$. Now suppose $F \in \mathcal{F}_{s+}$. We shall let $\epsilon \to 0$ in (2). If $t \geq 1$, then $H_{\epsilon}(X_{t+\epsilon}) = W_{X(t+\epsilon)}(g(X_{s_1-\epsilon}, \ldots, X_{s_n-\epsilon}))$, and it is clear that $H_{\epsilon}(X_{t+\epsilon})(w) \to H_0(X_t)(w)$ as $\epsilon \to 0$ for all $w \in \mathcal{C}_{\infty}$.

If t < 1, then by Proposition 2.3, we have $\widetilde{W}^U_{w(t+\epsilon),1-t-\epsilon} \Rightarrow \widetilde{W}^U_{w(t),1-t}$ as $\epsilon \to 0$ for all $w \in \mathcal{C}_{\infty}$ such that $w(t) \in U$. Since the collection of those w forms a set of full $\widetilde{W}^U_{x_0,1}$ -measure, we see that $H_{\epsilon}(X_{t+\epsilon})$ converges $\widetilde{W}^U_{x_0,1}$ -almost surely to $H_0(X_t)$ as $\epsilon \to 0$.

In each case, the dominated convergence theorem gives

$$\lim_{\epsilon \to 0} \widetilde{W}_{x_0,1}^U(F; H_{\epsilon}(X_{t+\epsilon})) = \widetilde{W}_{x_0,1}^U(F; H_0(X_t)) .$$

Hence

$$\widetilde{W}_{r_0,1}^U(F;G\circ\theta_t)=\widetilde{W}_{r_0,1}^U(F;H_0(X_t))=\widetilde{W}_{r_0,1}^U(F;\widetilde{W}_{X_t,1-t}^U(G))$$
.

Here again, by density, the result holds for all $G \in \mathcal{F}$.

It remains to prove that $\widetilde{W}_{x_0,1}^U(\tau_U > 1) = 1$. By the Markov property of $\widetilde{W}_{x_0,1}^U$, we have

$$\widetilde{W}^U_{x_0,1}(\forall t\in(\epsilon,1],X_t\in U)=\widetilde{W}^U_{x_0,1}\left(\widetilde{W}^U_{X_{\epsilon,1-\epsilon}}(\tau_U>1-\epsilon)\right)=1$$

for all $\epsilon \in (0,1)$. The expected result follows by letting $\epsilon \to 0$.

Remark 2.6. Letting $t \to 0$ in (1) would give a zero-one law for $\widetilde{W}^U_{x_0,1}$ (i.e. $\widetilde{W}^U_{x_0,1}(A) = 0$ or 1 if $A \in \mathcal{F}_{0^+}$) if we had the stronger assumption that $\widetilde{W}^U_{x,t} \Rightarrow \widetilde{W}^U_{x_0,1}$ as $(x,t) \to (x_0,1)$. Note that in the special case where U = C is a cone, the last convergence follows from the hypothesis $\widetilde{W}^C_{x,1} \Rightarrow \widetilde{W}^C_{x_0,1}$ because of the scaling property of Brownian motion. More precisely, let K_t be the scaling operator defined for all $w \in \mathcal{C}_{\infty}$ by

$$K_t(w)(s) = \sqrt{t}w(s/t)$$
.

Recall that W_0 is K_t -invariant. From the scaling invariance of the cone C, it is easily checked that

$$\widetilde{W}_{x,t}^C = \widetilde{W}_{x/\sqrt{t},1}^C \circ K_t^{-1}$$
.

If $(x,t) \to (x_0,1)$, then $x/\sqrt{t} \to x_0$ and, by the continuous mapping theorem, we get

$$\widetilde{W}_{x,t}^C \Rightarrow \widetilde{W}_{x_0,1}^C \circ K_1^{-1} = \widetilde{W}_{x_0,1}^C \; .$$

Therefore, the zero-one law follows under the hypothesis of Proposition 2.5.

3. The half-space case

3.1. Brownian motion conditioned to stay positive. We will now recall the one dimensional theorem of Durret, Iglehart and Miller ([5], Theorem 2.1) and give a sketch of their proof. Auxiliary results such as Propositions 3.2 and 3.3 shall also be used in Section 4. Throughout this section we set d = 1 and we denote by $\tau_+ = \inf\{t > 0 : X_t \le 0\}$ the first exit time from the half-line $(0, +\infty)$. The related conditional laws will be denoted by $\widetilde{W}_{\sigma,1}^+$.

The Brownian meander is an inhomogeneous Markov process with continuous path that is obtained from Brownian motion by the following path transformation:

Let $\sigma = \max\{t < 1 : X_t = 0\}$ be the time of the last zero before time 1, and

$$\widetilde{X}_t = \frac{1}{\sqrt{1-\sigma}} |X(\sigma + t(1-\sigma))|$$
.

Then, with respect to Wiener measure W_0 , the process $(\widetilde{X}_t)_{t\geq 0}$ is the Brownian meander. Let $\widetilde{W}_{0,1}^+$ be the law of the Brownian meander on \mathcal{C}_{∞} . We have the following theorem.

Theorem 3.1 (Durrett, Iglehart, Miller). As x > 0 tends to 0, $\widetilde{W}_{x,1}^+ \Rightarrow \widetilde{W}_{0,1}^+$.

Note that Theorem 3.1 gives a precise meaning to the statement that Brownian meander is Brownian motion conditioned to stay positive until time 1.

The idea of the proof of Theorem 3.1 is to turn the conditioned laws into unconditioned ones by the mean of well-chosen sections of the original process. Let us give some details.

For all $x \geq 0$, introduce the random time

$$T_x = \inf\{t \ge 0 : X_t = x \text{ and } X_s > 0 \text{ for all } s \in (t, t+1]\}$$
.

These times are W_0 -almost surely finite. To see this, suppose first x > 0, and let $h = \inf\{t > 0 : X_t = x\}$ be the first hitting time of x and $g = \inf\{t > h : X_t = 0\}$ be the time of the first return to 0 after time h. Note that they are both stopping times. By the Markov property of Brownian motion, we have

$$W_0(g-h>1)=W_x(\tau_+>1)>0$$
.

We construct an increasing sequence of stopping times $h_1 < g_1 < h_2 < g_2 < \cdots$, setting $h_1 = h$, $g_1 = g$ and for all $i \ge 1$, $h_{i+1} = g_i + h \circ \theta_{g_i}$ and $g_{i+1} = g_i + g \circ \theta_{g_i}$. Then $(g_i - h_i)_{i \ge 1}$ is a sequence of i.i.d. random variables. Since $\sum W_0(g_i - h_i > 1) = +\infty$, we get $W_0(g_i - h_i > 1 \text{ i.o.}) = 1$ by the Borel-Cantelli lemma. In particular $W_0(T_x < +\infty) = 1$. Finally, we remark that $W_0(T_0 \le T_x) = 1$ for every x > 0, so that T_0 is also W_0 -almost surely finite.

To each time T_x we associate the shift operator $\phi_x := \theta_{T_x}$ acting on \mathcal{C}_{∞} . We then have:

Proposition 3.2. For every x > 0, $\widetilde{W}_{x,1}^+ = W_0 \circ \phi_x^{-1}$.

Proof. Consider a finite sequence of times $0 \le t_1 < t_2 < \cdots < t_n$, a sequence B_1, B_2, \ldots, B_n of Borel subsets of \mathbb{R} , and set $B = \{w \in \mathcal{C}_{\infty} : w(t_i) \in B_i, \forall i = 1, \dots, n\}$

Let $h_x = \inf\{t > 0 : X_t = x\}$ be the hitting time of x and write

(3)
$$W_0(\phi_x \in B) = W_0(\phi_x \in B; T_x = h_x) + W_0(\phi_x \in B; T_x > h_x)$$
.

By the Strong Markov property, we have

(4)
$$W_0(\phi_x \in B; T_x = h_x) = W_0(\theta_{h_x}^{-1}\{B; \tau_+ > 1\}) = W_x(B; \tau_+ > 1)$$
.

In order to compute the second term of the right hand side of (3), observe that

$$\{\phi_x \in B; T_x > h_x\} = \{\tau_+ \circ \theta_{h_x} \le 1; \theta_{T_x} \circ \theta_{\tau_+} \circ \theta_{h_x} \in B\},$$

so that applying successively the strong Markov property to h_x and then to τ_+ gives

$$W_{0}(\phi_{x} \in B; T_{x} > h_{x}) = W_{x}(\tau_{+} \leq 1; \theta_{T_{x}} \circ \theta_{\tau_{+}} \in B)$$

$$= W_{x}(\tau_{+} \leq 1; W_{0}(\theta_{T_{x}} \in B))$$

$$= W_{x}(\tau_{+} \leq 1)W_{0}(\phi_{x} \in B).$$
(5)

Putting (4) and (5) in equation (3) gives

$$W_0(\phi_x \in B) = W_x(B; \tau_+ > 1) + \{1 - W_x(\tau_+ > 1)\}W_0(\phi_x \in B)$$

Therefore

$$W_0(\phi_x \in B) = \frac{W_x(B; \tau_+ > 1)}{W_x(\tau_+ > 1)} = \widetilde{W}_{x,1}^+(B) .$$

Proposition 3.2 gives an "unconditioned" representation of the laws $\widetilde{W}_{x,1}^+$, x > 0. It is noteworthy that $W_0 \circ \phi_x^{-1}$ also make sense for x = 0 whereas the definition of $\widetilde{W}_{x,1}^+$ does not.

Proposition 3.3. As $x \to 0$, ϕ_x converges almost surely to ϕ_0 with respect to W_0 .

Proof. It suffices to prove that T_x almost surely converges to T_0 with respect to W_0 . Let w be a continuous path such that w(0) = 0 and $T_0 := T_0(w) < +\infty$. By definition of T_0 and continuity of w, if ϵ is small enough, then w(t) > 0 for all $t \in (T_0, T_0 + 1 + \epsilon]$. Put $\eta = w(T_0 + \epsilon) > 0$. For all $x \in [0, \eta)$, there exists $t \in [T_0, T_0 + \epsilon)$ such that w(t) = x. Since w(s) > 0 for all $s \in (t, t + 1]$, we see that $T_x(w) \leq T_0 + \epsilon$. On the other hand it is clear that $T_0 \leq T_x(w)$. Hence $T_x(w) \to T_0$ as $x \to 0$.

From Propositions 3.2 and 3.3, it follows by the dominated convergence theorem that

$$\widetilde{W}_{x,1}^+ = W_0 \circ \phi_x^{-1} \Rightarrow W_0 \circ \phi_0^{-1}$$

as $x \to 0$. Note that the limit law clearly satisfies $W_0 \circ \phi_0^{-1}(\tau_+ > 1) = 1$; hence the convergence is neat. In order to prove Theorem 3.1, it remains to identify the limit with the Brownian meander. This can be done by computing the limit of the finite-dimensional distributions of the laws $\widetilde{W}_{x,1}^+$ which are easily derived from a classical formula for the joint distribution of Brownian motion and its minimum. We do not give further detail since no expression of these finite-dimensional distributions will be needed in what follows.

