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EXPLICIT EXPRESSION FOR THE STATIONARY DISTRIBUTION OF REFLECTED BROWNIAN MOTION IN A WEDGE

S. FRANCESCHI AND K. RASCHEL

ABSTRACT. For Brownian motion in a (two-dimensional) wedge with negative drift and oblique reflection on the axes, we derive an explicit formula for the Laplace transform of its stationary distribution (when it exists), in terms of Cauchy integrals and generalized Chebyshev polynomials. To that purpose we solve a Carleman-type boundary value problem on a hyperbola, satisfied by the Laplace transforms of the boundary stationary distribution.

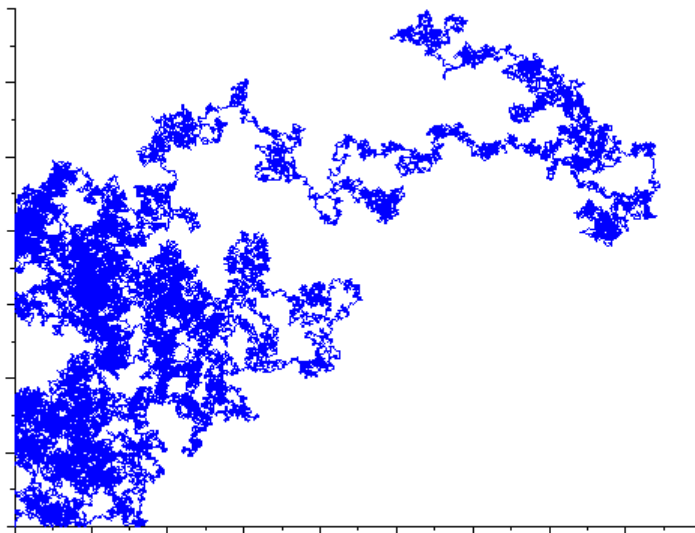


FIGURE 1. An example of path of reflected Brownian motion in the quadrant

1. INTRODUCTION AND MAIN RESULTS

Since its introduction in the eighties by Harrison, Reiman, Varadhan and Williams [31, 32, 54, 55, 56], reflected Brownian motion in the quarter plane has received a lot of attention from the probabilistic community. However, and surprisingly, finding a general explicit expression of the stationary distribution has been left as an open problem. The present paper solves this problem in a complete and unified way.

Reflected Brownian motion in two-dimensional cones. The semimartingale reflected Brownian motion with drift in the quarter plane $\mathbb{R}_+^2 := [0, \infty)^2$ (or equivalently in arbitrary convex wedges, by performing a simple linear transformation, cf. Appendix A) can be written as

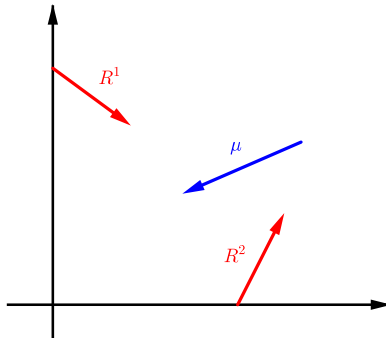
$$Z(t) = Z_0 + B(t) + \mu \cdot t + R \cdot L(t), \quad \forall t \geq 0, \quad (1)$$

where

- Z_0 is any initial point in the quadrant,

Key words and phrases. Reflected Brownian motion in a wedge; Stationary distribution; Laplace transform; Carleman-type boundary value problem; Boundary value problem with shift; Conformal mapping.

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FIGURE 2. Drift μ and reflection vectors R^1 and R^2

- B is a Brownian motion with covariance $\Sigma = \begin{pmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{12} & \sigma_{22} \end{pmatrix}$ starting from the origin,
- $\mu = \begin{pmatrix} \mu_1 \\ \mu_2 \end{pmatrix}$ denotes the interior drift,
- $R = (R^1, R^2) = \begin{pmatrix} r_{11} & r_{12} \\ r_{21} & r_{22} \end{pmatrix}$ is the reflection matrix,
- $L = \begin{pmatrix} L^1 \\ L^2 \end{pmatrix}$ is the local time.

For $i \in \{1, 2\}$, the local time $L^i(t)$ is a continuous non-decreasing process starting from 0 (i.e., $L^i(0) = 0$), increasing only at time t such that $Z^i(t) = 0$, viz., $\int_0^t \mathbb{1}_{\{Z^i(s) \neq 0\}} dL^i(s) = 0$, for all $t \geq 0$. The columns R^1 and R^2 represent the directions in which the Brownian motion is pushed when the axes are reached, see Figure 2.

The reflected Brownian motion $(Z(t))_{t \geq 0}$ associated with (Σ, μ, R) is well defined [54, 56], and is a fundamental stochastic process. This process has been extensively explored, and its multidimensional version (a semimartingale reflected Brownian motion in the positive orthant \mathbb{R}_+^d , as well as in convex polyhedrons) as well. It has been studied in depth, with focuses on its definition and semimartingale properties [54, 55, 16, 57], its recurrence or transience [55, 13, 35, 9, 7, 6, 15], the possible particular (e.g., product) form of its stationary distribution [34, 19, 46], its Lyapunov functions [23], its links with other stochastic processes [41, 22, 42], its use to approximate large queuing networks [26, 2, 33, 38, 37], the asymptotics of its stationary distribution [30, 17, 28, 48], numerical methods to compute the stationary distribution [13, 14], links with complex analysis [26, 2, 8, 29], comparison and coupling techniques [50, 49], etc.