3.2. Brownian motion conditioned to stay in a half-space. Theorem 3.1 can easily be extended to multidimensional Brownian motion conditioned to stay in a half-space. Let $d \geq 2$. Because of invariance properties of d-dimensional Brownian motion we need only to study the case of the half-space $D = \{x \in \mathbb{R}^d : x_1 > 0\}$. Let BM be a Brownian meander and B_2, \ldots, B_d be one-dimensional Brownian motions such that BM, B_2, \ldots, B_d

are mutually independent. The d-dimensional process (BM, B_2, \ldots, B_d) will be called D-Brownian meander and its law will be denoted by $\widetilde{W}_{0,1}^D$.

Corollary 3.4. As
$$x \in D \to 0$$
, $\widetilde{W}_{x,1}^D \Rightarrow \widetilde{W}_{0,1}^D$.

Proof. A Brownian motion conditioned to stay in the half-space D is a Brownian motion whose first coordinate is conditioned to stay positive. Since the coordinates are independent one-dimensional Brownian motions, the result follows immediately from Theorem 3.1.

Remark 3.5. It is clear from the definition of the D-Brownian meander that it satisfies $\widetilde{W}_{0,1}^D(\tau_D > 1) = 1$; thus the convergence in Theorem 3.4 is neat. Therefore, $\widetilde{W}_{0,1}^D$ has the Markov property of Proposition 2.5. Moreover, since D is a cone, it follows from Remark 2.6 that we also have a zero-one law with respect to $\widetilde{W}_{0,1}^D$.

Remark 3.6. For
$$x = (x_1, x_2, \dots, x_d) \in \overline{D}$$
, define a process Z_x by $\forall t \geq 0$, $Z_x(t) = x + X(T_{x_1}(X_1) + t) - X(T_{x_1}(X_1))$.

It is not difficult to verify that, under the Brownian distribution, the law of the random process Z_x is precisely $\widetilde{W}_{x,1}^D$. This alternative representation of the laws $\widetilde{W}_{x,1}^D$ will be used in the next section.

4. Preconditioning

We shall now use the results of Section 3 in order to obtain a convergence theorem for the Brownian motion conditioned to stay in a set satisfying some regularity and convexity assumptions (Theorem 4.11). Section 4.1 introduces the idea of preconditioning and explains how it could be applied to the convergence problem. The proposed method requires two estimates that are studied in Sections 4.2 and 4.3. This finally leads us to introduce the class of *nice sets* for which we solve the convergence problem in Section 4.4.

4.1. Changing laws for the convergence problem. Let $U \subset \mathbb{R}^d$ be a co-regular open set with $0 \in \partial U$. Recall that the definition of $\widetilde{W}_{x,1}^U$ by the formula

$$\widetilde{W}_{x,1}^{U}(*) = \frac{W_x(*; \tau_U > 1)}{W_x(\tau_U > 1)},$$

does not make any sense for x = 0 since $W_0(\tau_U > 1) = 0$.

Now suppose that U is contained in the half-space D. Then a Brownian motion conditioned to stay in U is also a Brownian motion conditioned to stay in D and then conditioned to stay in U, that is:

(6)
$$\widetilde{W}_{x,1}^{U}(*) = \widetilde{W}_{x,1}^{D}(* \mid \tau_{U} > 1) = \frac{\widetilde{W}_{x,1}^{D}(*; \tau_{U} > 1)}{\widetilde{W}_{x,1}^{D}(\tau_{U} > 1)}.$$

This simple identity is what we call *preconditioning*, for if we take it as a definition, it is the same as before except we have changed the initial law

of the paths $(W_x \leftrightarrow \widetilde{W}_{x,1}^D)$ which are now preconditioned to stay in D. The gain is that, although $W_0(\tau_U > 1) = 0$, we might have $\widetilde{W}_{0,1}^D(\tau_U > 1) > 0$ if the boundary of U is smooth enough at 0. Proposition 4.3 in Section 4.2 gives a sufficient condition on U which ensures that $\widetilde{W}_{0,1}^D(\tau_U > 1) > 0$. If so, we will set $\widetilde{W}_{0,1}^U(*) := \widetilde{W}_{0,1}^D(*|\tau_U > 1)$.

At this point, we have to note that $\widetilde{W}_{0,1}^D(\tau_U > 1) > 0$ is equivalent to $\widetilde{W}_{0,1}^D(\tau_U > 0) = 1$. The last condition is necessary for if $\widetilde{W}_{0,1}^D(\tau_U > 0)$ was < 1, then it would be 0 by the zero-one law, and $\widetilde{W}_{0,1}^D(\tau_U > 1)$ would also be equal to 0. The condition is also sufficient: It implies that there exists an $\epsilon > 0$ such that $\widetilde{W}_{0,1}^D(\tau_U > \epsilon) > 0$, and by the Markov property we then have

$$\widetilde{W}_{0,1}^D(\tau_U > 1) = \widetilde{W}_{0,1}^D\left(\tau_U > \epsilon; \widetilde{W}_{X_{\epsilon}, 1 - \epsilon}^D(\tau_U > 1 - \epsilon)\right) > 0$$

since $W_x(\tau_U > 1 - \epsilon)$ is > 0 for all $x \in U$.

The problem we then have to solve is the following: Since $\widetilde{W}_{x,1}^D \Rightarrow \widetilde{W}_{0,1}^D$ and $\widetilde{W}_{0,1}^D(\tau_U > 1) > 0$, do we have

$$\widetilde{W}_{x,1}^{D}(*|\tau_{U}>1) \Rightarrow \widetilde{W}_{0,1}^{D}(*|\tau_{U}>1)$$
?

The next lemma shows that the answer is positive when U is locally a half-space at 0.

Lemma 4.1. Let $V \subset \mathbb{R}^d$ be an open co-regular set such that $0 \in V$. For every bounded continuous function $f: \mathcal{C}_{\infty} \to \mathbb{R}$, we have

$$\lim_{x \to 0} \widetilde{W}_{x,1}^{D}(f; \tau_V > 1) = \widetilde{W}_{0,1}^{D}(f; \tau_V > 1) .$$

In addition, we have $\widetilde{W}_{0,1}^D(\tau_V > 1) > 0$; thus, as $x \to 0$,

$$\widetilde{W}_{x,1}^{D\cap V}\Rightarrow\widetilde{W}_{0,1}^{D\cap V}$$
 .

Proof. Set $\Omega = \{X_0 = 0, \tau_V = \tau_{\overline{V}} \neq 1\}.$

We first prove that $\widetilde{W}_{0,1}^D(\Omega) = 1$. For all $\epsilon \in (0,1)$, we have

$$\widetilde{W}_{0,1}^D(\epsilon < \tau_V < \tau_{\overline{V}}) = \widetilde{W}_{0,1}^D(\widetilde{W}_{X_\epsilon,1-\epsilon}^D(\tau_V < \tau_{\overline{V}})) = 0$$

since $W_x(\tau_V < \tau_{\overline{V}}) = 0$ for all $x \in \mathbb{R}^d$. Hence, letting $\epsilon \to 0$ gives

$$\widetilde{W}_{0,1}^D(\tau_V < \tau_{\overline{V}}) = 0 \ .$$

To see that $\widetilde{W}_{0,1}^D(\tau_V=1)=0$, fix $\epsilon\in(0,1)$ and write

$$\widetilde{W}_{0,1}^D(\tau_V=1) = \widetilde{W}_{0,1}^D(\widetilde{W}_{X_{\epsilon,1-\epsilon}}^D(\tau_V=1-\epsilon))$$
.

The result follows immediatly because $W_x(\tau_V = 1 - \epsilon) = 0$ for all $x \in \mathbb{R}^d$.

Then, we prove the following statement: For all $w \in \Omega$ and every sequence $(w_n) \in C_{\infty}$ such that $w_n \to w$, we have

(7)
$$\mathbb{1}_{\{\tau_V > 1\}}(w_n) \to \mathbb{1}_{\{\tau_V > 1\}}(w)$$
.

If $\tau_V(w) > 1$, then $w(t) \in U$ for all $t \in [0, 1]$ and for n large enough $w_n(t) \in U$ for all $t \in [0, 1]$; thus $\tau_V(w_n) > 1$.

If $\tau_V(w) = \tau_{\overline{V}}(w) < 1$, then there exists $t \in (0,1)$ such that $w(t) \notin \overline{V}$. For n large enough we then have $w_n(t) \notin \overline{V}$ and so $\tau_V(w_n) < 1$. This proves (7).

We can now apply the continuous mapping theorem which gives

$$\widetilde{W}_{x,1}^D(f;\tau_V>1) \to \widetilde{W}_{0,1}^D(f;\tau_V>1), \quad \text{as } x\to 0 \ ,$$

for any bounded continuous function $f: \mathcal{C}_{\infty} \to \mathbb{R}$, and the lemma is proved.

If U is not locally a half-space at 0, then τ_U is everywhere discontinuous. A way to get round this difficulty is given by the next lemma which will be the basic tool for applying preconditioning to the convergence problem.

Lemma 4.2. Suppose $\widetilde{W}_{0,1}^D(\tau_U > 1) > 0$ and let (x_n) be a sequence of points in U such that $x_n \to 0$. If we have

$$\lim_{s\to 0} \limsup_{n} \widetilde{W}_{x_n,1}^D(\tau_U \le s) = 0 ,$$

then

$$\lim_{n \to 0} \widetilde{W}_{x_n,1}^D(f; \tau_U > 1) = \widetilde{W}_{0,1}^D(f; \tau_U > 1)$$

for every bounded continuous function $f: \mathcal{C}_{\infty} \to \mathbb{R}$. As a consequence $\widetilde{W}_{x_n,1}^U \Rightarrow \widetilde{W}_{0,1}^U$.

Proof. We follow Shimura's proof ([7], Lemma 4.1). Fix $s \in (0,1)$ and introduce the random times $\tau_U^s = \inf\{t \geq s : X_t \notin U\}$ and $\tau_{\overline{U}}^s = \inf\{t \geq s : X_t \notin \overline{U}\}$. Set $\Omega = \{\tau_U^s = \tau_{\overline{U}}^s \neq 1\}$.

We shall prove that $\widetilde{W}_{0,1}^D(\Omega) = 1$. By the Markov property we have

$$\widetilde{W}_{0,1}^{D}(\tau_{U}^{s}=1) = \widetilde{W}_{0,1}^{D}\left(X_{s} \in U; \widetilde{W}_{X_{s},1-s}^{D}(\tau_{U}=1-s)\right) = 0$$

because $W_x(\tau_U = 1 - s) = 0$ for all $x \in \mathbb{R}^d$. We have also

$$\widetilde{W}_{0,1}^D(\tau_U^s < \tau_{\overline{U}}^s)$$

$$= \ \widetilde{W}_{0,1}^D(X_s \in U; \tau_U^s < \tau_{\overline{U}}^s) + \widetilde{W}_{0,1}^D(X_s \notin U; s < \tau_{\overline{U}}^s)$$

$$= \ \widetilde{W}_{0,1}^D \left(X_s \in U; \widetilde{W}_{X_s,1-s}^D(\tau_U < \tau_{\overline{U}}) \right) + \widetilde{W}_{0,1}^D \left(X_s \in \partial U; \widetilde{W}_{X_s,1-s}^D(\tau_{\overline{U}} > 0) \right)$$

= 0

since $W_x(\tau_U < \tau_{\overline{U}}) = 0$ for all $x \in \mathbb{R}^d$ and $W_x(\tau_{\overline{U}} > 0) = 0$ for all $x \in \partial U$ (remember that U is co-regular). Hence $\widetilde{W}_{0,1}^D(\Omega) = 1$.

Then, we shall prove the following statement: For all $w \in \Omega$ and every sequence $(w_n) \in \mathcal{C}_{\infty}$ such that $w_n \to w$, we have

(8)
$$\mathbb{1}_{\{\tau_U^s > 1\}}(w_n) \to \mathbb{1}_{\{\tau_U^s > 1\}}(w)$$
.

If $\tau_U^s(w) > 1$, then $w(t) \in U$ for all $t \in [s, 1]$ and for n large enough $w_n(t) \in U$ for all $t \in [s, 1]$; thus $\tau_U^s(w_n) > 1$.

If $\tau_{\overline{U}}^s(w) = \tau_{\overline{U}}^s(w) < 1$, then there exists $t \in [s,1)$ such that $w(t) \notin \overline{U}$. For n large enough we then have $w_n(t) \notin \overline{U}$ and so $\tau_{\overline{U}}^s(w_n) < 1$. This proves (8).

For any bounded continuous function $f: \mathcal{C}_{\infty} \to \mathbb{R}$, the continuous mapping theorem gives

$$\widetilde{W}_{x_{n},1}^{D}(f;\tau_{U}^{s}>1) \to \widetilde{W}_{0,1}^{D}(f;\tau_{U}^{s}>1)$$
.

If B is a bound for |f|, then

$$\begin{split} |\widetilde{W}_{x_{n},1}^{D}(f;\tau_{U}>1) - \widetilde{W}_{0,1}^{D}(f;\tau_{U}>1)| \\ & \leq |\widetilde{W}_{x_{n},1}^{D}(f;\tau_{U}^{s}>1) - \widetilde{W}_{0,1}^{D}(f;\tau_{U}^{s}>1)| \\ & + B\left(\widetilde{W}_{x_{n},1}^{D}(\tau_{U}\leq s) + \widetilde{W}_{0,1}^{D}(\tau_{U}\leq s)\right) \; . \end{split}$$

Hence

$$\limsup_{n} |\widetilde{W}_{x_{n},1}^{D}(f;\tau_{U}>1) - \widetilde{W}_{0,1}^{D}(f;\tau_{U}>1)|$$

$$\leq B \left(\limsup_{n} \widetilde{W}_{x_{n},1}^{D}(\tau_{U}\leq s) + \widetilde{W}_{0,1}^{D}(\tau_{U}\leq s)\right).$$

Therefore, using the hypothesis of the Lemma and the fact that $\widetilde{W}_{0,1}^D(\tau_U = 0) = 0$, we obtain the announced result by letting s go to 0.

4.2. **An irregularity criterion.** In order to apply the ideas of preconditioning to a set U, we have to know whether $\widetilde{W}_{0,1}^D(\tau_U > 0) = 1$ or not. Such a criterion was discovered independently by Shimura ([7], Lemma 3.1) in the two dimensional case, and by Burdzy ([4], Corollary 3.1) when studying excursions from hyperplanes and smooth surfaces.

Suppose $d \geq 2$. Let $h: [0, +\infty) \to [0, +\infty)$ be a continuous function and set

$$U_h = \left\{ x \in \mathbb{R}^d : x_1 > h\left(\sqrt{x_2^2 + \dots + x_d^2}\right) \right\} .$$

We then have:

Proposition 4.3. If h(r)/r, r > 0, is non-decreasing in a neighborhood of 0, then $\widetilde{W}_{0,1}^D(\tau_{U_h} > 0) = 1$ if and only if $\int_0^1 h(r)r^{-2} dr < \infty$.

For the reader convenience, we will now sketch a proof of Proposition 4.3. We go back first to dimension d=1. Let n>0 and let $h_n=\inf\{t>0:X_t=1/n\}$ be the first hitting time of 1/n. Let $a_{1,n}=\max\{t< h_n: X_t=0\}$ be the time of the last return to 0 before time h_n and let $b_{1,n}=\inf\{t>h_n:X_t=0\}$ be the first hitting time of 0 after time h_n . Note that $b_{1,n}$ is a stopping time. We construct an increasing sequence of random time $a_{1,n}< b_{1,n}< a_{2,n}< b_{2,n}<\cdots$ by setting $a_{j+1,n}=a_{1,n}\circ\theta_{b_{j,n}}+b_{j,n}$ and $b_{j+1,n}=b_{1,n}\circ\theta_{b_{j,n}}+b_{j,n}$. Since the Brownian motion starts afresh at $b_{j,n}$, the distribution of the process $X(a_{j,n}+\cdot)$ (w.r.t. W_0) does not depend on j.

Now recall that Brownian meander is the section of Brownian motion which starts at time

$$T_0 = \inf\{t > 0 : X_t = 0 \text{ and } X_s > 0 \text{ for all } s \in (t, t+1]\}$$
.

If $w \in$

 C_{∞} is such that $T_0(w) < +\infty$, then there exists an integer n > 0 such that, after time $T_0(w)$, the path w hits 1/n before hitting 0. Hence, $T_0(w)$ is one of the times $a_{j,n}(w)$.

Let $d \geq 2$ and consider the d-dimensional processes $Z_{j,n}$ and Z_0 defined by

$$Z_{j,n}(t) = (X_1(a_{j,n}(X_1) + t), X_2(t), \dots, X_d(t))$$

and

$$Z_0(t) = (X_1(T_0(X_1) + t), X_2(t), \dots, X_d(t)).$$

We know that Z_0 has the law of a D-Brownian meander with respect to W_0 . Hence, for an open set U, we have

$$\widetilde{W}_{0,1}^D(\tau_U > 0) = W_0(\tau_U(Z_0) > 0)$$
.

Since there almost surely exists j, n > 0 such that $T_0 = a_{j,n}$, the condition $W_0(\tau_U(Z_0) > 0) = 1$ is implied by the condition

$$W_0(\tau_U(Z_{j,n}) > 0 \text{ for all } j, n > 0) = 1$$

which is equivalent to

$$\forall j, n > 0, \quad W_0(\tau_U(Z_{j,n}) > 0) = 1.$$

As the distribution of the processes $Z_{j,n}$ does not depend on j, this is also equivalent to

$$\forall n > 0, \quad W_0(\tau_U(Z_{1,n}) > 0) = 1.$$

Recall that $a_{1,n}(X_1)$ is the time of the last return to 0 of X_1 before it hits 1/n. We know from Williams (see [10]) that a Brownian motion taken between those two times is a 3-dimensional Bessel process taken between time 0 and its first hitting time of 1/n. Therefore

$$W_0(\tau_U(Z_{1,n}) > 0) = W_0(\tau_U(BS, X_2, \dots, X_d) > 0)$$
,

where BS is a 3-dimensional Bessel process independent from X_2, \ldots, X_d . Since BS is the norm of a 3-dimensional Brownian motion, we can rewrite the last equality as

$$W_0(\tau_U(Z_{1,n}) > 0) = W_0(\tau_U(\sqrt{X_1^2 + X_2^2 + X_3^2}, X_4, \dots, X_{d+2}) > 0)$$
.

Now, for any set $U \subset \mathbb{R}^d$, put

$$U^* = \{x \in \mathbb{R}^{d+2} : (\sqrt{x_1^2 + x_2^2 + x_3^2}, x_4, \dots, x_{d+2}) \in U\}$$
.

We have just proved that $W_0(\tau_{U^*} > 0) = 1$ implies $\widetilde{W}_{0,1}^D(\tau_U > 1) = 1$. The same arguments hold if the event $\{\tau_U > 0\}$ is replaced by $\{\tau_U = 0\}$, and we have in fact an equivalence:

$$W_0(\tau_{U^*} > 0) = 1 \Leftrightarrow \widetilde{W}_{0,1}^D(\tau_U > 1) = 1$$
.

Thanks to this duality, Proposition 4.3 simply follows from a classical irregularity criterion for Brownian motion (e.g. Port and Stone [6], Proposition 3.5).

Remark 4.4. Although the criterion given in Proposition 4.3 is designed for the sets U_h , it can be applied to any set U that coincides with a set U_h in a neighborhood of 0: If V is a neighborhood of 0 then $\widetilde{W}_{0,1}^D(\tau_V > 0) = 1$ and by consequence

$$\widetilde{W}_{0,1}^D(\tau_U > 0) = \widetilde{W}_{0,1}^D(\tau_{U \cap V} > 0)$$
.

Thus, if $U \cap V = U_h \cap V$,

$$\widetilde{W}_{0,1}^{D}(\tau_{U}>0)=\widetilde{W}_{0,1}^{D}(\tau_{U_{h}}>0)$$
.

As an example of application, the ball B with center at $(1,0,\ldots,0)$ and radius 1 verifies $\widetilde{W}_{0,1}^D(\tau_B>0)=1$. On the contrary, for a proper cone $C\subset D$ with vertex at the origin, Proposition 4.3 shows that $\widetilde{W}_{0,1}^D(\tau_C>0)=0$, hence a law $\widetilde{W}_{0,1}^C$ can not be defined directly.

4.3. **The ball estimate.** In this section we will prove a fundamental estimate (Lemma 4.6) which solves the convergence problem for the Brownian motion conditioned to stay in a ball. This estimate will also play a key role in the study of Brownian motion conditioned to stay in nice sets (see Section 4.4).

Fix $d \geq 2$. Let D be the half-space $\{x \in \mathbb{R}^d : x_1 > 0\}$ and B the open ball with center at $e_1 = (1, 0, \dots, 0)$ and radius 1. As we have already seen, Proposition 4.3 shows that $\widetilde{W}_{0,1}^D(\tau_B > 0) = 1$.