The main contribution of the present paper is to find an exact expression for the stationary distribution (via its Laplace transforms, to be introduced in (3) and (4)), thanks to the theory of boundary value problems (BVPs), see our Theorem 1. This is one of the first attempts to apply boundary value techniques to diffusions in the quadrant, after [27] (under the symmetry conditions $\mu_1 = \mu_2$, $\sigma_{11} = \sigma_{22}$, and symmetric reflection vectors in (1)), [26] (which concerns very specific cases of the covariance matrix, essentially the identity), [2] (on diffusions with a special behavior on the boundary), [29] (orthogonal reflections, solved by Tutte's invariant approach [53, 3]); see also the introduction of [17], where the authors allude to the possibility of such an approach.

Laplace transforms and functional equation. The reflected Brownian motion defined in (1) exists if and only if

$$\{r_{11} > 0, \quad r_{22} > 0, \quad \det R > 0\} \quad \text{or} \quad \{r_{11} > 0, \quad r_{22} > 0, \quad r_{12} > 0, \quad r_{21} > 0\}.$$

This condition is equivalent for R to be completely- \mathcal{S} , which actually is a necessary and sufficient condition to the existence of reflected Brownian motion in arbitrary dimension, see [52, 47].

As for the stationary distribution, it exists if and only if

$$r_{11} > 0, \quad r_{22} > 0, \quad r_{11}r_{22} - r_{12}r_{21} > 0, \quad r_{22}\mu_1 - r_{12}\mu_2 < 0, \quad r_{11}\mu_2 - r_{21}\mu_1 < 0, \quad (2)$$

see [30], and in that case it is absolutely continuous w.r.t. the Lebesgue measure [13, 14, 33], with density denoted by $\pi(x) = \pi(x_1, x_2)$. Assumption (2) implies in particular that R is invertible and $R^{-1}\mu < 0$, which turns out to be a necessary condition for the existence of the stationary distribution in any dimension, see [33]. From now, we will assume that (2) is satisfied. Let the Laplace transform of π be defined by

$$\varphi(\theta) = \mathbb{E}_\pi[\exp(\theta \cdot Z)] = \iint_{\mathbb{R}_+^2} \exp(\theta \cdot x)\pi(x)dx. \quad (3)$$

Furthermore we define two finite boundary measures ν_1 and ν_2 such that, for $A \subset \mathbb{R}_+$,

$$\nu_1(A) = \mathbb{E}_\pi \left[\int_0^1 \mathbb{1}_{\{Z(t) \in \{0\} \times A\}} dL^1(t) \right], \quad \nu_2(A) = \mathbb{E}_\pi \left[\int_0^1 \mathbb{1}_{\{Z(t) \in A \times \{0\}\}} dL^2(t) \right].$$

The ν_i have their supports on the axes and may be viewed as boundary invariant measures. They are continuous w.r.t. the Lebesgue measure, see [34]. We define their Laplace transform by

$$\varphi_1(\theta_2) = \int_{\mathbb{R}_+} \exp(\theta_2 x_2)\nu_1(x_2)dx_2, \quad \varphi_2(\theta_1) = \int_{\mathbb{R}_+} \exp(\theta_1 x_1)\nu_2(x_1)dx_1. \quad (4)$$

The following functional equation relates the Laplace transforms:

$$-\gamma(\theta)\varphi(\theta) = \gamma_1(\theta)\varphi_1(\theta_2) + \gamma_2(\theta)\varphi_2(\theta_1), \quad (5)$$

where we have noted

$$\begin{cases} \gamma(\theta) = \frac{1}{2}(\theta \cdot \Sigma\theta) + \theta \cdot \mu = \frac{1}{2}(\sigma_{11}\theta_1^2 + 2\sigma_{12}\theta_1\theta_2 + \sigma_{22}\theta_2^2) + \mu_1\theta_1 + \mu_2\theta_2, \\ \gamma_1(\theta) = R^1 \cdot \theta = r_{11}\theta_1 + r_{21}\theta_2, \\ \gamma_2(\theta) = R^2 \cdot \theta = r_{12}\theta_1 + r_{22}\theta_2. \end{cases} \quad (6)$$

The Laplace transforms (3) and (4) (resp. Equation (5)) exist (resp. holds) at least for values of $\theta = (\theta_1, \theta_2)$ with $\Re\theta_1 \leq 0$ and $\Re\theta_2 \leq 0$. To prove the functional equation (5), the main idea is to use an identity called basic adjoint relationship (BAR); see [29, Section 2.1] and [17, 26] for details.

Kernel and associated quantities. In this paragraph we introduce necessary notation to state our main results. By definition, the kernel of (5) is the second degree polynomial γ . With this terminology, (5) is sometimes referred to as a kernel equation. The equality $\gamma(\theta_1, \theta_2) = 0$ with $\theta_1, \theta_2 \in \mathbb{C}$ defines algebraic functions $\Theta_1^\pm(\theta_2)$ and $\Theta_2^\pm(\theta_1)$ by

$$\gamma(\Theta_1^\pm(\theta_2), \theta_2) = \gamma(\theta_1, \Theta_2^\pm(\theta_1)) = 0.$$

Solving these equations readily yields

$$\begin{cases} \Theta_1^\pm(\theta_2) = \frac{-(\sigma_{12}\theta_2 + \mu_1) \pm \sqrt{\theta_2^2(\sigma_{12}^2 - \sigma_{11}\sigma_{22}) + 2\theta_2(\mu_1\sigma_{12} - \mu_2\sigma_{11}) + \mu_1^2}}{\sigma_{11}}, \\ \Theta_2^\pm(\theta_1) = \frac{-(\sigma_{12}\theta_1 + \mu_2) \pm \sqrt{\theta_1^2(\sigma_{12}^2 - \sigma_{11}\sigma_{22}) + 2\theta_1(\mu_2\sigma_{12} - \mu_1\sigma_{22}) + \mu_2^2}}{\sigma_{22}}. \end{cases} \quad (7)$$

The polynomials under the square roots in (7) have two zeros (called branch points), real and of opposite signs. They are denoted by θ_2^\pm and θ_1^\pm , respectively:

$$\begin{cases} \theta_2^\pm = \frac{(\mu_1\sigma_{12} - \mu_2\sigma_{11}) \pm \sqrt{(\mu_1\sigma_{12} - \mu_2\sigma_{11})^2 + \mu_1^2 \det \Sigma}}{\det \Sigma}, \\ \theta_1^\pm = \frac{(\mu_2\sigma_{12} - \mu_1\sigma_{22}) \pm \sqrt{(\mu_2\sigma_{12} - \mu_1\sigma_{22})^2 + \mu_2^2 \det \Sigma}}{\det \Sigma}. \end{cases} \quad (8)$$

Our next important definition is the curve

$$\mathcal{R} = \{\theta_2 \in \mathbb{C} : \gamma(\theta_1, \theta_2) = 0 \text{ and } \theta_1 \in (-\infty, \theta_1^-)\}.$$

As it will be seen in Lemma 5 (see also Figure 5), \mathcal{R} is a hyperbola. We denote by

$$\mathcal{R}^- = \{\theta_2 \in \mathcal{R} : \Im \theta_2 \leq 0\} \quad (9)$$

the negative imaginary part of \mathcal{R} oriented from the vertex to infinity.

We further define the function

$$w(\theta_2) = T_{\frac{\pi}{\beta}} \left(-\frac{2\theta_2 - (\theta_2^+ + \theta_2^-)}{\theta_2^+ - \theta_2^-} \right), \quad (10)$$

where $\beta = \arccos -\frac{\sigma_{12}}{\sqrt{\sigma_{11}\sigma_{22}}}$ and for $a \geq 0$, T_a is the so-called generalized Chebyshev polynomial

$$T_a(x) = \cos(a \arccos x) = \frac{1}{2} \left\{ (x + \sqrt{x^2 - 1})^a + (x - \sqrt{x^2 - 1})^a \right\}. \quad (11)$$

The function w plays a special role regarding the curve \mathcal{R} , as for $\theta_2 \in \mathcal{R}$ it satisfies $w(\theta_2) = w(\overline{\theta_2})$, see Lemma 9. (Here and throughout, $\overline{\theta_2}$ denotes the complex conjugate number of θ_2 .)

Finally, let G be the function whose expression is

$$G(\theta_2) = \frac{\gamma_1}{\gamma_2} (\Theta_1^-(\theta_2), \theta_2) \frac{\gamma_2}{\gamma_1} (\Theta_1^-(\overline{\theta_2}), \overline{\theta_2}). \quad (12)$$

Main results. In addition to (2), which guarantees existence and uniqueness of the stationary distribution, we shall assume that both coordinates of the drift are negative:

$$\mu_1 < 0, \quad \mu_2 < 0. \quad (13)$$

This hypothesis (also done in [27, 26, 28]) is only technical, and allows us to reduce the number of cases to handle. In Section 3.6 we comment on the case of a drift μ with one non-negative coordinate (having two non-negative coordinates is obviously incompatible with (2)). Our main result can be stated as follows.

Theorem 1. *Under the assumptions (2) and (13), the Laplace transform φ_1 in (4) is equal to*

$$\varphi_1(\theta_2) = \nu_1(\mathbb{R}_+) \left(\frac{w(0) - w(p)}{w(\theta_2) - w(p)} \right)^{-\chi} \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}^-} \log G(\theta) \left[\frac{w'(\theta)}{w(\theta) - w(\theta_2)} - \frac{w'(\theta)}{w(\theta) - w(0)} \right] d\theta \right\}, \quad (14)$$

where

- w , G and \mathcal{R}^- are defined in (10), (12) and (9), respectively,
- $\nu_1(\mathbb{R}_+) = \frac{r_{12}\mu_2 - r_{22}\mu_1}{\det R}$,
- the index χ is given by $\chi = \begin{cases} 0 & \text{if } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \leq 0, \\ -1 & \text{if } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) > 0, \end{cases}$
- $p = \frac{2r_{11}(\mu_1 r_{21} - \mu_2 r_{11})}{r_{11}^2 \sigma_{22} - 2r_{11} r_{21} \sigma_{12} + r_{21}^2 \sigma_{11}}$,
- to define the function $\log G(\theta)$ on \mathcal{R}^- , we use the determination of the logarithm taking a value in $i \cdot (-\pi, \pi]$ at the vertex of \mathcal{R}^- and varying continuously over the curve \mathcal{R}^- .

The function $\varphi_2(\theta_1)$ equals $\varphi_1(\theta_1)$ in (14) after the change of parameters

$$\sigma_{11} \leftrightarrow \sigma_{22}, \quad \mu_1 \leftrightarrow \mu_2, \quad r_{11} \leftrightarrow r_{22}, \quad r_{12} \leftrightarrow r_{21}.$$

The functional equation (5) finally gives an explicit formula for the bivariate Laplace transform φ .

Let us now give some comments around Theorem 1.