Set $\widetilde{W}_{0,1}^B(*) = \widetilde{W}_{0,1}^D(* | \tau_B > 1)$. According to Lemma 4.2, in order to prove that $\widetilde{W}_{x,1}^B \Rightarrow \widetilde{W}_{0,1}^B$ as $x \in B \to 0$, it would be sufficient that

(9)
$$\lim_{s \to 0} \limsup_{n \to \infty} \widetilde{W}_{x_n,1}^D(\tau_B \le s) = 0.$$

for every sequence $(x_n) \in B$ such that $x_n \to 0$. But (9) does not hold in general: for any $x \in \partial B \setminus \{0\}$ we have $\widetilde{W}_{x,1}^D(\tau_B = 0) = 1$, thus it is possible to find sequences $(x_n) \in B$ and (s_n) such that $x_n \to 0$, $s_n \to 0$ and $\widetilde{W}_{x_n,1}^D(\tau_B \leq s_n) \to 1$.

To overcome this difficulty, we will only prove that (9) holds for every sequence $(\lambda_n e_1)$ with $\lambda_n \to 0$ and $e_1 = (1, 0, \dots, 0)$, and then deduce the convergence of $\widetilde{W}_{x_n,1}^B$ -for a general sequence (x_n) - from the convergence of $\widetilde{W}_{\lambda_n e_1,1}^B$ thanks to invariance properties of Brownian motion.

We start with elementary geometry. Let E be the set of all (d-1)-uples $(\epsilon_2, \ldots, \epsilon_d)$ with $\epsilon_i = \pm 1$. For all $\epsilon = (\epsilon_2, \ldots, \epsilon_d) \in E$, let $\overline{\epsilon} = (-\epsilon_2, \ldots, -\epsilon_d)$ be the opposite of ϵ . We define a family of 2^{d-1} disjoint subsets of D indexed by E by setting

$$\forall \epsilon \in E, \quad D_{\epsilon} = \{x \in D : \epsilon_2 x_2, \dots, \epsilon_d x_d > 0\} .$$

Let H be the hyperplane $\{x_1 = 1\}$ and, for all $x \in \mathbb{R}^d$, let B(x) be the open ball with center at x and radius 1.

Lemma 4.5. If $x \in D_{\epsilon} \cap H$, then $B(x) \cap B^c \cap D_{\overline{\epsilon}} = \emptyset$.

Proof. Let $x \in D_{\epsilon} \cap H$. For any $y \in B(x) \cap D_{\overline{\epsilon}}$ we have

$$1 > (y_1 - 1)^2 + (y_2 - x_2)^2 + \dots + (y_d - x_d)^2$$

$$= (y_1 - 1)^2 + y_2^2 + \dots + y_d^2 + \sum_{i=2}^d (x_i^2 - 2y_i x_i)$$

$$> (y_1 - 1)^2 + y_2^2 + \dots + y_d^2$$

since $x \in D_{\epsilon}$ and $y \in D_{\overline{\epsilon}}$ imply $-x_i y_i > 0$ for all $i = 2 \dots d$. Hence y belongs to B.

We now come to the estimate which is the heart of this section:

Lemma 4.6.

$$\lim_{s\to 0} \limsup_{\lambda\to 0} \widetilde{W}_{\lambda e_1,1}^D(\tau_B \le s) = 0.$$

Proof. We will show that

(10)
$$\limsup_{\lambda \to 0} \widetilde{W}_{\lambda e_1, 1}^D(\tau_B \le s) \le 2^{d-1} \, \widetilde{W}_{0, 1}^D(\tau_B \le s)$$

for all s > 0, and the lemma will then follow by letting $s \to 0$ since $\widetilde{W}_{0,1}^D(\tau_B = 0) = 0$.

For $\lambda > 0$, set

$$T_{\lambda} = \inf\{t \geq 0 : X_1(t) = \lambda \text{ and } X_1(s) > 0, \forall s \in (t, t+1]\}$$

and consider the process Z_{λ} defined by

$$\forall t \geq 0, \quad Z_{\lambda}(t) = X(T_{\lambda} + t) - X(T_{\lambda}) + \lambda e_1.$$

Remember that Z_{λ} has the distribution $\widetilde{W}_{\lambda e_1,1}^D$ with respect to W_0 (see Lemma 3.6).

Write

(11)
$$W_0(\tau_B(Z_\lambda) \le s) \le W_0(\tau_B(Z_0) \le s + T_\lambda - T_0) + W_0(\tau_B(Z_0) > s + T_\lambda - T_0; \tau_B(Z_\lambda) \le s)$$
.

For convenience, we set $u = \tau_B(Z_\lambda)$. If $\tau_B(Z_0) > s + T_\lambda - T_0$ and $u \le s$, then $Z_0(u + T_\lambda - T_0) = X(T_\lambda + u) - X(T_0)$ belongs to B; this means that

 $Z_{\lambda}(u)$ belongs to $B(Y_{\lambda})$, where we have put

$$Y_{\lambda} = X(T_0) - X(T_{\lambda}) + (1 + \lambda)e_1$$

= $(1, X_2(T_0) - X_2(T_{\lambda}), \dots, X_d(T_0) - X_d(T_{\lambda}))$.

Note that $Y_{\lambda} \in H$. Since $Z_{\lambda}(u) \notin B$, we see by Lemma 4.5 that $Z_{\lambda}(u) \notin D_{\overline{\epsilon}}$ as soon as $Y_{\lambda} \in D_{\epsilon}$. Therefore

$$W_0(\tau_B(Z_0) > s + T_\lambda - T_0; \tau_B(Z_\lambda) \le s)$$

$$\le \sum_{\epsilon \in E} W_0(Y_\lambda \in D_\epsilon; \tau_B(Z_\lambda) \le s; Z_\lambda(u) \notin D_{\overline{\epsilon}}) .$$

Now, it is easily seen that Y_{λ} is independent of Z_{λ} conditionally to X_1 . In addition, we have $W_0(Y_{\lambda} \in D_{\epsilon}|X_1) = 1/2^{d-1}$. Thus

$$W_0(\tau_B(Z_0) > s + T_\lambda - T_0; \tau_B(Z_\lambda) \le s)$$

$$\le \frac{1}{2^{d-1}} \sum_{\epsilon \in E} W_0(\tau_B(Z_\lambda) \le s; Z_\lambda(u) \notin D_{\overline{\epsilon}})$$

$$= \frac{2^{d-1} - 1}{2^{d-1}} W_0(\tau_B(Z_\lambda) \le s) .$$

Combining this inequality with equation (11) gives

$$W_0(\tau_B(Z_\lambda) \le s) \le 2^{d-1} W_0(\tau_B(Z_0) \le s + T_\lambda - T_0)$$

and the result follows by letting $\lambda \to 0$ since $\lim_{\lambda \to 0} T_{\lambda} = T_0$ almost surely (see Lemma 3.3).

Remark 4.7. Note that the proof of Lemma 4.6 does not involve the "size" of B. Hence the result holds if B is replaced by any open ball B' that is tangent to ∂D at 0.

The following proposition and its proof illustrate how the ball estimate can be combined to invariance properties of Brownian motion in order to solve the convergence problem.

Proposition 4.8. As $x \in B \to 0$, $\widetilde{W}_{x,1}^B \Rightarrow \widetilde{W}_{0,1}^B$.

Proof. It follows from Lemmas 4.2 and 4.6 that

$$\widetilde{W}_{\lambda e_1,1}^B \Rightarrow \widetilde{W}_{0,1}^B \text{ as } \lambda \to 0$$
.

Let $(x_n) \in B$ be a sequence such that $x_n \to 0$ and set $\lambda_n = 1 - ||x_n - e_1||$. For each n, choose a rotation R_n in the plane $(0, e_1, x_n)$ with center at e_1 such that $R_n(\lambda_n e_1) = x_n$. From the invariance properties of Brownian motion, it is easily seen that

$$\widetilde{W}_{x_n,1}^B = \widetilde{W}_{\lambda_n e_1,1}^B \circ R_n^{-1} \; .$$

Now, by elementary geometry,

$$||R_n(y) - y|| \le ||y - e_1|| \frac{||x_n - \lambda_n e_1||}{||x_n - e_1||}$$

for all $y \in \mathbb{R}^d$. Thus, for any $w \in \mathcal{C}_{\infty}$ and any sequence $(w_n) \in \mathcal{C}_{\infty}$ such that $w_n \to w$, we get

$$R_n(w_n) \to w$$
.

It therefore follows from the continuous mapping theorem that

$$\widetilde{W}_{x_n,1}^B = \widetilde{W}_{\lambda_n e_1,1}^B \circ R_n^{-1} \Rightarrow \widetilde{W}_{0,1}^B \ .$$

4.4. **Application to nice sets.** In this section, we introduce the notion of *nice* sets and solve the convergence problem for those sets.

4.4.1. Convergence with variable sets. For any set $U \subset \mathbb{R}^d$ and any $\epsilon > 0$, put

$$U_{\epsilon^{+}} = \{ x \in \mathbb{R}^{d} : d(x, U) \le \epsilon \}$$

and

$$U_{\epsilon^{-}} = \{ x \in U : d(x, U^c) \ge \epsilon \} .$$

If (U_n) is a sequence of subsets of \mathbb{R}^d , we will say that (U_n) converges to U and write $U_n \to U$ if for all $\epsilon > 0$ there exist a n_0 such that

$$n \ge n_0 \Rightarrow U_{\epsilon^-} \subset U_n \subset U_{\epsilon^+}$$
.

Let D be the half-space $\{x \in \mathbb{R}^d : x_1 > 0\}$ and let B be an open ball tangent to ∂D at 0. Set $e_1 = (1, 0, \dots, 0)$.

Proposition 4.9. Let U be an open co-regular set such that $B \subset U \subset D$ and let (U_n) be a sequence of sets such that:

- (1) For all $n, B \subset U_n \subset D$;
- (2) For all R > 0, $U_n \cap B(0, R) \to U \cap B(0, R)$.

Then, for all sequence (λ_n) of positive numbers converging to 0,

$$\widetilde{W}_{\lambda_n e_1, 1}^{U_n} \Rightarrow \widetilde{W}_{0, 1}^{U}$$
.

Proof. Since $\widetilde{W}_{0,1}^D(\tau_U > 1) \ge \widetilde{W}_{0,1}^D(\tau_B > 1) > 0$, it suffices to prove that

(12)
$$\lim_{n \to \infty} \widetilde{W}_{\lambda_n e_1, 1}^D(f; \tau_{U_n} > 1) = \widetilde{W}_{0, 1}^D(f; \tau_U > 1)$$

for all bounded continuous function $f: C_{\infty} \to \mathbb{R}$.