- Theorem 1 completely generalizes the results of [27] (with symmetry conditions), [26] (with the identity covariance matrix Σ) and [29] (orthogonal reflections on the axes). It offers the first explicit expression of the Laplace transforms, covering all the range of parameters (Σ, μ, R) satisfying to (2) and (13), thereby solving an old open problem.
- We obtain three corollaries of Theorem 1, each of those corresponds to an already known result: we first compute the asymptotics of the boundary densities (Section 4.1, initially obtained in [18]); second we derive the product form expression of the density in the famous skew-symmetric case (Section 5.2, result originally proved in [34]); finally we compute the expression of the Laplace transforms in the case of orthogonal reflections (Section 5.4, see [29] for the first derivation).
- It is worth remarking that the expression (14) is intrinsically non-continuous in terms of the parameters: the index χ can indeed take two different values (namely, 0 and 1). For this reason, (14) actually contains two different formulas. See Remarks 7 and 18 for further related comments.
- The paper [28] obtains the exact asymptotics of the stationary distribution along any direction in the quarter plane, see [28, Theorems 22–28]. Constants in these asymptotics involve the functions φ_1 and φ_2 in (4), and can thus be made explicit with Theorem 1.
- It is also interesting to dissect (14) and to notice that certain quantities in that formula depend only on the behavior of the process in the interior of the quadrant (w and \mathcal{R}^-), while the remaining ones mix properties of the interior and boundary of the quarter plane ($\nu_1(\mathbb{R}_+)$, p , χ and G).
- The statement of Theorem 1 (namely, an expression of the Laplace transform as a Cauchy integral), as well as the techniques we shall employ to prove it (viz., reduction to BVPs with shift), are reminiscent of the results and methods used for discrete random walks in the quarter plane, see [25] for a modern reference, and [44, 24, 39] for historical breakthroughs.
- Altogether, Theorem 1 illustrates that the analytic approach consisting in solving quarter plane problems via BVPs is better suited for diffusions than for discrete random walks. We can actually treat any wedge, covariance matrix, drift vector and reflection vectors (see Corollary 2 below), whereas in the discrete case, hypotheses should be done on the boundedness of the jumps (only small steps are considered in [25, 39, 5, 3]) and on the cone (typically, half and quarter planes only).

Theorem 1 further leads to an explicit expression for the Laplace transform of the stationary distribution of reflected Brownian motion in an arbitrary convex wedge, as stated in the following corollary, whose proof is postponed to Appendix A.

Corollary 2. *Let \tilde{Z} be a reflected Brownian motion in a wedge of angle $\beta \in (0, \pi)$, of covariance matrix $\tilde{\Sigma}$, drift $\tilde{\mu}$ and reflection matrix $\tilde{R} = (\tilde{R}^1, \tilde{R}^2)$, corresponding to the angles of reflection δ and ε on Figure 10. Assume it is recurrent and note $\tilde{\pi}$ the stationary distribution and $\tilde{\varphi}$ its Laplace transform. Let*

$$T_1 = \begin{pmatrix} \frac{1}{\sin \beta} & \cot \beta \\ 0 & 1 \end{pmatrix}, \quad T_1^{-1} = \begin{pmatrix} \sin \beta & -\cos \beta \\ 0 & 1 \end{pmatrix}.$$

Then

$$\tilde{\varphi}(\tilde{\theta}) = \varphi(T_1^\top \tilde{\theta}),$$

where φ is the Laplace transform (14) of Theorem 1 associated to (Σ, μ, R) , with

$$\Sigma = T_1 \tilde{\Sigma} T_1^\top, \quad \mu = T_1^{-1} \tilde{\mu} \quad \text{and} \quad R = T_1^{-1} \tilde{R}.$$

Structure of the paper.

- Section 2: analytic preliminaries, continuation of the Laplace transforms and definition of an important hyperbola

- Section 3: statement and proof that the Laplace transforms satisfy BVP of Carleman-type on branches of hyperbolas, transformation of the Carleman BVP with shift into a (more classical) Riemann BVP, study of the conformal mapping allowing this transformation, resolution of the BVP
- Section 4: asymptotics of the stationary distribution, links with Dai and Miyazawa's [18] asymptotic results
- Section 5: simplifications of the integral expression of Theorem 1 for models satisfying the skew-symmetric condition, for orthogonal reflections (leading to a new proof of the results of [29]), links with Dieker and Moriarty's results [19]
- Appendix A: equivalence between Brownian motion in the quarter plane and Brownian motion in convex wedges

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2. METHODOLOGY AND ANALYTIC PRELIMINARIES

2.1. Methodology and positioning of our work regarding [27, 26, 2, 17]. Schematically, our argument for the proof of Theorem 1 is composed of the following steps:

- (1) presentation of the functional equation, analytic preliminaries, meromorphic continuation of the Laplace transforms (Section 2);
- (2) statement of a Carleman BVP with shift satisfied by the Laplace transforms (Section 3.1);
- (3) introduction of a conformal mapping, allowing to transform the latter BVP into a more classical Riemann BVP (Section 3.2), see Figure 6;
- (4) statement of the Riemann BVP (Section 3.3);
- (5) definition and study of the index (denoted by χ in Theorem 1), which turns out to have a crucial role in solving the Riemann BVP (Section 3.4);
- (6) resolution of the BVP (Section 3.5).

Except for the study of the index (item (5), which is specific to our problem at hand), the above structuration of the proof dates back to (in chronological order) [24] (for the discrete setting, which later led to the book [25]), [27], [26] and [2].

To begin the discussion, let us remind that the process in [2] is absorbed on the boundary and then relaunched in the quadrant after an exponential time; it is therefore totally different from ours. Let us also recall that the analytic study of reflected Brownian motion in the quadrant was initiated in [27, 26], but both [27] and [26] did quite restrictive assumptions on the process (symmetry conditions $\mu_1 = \mu_2$, $\sigma_{11} = \sigma_{22}$, and symmetric reflection vectors in [27], identity covariance matrix in [26]).