First, suppose that $U_n \to U$. Since each set U_n contains the ball B, we have

$$\widetilde{W}_{\lambda_n e_1, 1}^D(\tau_{U_n} \le s) \le \widetilde{W}_{\lambda_n e_1, 1}^D(\tau_B \le s)$$
.

Hence

(13)
$$\lim_{s \to 0} \limsup_{n} \widetilde{W}_{\lambda_n e_1, 1}^D(\tau_{U_n} \le s) = 0$$

by Lemma 4.6. Now, a minor modification of the proof of Lemma 4.2 shows that (12) holds. Let us give some details. Because of (13) it is sufficient to prove that

(14)
$$\lim_{n \to \infty} \widetilde{W}_{\lambda_n e_1, 1}^D(f; \tau_{U_n}^s > 1) = \widetilde{W}_{0, 1}^D(f; \tau_U^s > 1)$$

for all s > 0, where $\tau_{U_n}^s = \inf\{t \geq s : X_t \notin U_n\}$ and $\tau_U^s = \inf\{t \geq s : X_t \notin U_n\}$. Recall that $\Omega = \{\tau_U^s = \tau_{\overline{U}}^s \neq 1\}$ satisfies $\widetilde{W}_{0,1}^D(\Omega) = 1$. Let $w \in \Omega$ be given and let $(w_n) \in \mathcal{C}_{\infty}$ be a sequence such that $w_n \to w$.

If $\tau_U^s(w) > 1$, then $w(t) \in U$ for all $t \in [s, 1]$. Choose $\epsilon > 0$ such that $w(t) \in U_{2\epsilon^-}$ for all $t \in [s, 1]$. Then, for n large enough, $w_n(t) \in U_{\epsilon^-}$ for all $t \in [s, 1]$ and $U_n \supset U_{\epsilon^-}$; thus $\tau_{U_n}^s(w_n) > 1$.

If $\tau_U^s(w) = \tau_{\overline{U}}^s(w) < 1$, then there exists $t \in [s,1)$ such that $w(t) \notin \overline{U}$. Pick $\epsilon > 0$ such that $w(t) \notin U_{2\epsilon^+}$. Then, for n sufficiently large, $w_n(t) \notin U_{\epsilon^+}$ and $U_n \subset U_{\epsilon^+}$; hence $\tau_{U_n}^s(w_n) < 1$. This proves that in each case

$$1_{\{\tau_{U}^{s}>1\}}(w_{n}) \to 1_{\{\tau_{U}^{s}>1\}}(w)$$

and (14) follows from the continuous mapping theorem.

Now we turn to the general case, that is we consider the local convergence hypothesis 2. of Proposition 4.9. Fix $\epsilon > 0$ and choose R > 0 such that $\widetilde{W}_{0,1}^D(\tau_{B(0,R)} > 1) \geq 1 - \epsilon$. By Lemma 4.1, we have

$$\lim_{n \to \infty} \widetilde{W}_{\lambda_n e_1, 1}^D(\tau_{B(0,R)} > 1) = \widetilde{W}_{0,1}^D(\tau_{B(0,R)} > 1) .$$

Therefore

$$\widetilde{W}_{\lambda_n e_1, 1}^D(\tau_{B(0,R)} > 1) \ge 1 - 2\epsilon$$

for all large enough n. Set $U'_n = U_n \cap B(0,R)$ and $U' = U \cap B(0,R)$. Then

$$\begin{aligned} |\widetilde{W}_{\lambda_{n}e_{1},1}^{D}(f;\tau_{U_{n}}>1) - \widetilde{W}_{0,1}^{D}(f;\tau_{U}>1)| \\ &\leq |\widetilde{W}_{\lambda_{n}e_{1},1}^{D}(f;\tau_{U_{n}'}>1) - \widetilde{W}_{0,1}^{D}(f;\tau_{U'}>1)| + 3M\epsilon \end{aligned}$$

where M is a bound for |f|. By hypothesis $U'_n \to U'$, hence

$$\limsup_n |\widetilde{W}^D_{\lambda_n e_1,1}(f;\tau_{U'_n}>1) - \widetilde{W}^D_{0,1}(f;\tau_{U'}>1)| = 0$$

by the first step of this proof. Therefore

$$\limsup_{n} |\widetilde{W}_{\lambda_{n}e_{1},1}^{D}(f;\tau_{U_{n}}>1) - \widetilde{W}_{0,1}^{D}(f;\tau_{U}>1)| \le 4B\epsilon$$

and the desired result follows by letting $\epsilon \to 0$.

- 4.4.2. Nice sets. Let U be an open subset of \mathbb{R}^d and x_0 a boundary point of U. We will say that U is nice at x_0 if there exist a neighborhood V of x_0 and a number r > 0 such that the following conditions are satisfied:
 - (1) For all $x \in \partial U \cap V$ there exists a half-space $D_x \supset U$ such that:
 - $x \in \partial D_x$;
 - The ball $B_x \subset D_x$ with radius r which is tangent to ∂D_x at x is contained in U;

- The application c which maps x to the center c(x) of the ball B_x is continuous at x_0 .
- (2) For all $y \in U \cap V$ such that $d(y, \partial U) \leq r/2$, there exists a point $x = p(y) \in \partial U \cap V$ such that:
 - $y \in (x, c(x)];$
 - The mapping $y \mapsto p(y)$ is continuous.

Remark 4.10. It can be verified that regularity and convexity assumptions ensures the property of being a "nice set". If the open set U is convex and has a boundary of class C^2 in a neighborhood of x_0 then the set U is nice at x_0 .

Suppose that U is nice at x_0 . With the above notations, for any $x \in \partial U \cap V$, the point $x + c(x_0) - c(x)$ belongs to the boundary of the ball B_{x_0} ; thus we can choose a planar rotation R_x with center at $c(x_0)$ and such that $R_x(x+c(x_0)-c(x))=x_0$. Note that the angle of R_x tends to 0 as $x \to x_0$, since $c(x) \to c(x_0)$. Set $\phi_x(y) = R_x(y+c(x_0)-c(x))$ and $U_x = \phi_x(U)$. Then it can be seen that

$$(15) B_{x_0} \subset U_x \subset D_{x_0}$$

and

(16)
$$U_x \cap B(0,R) \to U_{x_0} \cap B(0,R), \text{ as } x \to x_0,$$

for all R > 0.

Theorem 4.11. Suppose that U is co-regular and nice at x_0 . Then, as $x \in U \to x_0$, we have

$$\widetilde{W}_{x,1}^{U} \Rightarrow \widetilde{W}_{x_0,1}^{U}$$
.

Proof. For y close to x_0 , set x = p(y). Since y belongs to (x, c(x)], the point $q(y) = \phi_x(y)$ belongs to $(x_0, c(x_0)]$. Moreover, q(y) tends to x_0 as $y \to x_0$. Thus, from (15) and (16) together with Proposition 4.9, we obtain

$$\widetilde{W}_{q(y),1}^{U_x} \Rightarrow \widetilde{W}_{x_0,1}^{U}$$
, as $y \to x_0$.

Now by the invariance properties of Brownian motion, we have

$$\widetilde{W}_{y,1}^U = \widetilde{W}_{q(y),1}^{U_x} \circ \phi_x$$
.

Since ϕ_x tends to the identity mapping as $x \to x_0$, uniformly on compact subsets of \mathbb{R}^d , it follows from the continuous mapping theorem that

$$\widetilde{W}_{y,1}^U \Rightarrow \widetilde{W}_{x_0,1}^U$$
, as $y \to x_0$.

5. Nice cones

Let $d \geq 2$ and let $C \subset \mathbb{R}^d$ be an open cone with vertex at 0. We will say that C is a *nice cone* if it is nice (see 4.4.2) at any point of its boundary, excepting 0. For example, any two-dimensional convex cone is nice. In higher dimension, any circular cone or ellipsoidal cone is nice.

We note two important facts about nice cones:

- (1) If C is a nice cone, then it is a Lipschitz cone;
- (2) If C is a nice cone, ∂C is a null set with respect to Lebesgue measure.

The proof of the first one is elementary but quite tedious, so we omit it here. Note that the second fact is a consequence of the first one.

5.1. **Proof of Theorem 1.1.** We begin with the following lemma:

Lemma 5.1. Suppose C is a nice cone. Let $x_0 \in \partial C \setminus \{0\}$ and $t_0 > 0$. As $(x,t) \to (x_0,t_0),$

$$\widetilde{W}_{x,t}^C \Rightarrow \widetilde{W}_{x_0/\sqrt{t_0},1}^C \circ K_{t_0}^{-1}$$
.

Proof. By the scaling property of \widetilde{W}^C (see Remark 2.6), we have

$$\widetilde{W}_{x,t}^C = \widetilde{W}_{x/\sqrt{t},1}^C \circ K_t^{-1}$$
.

The result simply follows from Theorem 4.11 together with the continuous mapping theorem.

5.1.1. Convergence of the finite-dimensional distributions. We will prove in this section that the finite-dimensional distributions of $\widetilde{W}_{x,1}^C$ converge weakly as $x \in C \to 0$. Recall that for any $t \in (0,1]$ the law $\widetilde{W}_{x,1}^C(X_t \in dy)$ has the density $e_x(t,y)$ given by

$$e_x(t,y) = \frac{p^C(t,x,y)}{W_x(\tau_C > 1)} W_y(\tau_C > 1 - t) .$$

By using an expansion of the heat kernel $p^{C}(t,x,y)$ of C that is given by R. Bañuelos and R. Smits in [1], we shall prove that $e_x(t,y)$ converges to a limit density e(t, y), as $x \in C \to 0$.

Before we recall their result, let us introduce some notations. Let C_{Σ} be the intersection of the cone C with the unit sphere \mathbb{S}^{d-1} and suppose that it is a regular set for the Dirichlet problem with respect to the Laplace-Beltrami operator L on \mathbb{S}^{d-1} . Then there exists a complete set of orthonormal eigenfunctions m_j with corresponding eigenvalues $0 < \lambda_1 < \lambda_2 \le \lambda_3 < \cdots$ satisfaying

$$\begin{cases} Lm_j = -\lambda_j m_j & \text{on } C_{\Sigma}; \\ m_j = 0 & \text{on } \partial C_{\Sigma} \end{cases}$$

 $\begin{cases} Lm_j = -\lambda_j m_j & \text{on } C_\Sigma \;; \\ m_j = 0 & \text{on } \partial C_\Sigma \;. \end{cases}$ We set $\alpha_j = \sqrt{\lambda_j + (\frac{d}{2} - 1)^2}$. We will use the following facts that are proved in [1]:

• there exists two constants $0 < c_1 < c_2$ such that

(17)
$$\forall j \ge 1, \quad c_1 j^{\frac{1}{d-1}} \le \alpha_j \le c_2 j^{\frac{1}{d-1}};$$

 \bullet there exists a constant c such that

(18)
$$\forall j \ge 1, \quad ||m_j||_{\infty} \le c\alpha_j^{\frac{d-1}{2}};$$

ullet if C is a Lipschitz cone, then there exists a constant c' such that

(19)
$$\forall j \ge 1, \forall \eta \in C_{\Sigma}, \quad m_j^2(\eta) \le \frac{c' m_1^2(\eta)}{I_{\alpha_j}(1)},$$

where I_{ν} is the modified Bessel function of order ν :

(20)
$$I_{\nu}(x) = \frac{2(\frac{x}{2})^{\nu}}{\sqrt{\pi}\Gamma(\nu + \frac{1}{2})} \int_{0}^{\frac{\pi}{2}} (\sin t)^{2\nu} \cosh(x \cos t) dt$$
$$= \sum_{m=0}^{\infty} \frac{(\frac{x}{2})^{\nu + 2m}}{m!\Gamma(\nu + m + 1)}.$$

Then we have the following lemma:

Proposition 5.2 (Bañuelos, Smits). Write $x = \rho\theta$, $y = r\eta$, ρ , r > 0, θ , $\eta \in C_{\Sigma}$. We have

$$p^{C}(t, x, y) = \frac{e^{-\frac{(r^{2} + \rho^{2})}{2t}}}{t(\rho r)^{\frac{d}{2} - 1}} \sum_{j=1}^{\infty} I_{\alpha_{j}}\left(\frac{\rho r}{t}\right) m_{j}(\theta) m_{j}(\eta) .$$

Together with the expression of I_{α_j} , this suggests that $p^C(t, x, y)$ is equivalent at x = 0 to the product g(x)h(t, y) where

$$g(x) = \rho^{\alpha_1 - (\frac{d}{2} - 1)} m_1(\theta)$$

and

$$h(t,y) = \frac{r^{\alpha_1 - (\frac{d}{2} - 1)} e^{-\frac{r^2}{2t}}}{2^{\alpha_1} \Gamma(\alpha_1 + 1) t^{\alpha_1 + 1}} m_1(\eta) .$$

In fact, we have the following:

Lemma 5.3. For $x = \rho\theta$, $y = r\eta$, ρ , r > 0, θ , $\eta \in C_{\Sigma}$, we have

$$\lim_{\rho \to 0} \frac{p^C(t, x, y)}{g(x)} = h(t, y) ,$$

uniformly in $(t, r, \theta, \eta) \in [T, +\infty) \times [0, R] \times C_{\Sigma} \times C_{\Sigma}$, for any positive constants T and R.

Proof. Throughout this proof, the letter κ will denote some positive constant whose value may change from line to line.

Set $M = \frac{\rho r}{t}$. We have

$$\frac{p^{C}(t,x,y)}{g(x)h(t,y)} = 2^{\alpha_1}\Gamma(\alpha_1+1)e^{-\frac{\rho^2}{2t}}\sum_{j=1}^{\infty} \frac{I_{\alpha_j}(M)}{M^{\alpha_1}} \frac{m_j(\theta)}{m_1(\theta)} \frac{m_j(\eta)}{m_1(\eta)} .$$

Using relation (19), we get

$$\left| \frac{I_{\alpha_j}(M)}{M^{\alpha_1}} \frac{m_j(\theta)}{m_1(\theta)} \frac{m_j(\eta)}{m_1(\eta)} \right| \le \frac{\kappa}{M^{\alpha_1}} \frac{I_{\alpha_j}(M)}{I_{\alpha_j}(1)} .$$

Now, using the integral expression for I_{α_i} , we obtain

$$I_{\alpha_j}(M) \le \frac{2(\frac{M}{2})^{\alpha_j}}{\sqrt{\pi}\Gamma(\alpha_j + \frac{1}{2})} \cosh(M) \int_0^{\frac{\pi}{2}} (\sin t)^{2\alpha_j} dt$$

and

$$I_{\alpha_j}(1) \ge \frac{2(\frac{1}{2})^{\alpha_j}}{\sqrt{\pi}\Gamma(\alpha_j + \frac{1}{2})} \int_0^{\frac{\pi}{2}} (\sin t)^{2\alpha_j} dt$$
.

Hence

$$\frac{I_{\alpha_j}(M)}{I_{\alpha_j}(1)} \le M^{\alpha_j} \cosh M ,$$

and so we get

$$\left| \frac{I_{\alpha_j}(M)}{M^{\alpha_1}} \frac{m_j(\theta)}{m_1(\theta)} \frac{m_j(\eta)}{m_1(\eta)} \right| \le \kappa M^{\alpha_j - \alpha_1} \cosh M.$$

From relation (17), it is easily seen that the series $\sum_{j} M^{\alpha_{j}-\alpha_{1}} \cosh M$ is uniformly convergent on $[0, 1-\epsilon]$. So, the series

$$\sum_{j} \frac{I_{\alpha_{j}}(M)}{M^{\alpha_{1}}} \frac{m_{j}(\theta)}{m_{1}(\theta)} \frac{m_{j}(\eta)}{m_{1}(\eta)}$$

is uniformly convergent for $(M, \theta, \eta) \in [0, 1 - \epsilon] \times C_{\Sigma} \times C_{\Sigma}$. Therefore we can take the limit term by term : since

$$\lim_{M \to 0} \frac{I_{\alpha_j}(M)}{M^{\alpha_1}} \frac{m_j(\theta)}{m_1(\theta)} \frac{m_j(\eta)}{m_1(\eta)} = \begin{cases} \frac{1}{2^{\alpha_1} \Gamma(\alpha_1 + 1)} & \text{si } j = 1 ; \\ 0 & \text{si } j \ge 2 ; \end{cases}$$

we get

$$\lim_{M\to 0}\sum_{j=1}^{\infty}\frac{I_{\alpha_j}(M)}{M^{\alpha_1}}\frac{m_j(\theta)}{m_1(\theta)}\frac{m_j(\eta)}{m_1(\eta)}=\frac{1}{2^{\alpha_1}\Gamma(\alpha_1+1)}\;,$$

where the convergence is uniform for $(\theta, \eta) \in C_{\Sigma} \times C_{\Sigma}$.

Lemma 5.4. The function of y

$$\sup_{\|x\|<\frac{1}{2}} \left| \frac{p^C(1,x,y)}{g(x)} \right|$$

is integrable on \mathbb{R}^d .

Proof. Throughout this proof, the letter κ will denote some positive constant whose value may change from line to line. Using relations (18) and (19), we get

$$\left| \frac{p^{C}(1, x, y)}{g(x)} \right| \leq \frac{e^{-\frac{(\rho^{2} + r^{2})}{2}}}{r^{\frac{d}{2} - 1} \rho^{\alpha_{1}}} \sum_{j=1}^{\infty} I_{\alpha_{j}}(\rho r) \left| \frac{m_{j}(\theta)}{m_{1}(\theta)} m_{j}(\eta) \right|$$

$$\leq \frac{e^{-\frac{r^{2}}{2}}}{r^{\frac{d}{2} - 1} \rho^{\alpha_{1}}} \sum_{j=1}^{\infty} \frac{I_{\alpha_{j}}(\rho r)}{I_{\alpha_{j}}(1)^{\frac{1}{2}}} \alpha_{j}^{\frac{d-1}{2}}.$$

From the integral expression for I_{α_i} , we find that

$$I_{\alpha_j}(\rho r) \le \frac{2(\frac{\rho r}{2})^{\alpha_j}}{\sqrt{\pi}\Gamma(\alpha_j + \frac{1}{2})} \cosh(\rho r) \int_0^{\frac{\pi}{2}} (\sin t)^{2\alpha_j} dt$$

and

$$I_{\alpha_j}(1) \ge \frac{2(\frac{1}{2})^{\alpha_j}}{\sqrt{\pi}\Gamma(\alpha_j + \frac{1}{2})} \int_0^{\frac{\pi}{2}} (\sin t)^{2\alpha_j} dt$$
.

So, if we set $\omega_{\alpha_i} = \int_0^{\frac{\pi}{2}} (\sin t)^{2\alpha_j} dt$, we get

$$\frac{I_{\alpha_j}(\rho r)}{I_{\alpha_j}(1)^{\frac{1}{2}}} \le \kappa \cosh(\rho r) \left(\frac{\rho r}{\sqrt{2}}\right)^{\alpha_j} \frac{\sqrt{\omega_{\alpha_j}}}{\Gamma(\alpha_j + \frac{1}{2})^{\frac{1}{2}}}.$$

The Wallis integral $\int_0^{\frac{\pi}{2}} (\sin t)^{2n} dt$ is equivalent to $cn^{-\frac{1}{2}}$, thus ω_{α_j} is equivalent to $c\alpha_j^{-\frac{1}{2}}$. From Stirling's Formula we also get $\Gamma(\alpha_j + \frac{1}{2}) \geq c\alpha_j^{\alpha_j} e^{-\alpha_j}$. Thus,

$$\frac{I_{\alpha_j}(\rho r)}{I_{\alpha_j}(1)^{\frac{1}{2}}} \le \kappa \cosh(\rho r) \left(\frac{\sqrt{e}\rho r}{\sqrt{2}}\right)^{\alpha_j} \frac{\alpha_j^{-\frac{1}{4}}}{\alpha_i^{\alpha_j/2}}.$$

Therefore,

$$(21) \qquad \left| \frac{p^C(1, x, y)}{g(x)} \right| \leq \kappa \frac{e^{-\frac{r^2}{2}}}{r^{\frac{d}{2} - 1} \rho^{\alpha_1}} \cosh(\rho r) \sum_{j=1}^{\infty} \left(\frac{\sqrt{e} \rho r}{\sqrt{2}} \right)^{\alpha_j} \frac{\alpha_j^{\frac{2d-3}{4}}}{\alpha_j^{\alpha_j/2}}.$$

Since $\alpha_j \geq \alpha_1$, the right-hand side of (21) is increasing with ρ , so

$$\sup_{\rho \leq \frac{1}{2}} \left| \frac{p^C(1,x,y)}{g(x)} \right| \leq \kappa \frac{e^{-\frac{r^2}{2}}}{r^{\frac{d}{2}-1}} \cosh\left(\frac{r}{2}\right) \sum_{j=1}^{\infty} \left(\frac{\sqrt{e}r}{2\sqrt{2}}\right)^{\alpha_j} \frac{\alpha_j^{\frac{2d-3}{4}}}{\alpha_j^{\alpha_j/2}} =: f(r) \ .$$

Because $\alpha_j > \left(\frac{d}{2} - 1\right)$, the function f is integrable on any compact subset of $[0, +\infty)$. We shall now find an upper bound for the sum that appears in the definition of f for large values of r. Let $M \ge 1$. For $2n \le \alpha_j \le 2n + 1$,

we have

$$M^{\alpha_j} \frac{\alpha_j^{\frac{2d-3}{4}}}{\alpha_j^{\alpha_j/2}} \le M^{2n+1} \frac{(2n+1)^{\frac{2d-3}{4}}}{(2n)^n} = M(M^2/2)^n \frac{(2n+1)^{\frac{2d-3}{4}}}{n^n} .$$

Since $\alpha_j > c_1 j^{\frac{1}{d-1}}$, the number of indices j for which $\alpha_j \leq 2n+1$ is bounded by $\left(\frac{2n+1}{c_1}\right)^{d-1}$. Thus, there exist K = K(d) > 0 such that

$$\sum_{j=1}^{\infty} M^{\alpha_j} \frac{\alpha_j^{\frac{2d-3}{4}}}{\alpha_j^{\alpha_j/2}} \le M \sum_{n=1}^{\infty} (M^2/2)^n \frac{n^K}{n^n} \le P(M) e^{M^2/2} ,$$

where P is a polynomial. Applying this result with $M=\frac{\sqrt{er}}{2\sqrt{2}}$ and $r\geq 2\sqrt{2/e}$ gives

$$f(r) \le \kappa \frac{e^{-(1-e/8)\frac{r^2}{2}}}{r^{\frac{d}{2}-1}}\widetilde{P}(r)\cosh\left(\frac{r}{2}\right)$$

where \widetilde{P} is a polynomial whose coefficients depend only on d. This is sufficient to conclude the proof of Lemma 5.4.