In this paper we follow the same steps (1)–(6) and use a synthetic approach of [24, 27, 26, 2], to eventually go further than the existing literature and prove Theorem 1. Some technical details of the present work are borrowed from [27, 26, 2] and more particularly from [2]. Although being different, the stochastic process of [2] and ours share the property of satisfying a functional equation with the same left-hand side (as in (5)). Accordingly, technical details will be very similar as soon as they concern the kernel; this is the case of Section 2.

On the other hand, many points of our analysis profoundly differ from that of [2]: the resolution of the BVP is different as the right-hand side of (5) is not comparable to that of [2]. Moreover we make our main result as tractable as possible: see our Section 5, where we present many techniques to simplify the analytic expression derived in Theorem 1. We further also propose asymptotic considerations, in close relation to those of [17].

2.2. Real points of the kernel. The set of real points of the zero set of the kernel

$$\{(\theta_1, \theta_2) \in \mathbb{R}^2 : \gamma(\theta_1, \theta_2) = 0\}$$

defines an ellipse, see Figure 3. Introduced in [17], this ellipse offers the possibility of presenting many analytical quantities in a clear and compact way:

- the branch points θ_1^- and θ_1^+ (resp. θ_2^- and θ_2^+) in (8) are the leftmost and rightmost (resp. bottommost and topmost) points on the ellipse;
- the drift μ is orthogonal to the tangent at the origin;
- the set of points where $\gamma_1(\theta_1, \theta_2) = 0$ and $\gamma_2(\theta_1, \theta_2) = 0$ are straight lines, orthogonal to the reflection vectors R^1 and R^2 . Their intersection points with the ellipse are easily computed (notice that these intersection points appear in the statement of Theorem 1, in particular in the index χ).

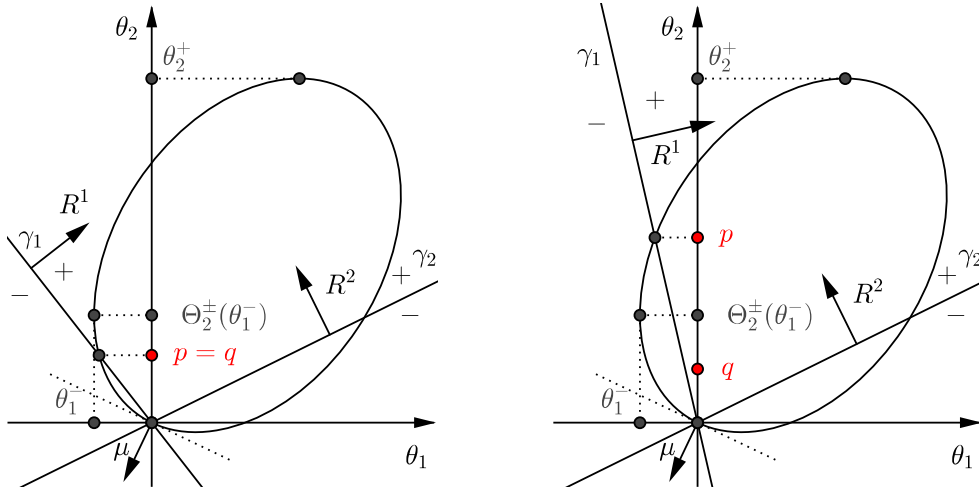


FIGURE 3. The ellipse and the two straight lines are the sets of real points (θ_1, θ_2) which cancel γ , γ_1 and γ_2 , respectively. Location of p (see (17)) and q (see (24)) according to the sign of $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-))$

2.3. Meromorphic continuation and poles of the Laplace transforms. In Section 3.1 we shall state a boundary condition for φ_1 , on a curve lying outside its natural domain of definition (namely, the half plane with negative real part). The statement hereafter (straightforward consequence of the functional equation (5), see Proposition 17 for an extended version) proposes a meromorphic continuation on a domain containing the latter curve.

Lemma 3. *The Laplace transform $\varphi_1(\theta_2)$ can be extended meromorphically to the open and simply connected set*

$$\{\theta_2 \in \mathbb{C} : \Re \theta_2 \leq 0 \text{ or } \Re \Theta_1^-(\theta_2) < 0\}, \quad (15)$$

by mean of the formula

$$\varphi_1(\theta_2) = \frac{\gamma_2}{\gamma_1}(\Theta_1^-(\theta_2), \theta_2) \varphi_2(\Theta_1^-(\theta_2)). \quad (16)$$

Proof. The domain (15) is simply connected by [28, §2.4]. The formula (16), see [28, Lemma 6], is a direct consequence of the functional equation (5) evaluated at $(\Theta_1^-(\theta_2), \theta_2)$, first on the (non-empty) open domain

$$\{\theta_2 \in \mathbb{C} : \Re \theta_2 < 0 \text{ and } \Re \Theta_1^-(\theta_2) < 0\}. \quad \square$$

Due to the continuation formula (16), the only possible pole of φ_1 in the domain (15) will come from a cancelation of the denominator γ_1 . More precisely, let p be the (unique, when it exists) non-zero point such that

$$\gamma_1(\Theta_1^-(p), p) = 0. \quad (17)$$

It follows from (17) that p satisfies a second degree polynomial equation with real coefficients. As one of the roots is 0 the other one must be p , which is then real and equals (when it exists)

$$p = \frac{2r_{11}(\mu_1 r_{21} - \mu_2 r_{11})}{r_{11}^2 \sigma_{22} - 2r_{11} r_{21} \sigma_{12} + r_{21}^2 \sigma_{11}}. \quad (18)$$

Formula (17) means that p is the ordinate of the intersection point between the ellipse $\gamma = 0$ and the line $\gamma_1 = 0$, see Figure 3. The intersection point always exists but sometimes its abscissa is associated to Θ_1^- (when p exists, i.e., $\gamma_1(\Theta_1^\pm(\theta_2^+), \theta_2^+) \geq 0$) and sometimes to Θ_1^+ (when p doesn't exist, i.e., $\gamma_1(\Theta_1^\pm(\theta_2^+), \theta_2^+) < 0$), the limit case being $p = \theta_2^+$, see Figure 4.

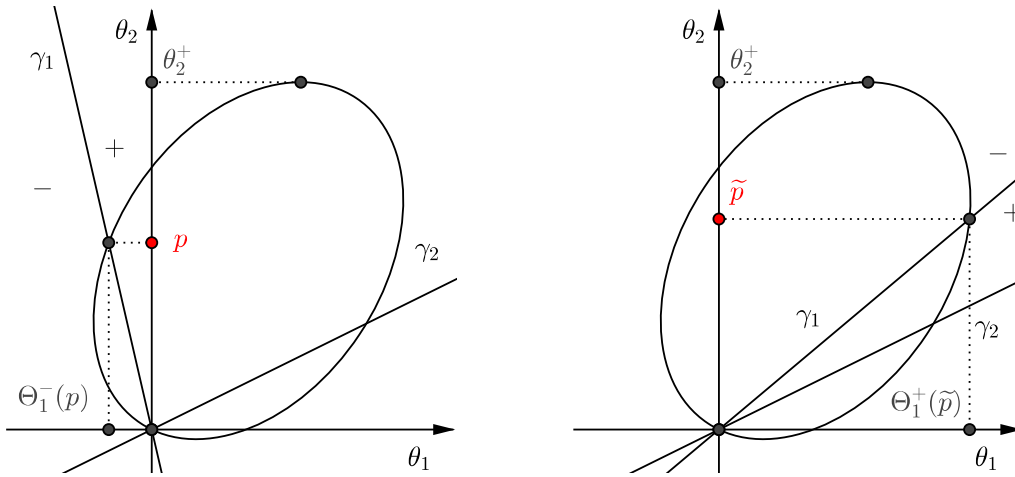


FIGURE 4. The intersection point between the ellipse $\gamma = 0$ and the straight line $\gamma_1 = 0$. On the left side p exists. On the right side p doesn't exist, although \tilde{p} , defined by that $\gamma_1(\Theta_1^+(\tilde{p}), \tilde{p}) = 0$, exists

Let us finally remark that the pole that φ_1 may have at p is necessarily simple, due to the expression (6) of γ_1 .

Remark 4. The continuation formula (16) uses the branch Θ_1^- in (7) and not Θ_1^+ . The reason is that Θ_1^- is the small branch of the algebraic function Θ_1^\pm , taking the value 0 at 0. Note, this is peculiar to the case of a drift with negative coordinates (our hypothesis (13)). More details can be found in the proof of Theorem 11 in [28].

2.4. An important hyperbola. For further use, we need to introduce the curve

$$\mathcal{R} = \{\theta_2 \in \mathbb{C} : \gamma(\theta_1, \theta_2) = 0 \text{ and } \theta_1 \in (-\infty, \theta_1^-)\} = \Theta_2^\pm((-\infty, \theta_1^-)). \quad (19)$$

It is symmetrical w.r.t. the horizontal axis, see Figure 5. Indeed, the discriminant of Θ_2^\pm (i.e., the polynomial under the square root in (7)) is positive on (θ_1^-, θ_1^+) and negative on $\mathbb{R} \setminus [\theta_1^-, \theta_1^+]$. Accordingly, the branches Θ_2^\pm take respectively real and complex conjugate values on the sets above. Furthermore, \mathcal{R} has a simple structure, as shown by the following elementary result:

Lemma 5 (Lemma 9 in [2]). *The curve \mathcal{R} in (19) is a branch of hyperbola, given by the equation*

$$\sigma_{22}(\sigma_{12}^2 - \sigma_{11}\sigma_{22})x^2 + \sigma_{12}^2\sigma_{22}y^2 - 2\sigma_{22}(\sigma_{11}\mu_2 - \sigma_{12}\mu_1)x = \mu_2(\sigma_{11}\mu_2 - 2\sigma_{12}\mu_1). \quad (20)$$

We denote the part of \mathcal{R} with negative imaginary part by \mathcal{R}^- , see (9) and Figure 5. We further denote by $\mathcal{G}_{\mathcal{R}}$ the open domain of \mathbb{C} containing 0 and bounded by \mathcal{R} , see again Figure 5. The closure of $\mathcal{G}_{\mathcal{R}}$ is equal to $\mathcal{G}_{\mathcal{R}} \cup \mathcal{R}$ and will be noted $\overline{\mathcal{G}_{\mathcal{R}}}$.

The domain in (15) strictly contains $\overline{\mathcal{G}_{\mathcal{R}}}$, see [28]. Thanks to Lemma 3, this implies that the Laplace transform $\varphi_1(\theta_2)$ can be extended meromorphically to a domain containing $\overline{\mathcal{G}_{\mathcal{R}}}$. To see the above inclusion we refer to [28], where the boundaries of the domain (15) (called $\Delta \cup \{s_0\}$ in [28]) is computed, see in particular [28, Figures 10 and 11]. The technique used in [28] is the parametrization of the zero set of the kernel (6).