Proposition 5.5. For any $t \in [0,1]$ and any R > 0, we have

$$\lim_{\|x\|\to 0} e_x(t,y) = \frac{h(t,y)}{\int h(1,z) \, dz} W_y(\tau_C > 1 - t) ,$$

uniformly on $\{y \in C : ||y|| \le R\}$.

Proof. First recall that $W_x(\tau_C > 1) = \int p^C(1, x, z) dz$. It follows from Lemmas 5.3, 5.4 and the dominated convergence theorem that

(22)
$$\lim_{\|x\|\to 0} \frac{W_x(\tau_C > 1)}{g(x)} = \int h(1, z) \, dz .$$

Since

$$e_x(t,y) = \frac{p^C(t,x,y)}{g(x)} \frac{g(x)}{W_x(\tau_C > 1)} W_y(\tau_C > 1 - t) ,$$

the desired result is easily deduced from proposition 5.3 and relation (22).

Remark 5.6. By scaling, it is easily deduced from (22) that

$$\lim_{t \to +\infty} \frac{W_x(\tau_C > t)}{g(t^{-1/2}x)} = \int h(1, z) \, dz \; ,$$

uniformly on $\{x \in C : ||x|| \leq R\}$, for any positive constant R. So, for a Lipschitz cone, the convergence is stronger than what was announced by R. Bañuelos and R. Smits ([1], Corollary 1).

Note that a straightforward computation gives

$$\int h(1,z) dz = 2^{-\frac{\alpha_1}{2} + \frac{d-2}{4}} \frac{\Gamma(\frac{\alpha_1}{2} + \frac{d+2}{4})}{\Gamma(\alpha_1 + 1)} \int_{C_{\Sigma}} m_1(\eta) \, \sigma(d\eta) ,$$

where σ is Lebesgue measure on the unit sphere \mathbb{S}^{d-1} .

Let e(t,y) be the limit of the densities $e_x(t,y)$ as $x \in C \to 0$. By Proposition 5.5, we have

(23)
$$e(t,y) = ct^{-\alpha_1 - 1} ||y||^{\alpha_1 - (\frac{d}{2} - 1)} e^{-\frac{||y||^2}{2t}} m_1(\vec{y}) W_y(\tau_C > 1 - t) ,$$

where $c^{-1} = 2^{\frac{\alpha_1}{2} + \frac{d-2}{4}} \Gamma(\frac{\alpha_1}{2} + \frac{d+2}{4}) \int m_1(\eta) \, \sigma(d\eta)$ and $\vec{y} = y/\|y\|$. Proposition 5.5 is not sufficient to prove weak convergence of the law

$$\widetilde{W}_{x,1}^C(X_t \in dy) = e_x(t,y) \, dy$$

as $x \in C$ goes to 0. Indeed, except for t = 1, we don't know a priori if e(t, y)is a probability density (except for the case t=1), and this seems hard to check by a direct computation. However, when C is a nice cone, we may handle this by proving the equi-integrability of the family $\{e_x(t,y): ||x|| \le$ 1}.

Lemma 5.7. For all $t \in [0,1)$, we have

$$\lim_{R \to \infty} \sup_{\|x\| < 1} \widetilde{W}_{x,1}^C(\|X_t\| > R) = 0.$$

In other words, the family of probability $\{\widetilde{W}_{x_1}^C(X_t \in dy) : x \in C, ||x|| \leq 1\}$ is tight.

Proof. Let $x \in C$ with $||x|| \le 1$ be given, and let R > 2. We denote by $\rho = \tau_{B(0,2)}$ the first exit time from the ball B(0,2). A continuous path started at x that is outside B(0,R) at time t must have left B(0,2) before that time, so

$$\begin{split} \widetilde{W}_{x,1}^{C}(\|X_{t}\| > R) \\ &= \widetilde{W}_{x,1}^{C}(\rho < t; \|X_{t}\| > R) \\ &= \widetilde{W}_{x,1}^{C}(\rho < t; \widetilde{W}_{X_{\rho},1-\rho}^{C}(\|X_{t-s}\| > R)_{|s=\rho}) \\ &\leq \sup \left\{ \widetilde{W}_{y,1-s}^{C}(\|X_{t-s}\| > R) : y \in C, \|y\| = 2 \text{ and } s \in [0,t] \right\} . \end{split}$$

Suppose that the last expression does not tend to 0 as $R \to \infty$; then there exist a sequence $(y_n) \in C$ with $||y_n|| = 2$ and a sequence $(s_n) \in [0,t]$ such that

(24)
$$\liminf_{n \to \infty} \widetilde{W}_{y_n, 1 - s_n}^C(\|X_{t - s_n}\| > n) > 0.$$

Without loss of generality, we can suppose that (y_n) converges to a point $y \in \overline{C}$ with ||y|| = 2, and that (s_n) converges to $s \in [0, t]$. But Lemma 5.1 (or Proposition 2.3 if $y \in C$) then implies that $\left(\widetilde{W}_{y_n,1-s_n}^C(X_{t-s_n} \in dy)\right)$ is a convergent sequence of probability measures: this contradicts (24).

Proposition 5.8. The finite-dimensional distributions of $\widetilde{W}_{x,1}^C$ converge weakly as $x \in C$ tends to 0.

Moreover, the limit distribution of the first transition law $W_{x,1}^C(X_t \in dy)$, $t \in (0,1]$, has density e(t,y) w.r.t. Lebesgue measure.

Proof. By Proposition 5.5, $e_x(t, y)$ converges to e(t, y) uniformly on every compact set as $x \in C$ tends to 0. Thus, for all R > 0, we have

$$\limsup_{\|x\| \to 0} \int |e_x(t, y) - e(t, y)| dy$$

$$\leq \sup_{\|x\| \leq 1} \int_{B(0, R)^c} e_x(t, y) dy + \int_{B(0, R)^c} e(t, y) dy.$$

Thanks to Lemma 5.7 and the integrability of e(t, y), letting $R \to +\infty$ gives

$$\limsup_{\|x\| \to 0} \int |e_x(t, y) - e(t, y)| \ dy = 0 \ .$$

This proves that the function $y \mapsto e(t,y)$ is a probability density and that the probability measures $\widetilde{W}_{x,1}^C(X_t \in dy)$ converge weakly to e(t,y) dy as $x \in C$ tends to 0. The weak convergence of the finite-dimensional distributions then follows from Proposition 2.4 since ∂C has Lebesgue measure 0.

5.1.2. Tightness. For any T > 0, the space \mathcal{C}_T of continuous paths $w : [0,T] \to \mathbb{R}^d$ is endowed with the topology generated by the supremum metric and the corresponding Borel σ -algebra.

Proposition 5.9. For any sequence (x_n) of points of C converging to 0 and for any T > 0, the sequence of probability measures $(\widetilde{W}_{x_n,1}^C)$ is tight in C_T .

Proof. Since the arguments do not depend on the value of T, we will only consider the case T=1. It suffices to prove that, for all $\epsilon > 0$,

$$\lim_{\delta \to 0} \limsup_{n \to \infty} \widetilde{W}^{C}_{x_n,1}(\chi(\delta,0,1) > \epsilon) = 0 \; ,$$

where

$$\chi(\delta, a, b)(w) = \sup\{\|w(s) - w(t)\| : |s - t| < \delta, s, t \in [a, b]\}$$

is the modulus of continuity of order δ of w on [a, b] (see Billingsley [2]).

Our proof is a modification of Shimura's one for the two-dimensional case ([8], Theorem 2).

Fix $\epsilon > 0$ and set s = 1/2. Since $\chi(\delta, \cdot, \cdot)$ is sub-additive when considered as a function on the set of intervals, we have

$$\widetilde{W}_{x_{n},1}^{C}(\chi(\delta,0,1) > 4\epsilon) \leq \underbrace{\widetilde{W}_{x_{n},1}^{C}\left(\chi(\delta,0,s) > 3\epsilon\right)}_{A_{n}(\delta)} + \underbrace{\widetilde{W}_{x_{n},1}^{C}\left(\chi(\delta,s,1) > \epsilon\right)}_{B_{n}(\delta)}.$$

Let us start with $B_n(\delta)$. It follows from Proposition 5.8 that

$$\lim_{r \to 0, R \to \infty} \liminf_{n \to \infty} \widetilde{W}_{x_n, 1}^{C} \left(r \le ||X_s|| \le R \right) = 1.$$

Hence we can fix $\alpha > 0$ and choose 0 < r < R such that

$$\inf_{n} \widetilde{W}_{x_{n},1}^{C} (r \leq ||X_{s}|| \leq R) \geq 1 - \alpha.$$

We then have

$$B_n(\delta) \le \widetilde{W}_{x_n,1}^C (r \le ||X_s|| \le R; \chi(\delta, s, 1) > \epsilon) + \alpha$$
,

So, by the Markov property,

$$B_{n}(\delta) \leq \widetilde{W}_{x_{n},1}^{C} \left(r \leq \|X_{s}\| \leq R; \widetilde{W}_{X_{s},1-s}^{C} (\chi(\delta,0,1-s) > \epsilon) \right) + \alpha$$

$$\leq \sup \left\{ \widetilde{W}_{y,s}^{C} \left(\chi(\delta,0,s) > \epsilon \right) \right) : y \in C \text{ and } r \leq \|y\| \leq R \right\} + \alpha.$$