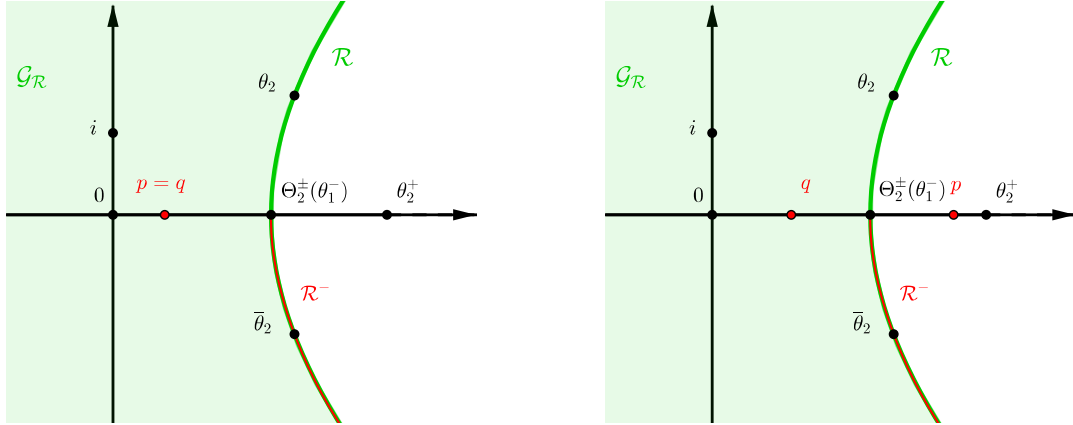


FIGURE 5. The curve \mathcal{R} in (19) is symmetric w.r.t. the horizontal axis, and $\mathcal{G}_{\mathcal{R}}$ is the domain in green. The curve \mathcal{R}^- is the half branch of \mathcal{R} with negative imaginary part. The points p and q are used to define the conformal mapping W , see (17) and (26). If the pole p is in $\mathcal{G}_{\mathcal{R}}$ then $q = p$ (figure on the left), otherwise $q = \frac{1}{2}\Theta_2^\pm(\theta_1^-)$ (figure on the right), see (24)

3. A PROOF OF THEOREM 1 VIA REDUCTION TO BVPS

3.1. **Carleman BVP.** For $\theta_2 \in \mathcal{R}$, define the function G as in (12):

$$G(\theta_2) = \frac{\gamma_1}{\gamma_2}(\Theta_1^-(\theta_2), \theta_2) \frac{\gamma_2}{\gamma_1}(\Theta_1^-(\bar{\theta}_2), \bar{\theta}_2).$$

Notice that $G(\theta_2)G(\bar{\theta}_2) = 1$ and that for $\theta_2 \in \mathcal{R}$ one has $\Theta_1^-(\bar{\theta}_2) = \Theta_1^-(\theta_2)$. Let us also recall that p is defined in (17).

Proposition 6 (Carleman BVP with shift). *The function φ_1 in (4)*

- (1) *is meromorphic on $\mathcal{G}_{\mathcal{R}}$,*
 - *without pole on $\mathcal{G}_{\mathcal{R}}$ if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) < 0$,*
 - *with a single pole on $\mathcal{G}_{\mathcal{R}}$ at p of order one if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) > 0$,*
 - *without pole on $\mathcal{G}_{\mathcal{R}}$ and with a single pole of order one on the boundary \mathcal{R} of $\mathcal{G}_{\mathcal{R}}$, at $p = \Theta_2^\pm(\theta_1^-)$, if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$,*
- (2) *is continuous on $\overline{\mathcal{G}_{\mathcal{R}}} \setminus \{p\}$ and bounded at infinity,*
- (3) *satisfies the boundary condition*

$$\varphi_1(\bar{\theta}_2) = G(\theta_2)\varphi_1(\theta_2), \quad \forall \theta_2 \in \mathcal{R}. \quad (21)$$

It is worth mentioning that the condition on the sign of $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-))$ has a clear geometric meaning: indeed, $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-))$ is negative (resp. positive) if and only if the straight line

corresponding to γ_1 crosses the ellipse below (resp. above) the ordinate $\Theta_2^\pm(\theta_1^-)$; see Figure 3, left (resp. right).

Remark 7. *Item (1) of Proposition 6 shows that according to the values of the parameters, various cases exist regarding the singularities of the Laplace transform in the domain $\mathcal{G}_{\mathcal{R}}$. This is the reason why there isn't a unique expression for the Laplace transform in our main Theorem 1, but two different expressions.*

Proof of Proposition 6. First of all, it follows from Lemma 3 that φ_1 is meromorphic in $\mathcal{G}_{\mathcal{R}}$ and may have a pole of order one at p . Indeed, due to the continuation formula (16), the only potential pole p of φ_1 in $\mathcal{G}_{\mathcal{R}}$ should be a zero of γ_1 . It is then on the real line and characterized by (17). Moreover, p defined by (17) is smaller than $\Theta_2^\pm(\theta_1^-)$ (i.e., $p \in \mathcal{G}_{\mathcal{R}}$) if and only if the geometric condition $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) > 0$ is satisfied, see Figure 3. This demonstrates the first item of Proposition 6.

The second item (in particular the fact that φ_1 is bounded at infinity) comes from Lemma 3 together with the fact that (4) implies that φ_1 (resp. φ_2) is bounded on the set $\{\theta_2 \in \mathbb{C} : \Re \theta_2 \leq 0\}$ (resp. $\{\theta_2 \in \mathbb{C} : \Re \theta_1 \leq 0\}$).