Now, if $D(\delta)$ did not tend to 0 as δ goes to 0, then we could find a sequence (δ_n) converging to 0 and a sequence (y_n) of points of C converging to a point $y \in \overline{C} \setminus \{0\}$ such that

$$\liminf_{n \to \infty} \widetilde{W}_{y_n,s}^C \left(\chi(\delta_n, 0, s) > \epsilon \right) > 0 ,$$

which would contradict the weak convergence of the sequence of probability measures $(\widetilde{W}_{y_n,s}^C)$ (Lemma 5.1 or Proposition 2.3 if $y \in C$). This proves that $\lim_{\delta \to 0} \limsup_n B_n(\delta) \leq \alpha$, and letting $\alpha \to 0$ then gives

$$\lim_{\delta \to 0} \limsup_{n \to \infty} B_n(\delta) = 0.$$

We now turn to $A_n(\delta)$. Let $\rho = \tau_{B(0,\epsilon)}$ be the exit time from the ball $B(0,\epsilon)$ with center at 0 and radius ϵ . Since the modulus of continuity of a path w is less than 2ϵ as long as it has not left the ball $B(0,\epsilon)$, we have

$$A_n(\delta) \leq \widetilde{W}_{x_n,1}^C \left(\rho < s; \chi(\delta, \rho, s) > \epsilon \right)$$

$$\leq \widetilde{W}_{x_n,1}^C \left(\rho < s; \widetilde{W}_{X(\rho),1-\rho}^C (\chi(\delta, 0, 1) > \epsilon) \right) .$$

Hence

 $\limsup_{n\to\infty} A_n(\delta)$

$$\leq \ \sup\{\widetilde{W}^C_{y,t}(\chi(\delta,0,1)>\epsilon): y\in C, \|y\|=\epsilon \text{ and } t\in[s,1]\}\;.$$

In the same way as above, we then get

$$\limsup_{n\to\infty} A_n(\delta) = 0 ,$$

which is sufficient to prove Proposition 5.9.

Together with Proposition 5.8, Proposition 5.9 proves that $W_{x,1}^C$ converges weakly on every C_T , T > 0, as $x \in C$ tends to 0. This is equivalent to weak convergence on C_{∞} ; thus Theorem 1.1 is proven.

The limit law will be denoted by $\widetilde{W}_{0,1}^C$ and called the law of C-Brownian meander. Thanks to Theorem 1.1 we shall say that the C-Brownian meander is Brownian motion conditioned to stay in C for a unit of time.

5.2. Properties of the C-Brownian meander. The C-Brownian meander starts from the vertex of the cone C and stays in it for a unit of time. From the Markov property and the expression of the densities of the entrance laws, the law of the exit time of C after time 1 is easily derived:

Proposition 5.10. For any t > 1, we have

$$\widetilde{W}_{0,1}^{C}(\tau_{C} > t) = t^{-\frac{\alpha_{1}}{2} + \frac{d-2}{4}}$$
.

Proof. By the Markov property, we have

$$\widetilde{W}_{0,1}^{C}(\tau_{C} > t) = \widetilde{W}_{0,1}^{C}(W_{X_{1}}(\tau_{C} > t - 1))$$

$$= \int_{C} e(1, y)W_{y}(\tau_{C} > t - 1) dy.$$

With the change of variables $y = \sqrt{tu}$, the integral becomes

$$\int_{C} e(1, \sqrt{tu}) W_{\sqrt{tu}}(\tau_{C} > t - 1) t^{\frac{d}{2}} du .$$

The scaling property of Brownian motion gives

$$W_{\sqrt{t}u}(\tau_C > t - 1) = W_u(\tau_C > 1 - 1/t)$$
,

and from relation (23) p. 27 it is easily seen that

$$e(1, \sqrt{tu})W_u(\tau_C > 1 - 1/t)t^{\frac{d}{2}} = t^{-\frac{\alpha_1}{2} + \frac{d-2}{4}}e(1/t, u)$$
.

The expected result follows from the fact that e(1/t, u) is a probability density.

5.3. Application to random walks. Let $(\xi_n)_{n\geq 1}$ be a sequence of i.i.d. random vectors of \mathbb{R}^d with $\mathbb{E}(\xi_n)=0$ and a covariance matrix equal to the identity. Form the random walk $S_n=\xi_1+\xi_2+\cdots+\xi_n$ and let $\{S_n(t),t\geq 0\}$ be the process with continuous paths for which $S_n(k/n)=S_k/\sqrt{n}$ and which is linearly interpolated elsewhere.

Consider a nice cone C for which $\mathbb{P}(\xi_1 \in C) > 0$, and let

$$T_C = \inf\{n > 0 : S_n \notin C\}$$

be the first exit time of the random walk from the cone C.

For d=1 and $C=(0,+\infty)$, Bolthausen proved in [3] that conditionally on $\{T_C>n\}$ the process \mathcal{S}_n converges in law to a Brownian meander; in other words, we have

$$\mathbb{P}(\mathcal{S}_n \in * | T_C > n) \Rightarrow \widetilde{W}_{0,1}^+.$$

This means that a random walk conditioned to stay positive converges in law to a Brownian motion conditioned to stay positive. By analogy, we are led to conjecture that a multidimensional random walk conditioned to stay in a cone should converge, in good cases, to a Brownian motion conditioned to stay in the same cone, that is,

(25)
$$\mathbb{P}(\mathcal{S}_n \in * | T_C > n) \Rightarrow \widetilde{W}_{0,1}^C.$$

Shimura proved it to be true in the two-dimensional case, under extra symmetry assumptions on the random walk and the cone (see [9]). It is clear that such an invariance principle can not hold without extra assumptions. For example, consider the simple random on \mathbb{Z}^3 and the cone $C = \{(x, y, z) : x \geq 0 \text{ and } x/2 \leq y \leq 2x\}$. Given $\{T_C > n\}$, $\{S_k, k = 1 \dots n\}$ is a simple random walk on $\{(0,0,z) : z \in \mathbb{Z}\}$. Hence it converges in law to a one-dimensional Brownian motion that lives on the z-axis. Note that in this case $\mathbb{P}(T_C > n) = 1/3^n$ decreases exponentially to 0. The next result states that a necessary condition for the above-mentioned limit theorem to hold is that $\mathbb{P}(T_C > n)$ decreases polynomialy to 0

Proposition 5.11. Suppose $d \geq 2$ and let C be a nice cone in \mathbb{R}^d . If

$$\mathbb{P}(\mathcal{S}_n \in * | T_C > n) \Rightarrow \widetilde{W}_{0,1}^C,$$

then

$$\mathbb{P}(T_C > n) = n^{-\frac{\alpha_1}{2} + \frac{d-2}{4}} L(n) ,$$

where L is a slowly varying function.

Proof. Set $\mu_n(*) = \mathbb{P}(S_n \in * | T_C > n)$. Let t > 1 and observe that

$$\frac{\mathbb{P}(T_C > [tn])}{\mathbb{P}(T_C > n)} = \mu_n(\tau_C > t_n) ,$$

where $t_n = [tn]/n$. Since $\mu_n \Rightarrow \widetilde{W}_{0,1}^C$ and $t_n \to t$, an argument similar to the one used in the proof of Lemma 4.2 shows that

$$\lim_{n\to\infty} \mu_n(\tau_C > t_n) = \widetilde{W}_{0,1}^C(\tau_C > t) .$$

By Proposition 5.10 we have therefore

$$\lim_{n \to \infty} \frac{\mathbb{P}(T_C > [tn])}{\mathbb{P}(T_C > n)} = t^{-\frac{\alpha_1}{2} + \frac{d-2}{4}}.$$

This proves that $L(n) = n^{\frac{\alpha_1}{2} - \frac{d-2}{4}} \mathbb{P}(T_C > n)$ is slowly varying. \square

For a two-dimensional random walk with bounded increments, the conditional limit theorem (25) holds under the assumption that $\mathbb{P}(T_C > n)$ is "slowly" decreasing to 0. We will present this result in a forthcoming paper.

6. Acknowledgements

This paper is related to the *Thèse de doctorat* of the author. We would like to thank our advisors, Emmanuel Lesigne and Marc Peigné for their constant help in preparing the thesis and the present paper.

APPENDIX A. THE CONTINUOUS MAPPING THEOREM

For the reader convenience we state here the continuous mapping theorem. The proof can be found for example in Billingsley [2].

Theorem A.1. Let X and Y be two metric spaces equipped with their Borel σ -algebras. Let (μ_n) , μ be probability measures on X such that $\mu_n \Rightarrow \mu$, and let (ϕ_n) , $\phi: X \to Y$ be measurable mappings.

If there exists a measurable subset X' of X such that $\mu(X') = 1$ and such that $\phi_n(x_n) \to \phi(x)$ for all sequence (x_n) converging to a point $x \in X'$, then

$$\mu_n \circ \phi_n^{-1} \Rightarrow \mu \circ \phi^{-1}$$
.

References

- [1] Bañuelos, R. and Smits, R.G. (1997). Brownian motion in cones. *Probab. Theory Related Fields* 108 (3): 299–319.
- [2] Billingsley, P. (1999). Convergence of Probability Measures. Second edition. Wiley, New York.
- [3] Bolthausen, E. (1976). On a functional central limit theorem for random walks conditioned to stay positive. Ann. Probab. 4 (3): 480–485.
- [4] Burdzy, K. (1986). Brownian excursions from hyperplanes and smooth surfaces. Trans. Amer. Math. Soc. 295 (1): 35–57.
- [5] Durrett, R. T., Iglehart, D. L. and Miller, D. R. (1977). Weak convergence to Brownian meander and Brownian excursion. *Ann. Probab.* 5 (1): 117–129.
- [6] Port, S. C. and Stone, C. J. (1978). Brownian Motion and Classical Potential Theory. Academic Press, New York.
- [7] Shimura, M. (1984). A limit theorem for two-dimensional conditioned random walk. Nagoya Math. J. 95: 105–116.
- [8] Shimura, M. (1985). Excursions in a cone for two-dimensional Brownian motion. J. Math. Kyoto Univ. 25 (3): 433-443.
- [9] Shimura, M. (1991). A limit theorem for two-dimensional random walk conditioned to stay in a cone. *Yokohama Math. J.* **39** (1): 21–36.
- [10] Williams, D. (1970). Decomposing the Brownian path. Bull. Amer. Math. Soc. 76: 871–873.

Laboratoire de Mathématiques Jean Leray, UMR CNRS 6629, Université de Nantes, BP 92208, 44322 Nantes Cedex 3, France.

E-mail address: rodolphe.garbit@univ-nantes.fr