To prove the boundary condition (21) (that we announced in [29, Proposition 7]), we consider θ_1 such that $\Re \theta_1 < 0$, and evaluate the functional equation (5) at $(\theta_1, \Theta_2^\pm(\theta_1))$. This implies

$$\frac{\gamma_1}{\gamma_2}(\theta_1, \Theta_2^\pm(\theta_1))\varphi_1(\Theta_2^\pm(\theta_1)) + \varphi_2(\theta_1) = 0,$$

which in turn yields

$$\frac{\gamma_1}{\gamma_2}(\theta_1, \Theta_2^+(\theta_1))\varphi_1(\Theta_2^+(\theta_1)) = \frac{\gamma_1}{\gamma_2}(\theta_1, \Theta_2^-(\theta_1))\varphi_1(\Theta_2^-(\theta_1)). \quad (22)$$

Restricting (22) to values of $\theta_1 \in (-\infty, \theta_1^-)$, for which $\Theta_2^+(\theta_1)$ and $\Theta_2^-(\theta_1)$ are complex conjugate (see Section 2), we reach the conclusion that

$$\varphi_1(\overline{\theta_2}) = \frac{\gamma_1}{\gamma_2}(\theta_1, \theta_2) \frac{\gamma_2}{\gamma_1}(\theta_1, \overline{\theta_2})\varphi_1(\theta_2),$$

which, by definition (12) of G , exactly coincides with the boundary condition (21). Although we do not exclude a priori the denominators in (22) to vanish, note that this does not happen for $\theta_1 \in (-\infty, \theta_1)$ since then the imaginary part of $\Theta_2^\pm(\theta_1)$ is non-zero. \square

The BVP established in Proposition 6 belongs to the class of homogeneous Carleman (or Riemann-Carleman) BVPs with shift, see [43], the shift being here the complex conjugation.

In some cases, the function G in (21) can be factorized, leading to an interesting particular case, that we comment below. As we shall see, the well-known skew-symmetric condition

$$2\Sigma = R \cdot \text{diag}(R)^{-1} \cdot \text{diag}(\Sigma) + \text{diag}(\Sigma) \cdot \text{diag}(R)^{-1} \cdot R^\top \quad (23)$$

(equivalent for the stationary distribution $\pi(x_1, x_2)$ to have a product-form) gives a family of examples where such a factorization holds. In (23) we have noted $\text{diag}(A)$ the diagonal matrix with the same diagonal coefficients as those of A .

Remark 8. *If there exists a rational function F such that*

$$G(\theta_2) = \frac{F(\theta_2)}{F(\overline{\theta_2})},$$

one can transform the boundary condition (21) for φ_1 with $G \neq 1$ into a boundary condition for $\varphi_1 \cdot F$ with $G = 1$, namely,

$$(\varphi_1 \cdot F)(\overline{\theta_2}) = (\varphi_1 \cdot F)(\theta_2).$$

Then the associated BVP could be solved using Tutte's invariants [3]. This is what has been done in [29], for the particular case of orthogonal reflections (corresponding to $F(\theta_2) = \frac{1}{\theta_2}$).

We show in Section 5.2 that such a rational function F always exists in the skew-symmetric case (23), and from this we derive a rational expression of the Laplace transform.

The existence of a rational function F factorizing G as above is reminiscent of the notion of decoupling function or telescoper, introduced in [3, 21].

However, a rational factorization term F as in Remark 8 does not exist in general, and it is still an open problem to characterize the parameters (Σ, μ, R) for which F exists. As a consequence, we cannot systematically use Tutte's invariants technique: we are left with transforming the BVP of Proposition 6 into a more classical one, using a certain conformal mapping having a very convenient gluing property.

3.2. Conformal gluing. Our main result in this section is to prove that the function W defined by Equation (26) below satisfies the properties of Lemma 9, allowing to transform the Carleman BVP with shift on the curve \mathcal{R} of Proposition 6 into a classical BVP on the segment $[0, 1]$.

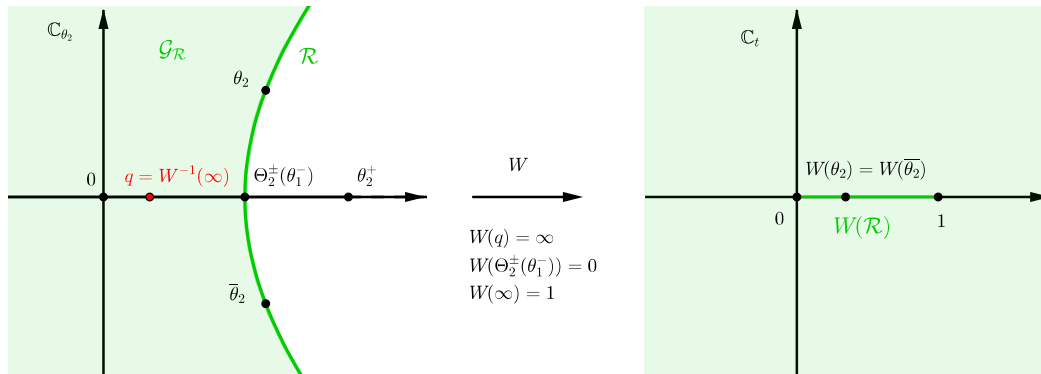


FIGURE 6. Domains, curves and points related to the Carleman BVP with shift on \mathcal{R} (left) and the standard BVP on $[0, 1]$ (right)

First we need to define q by

$$q = \begin{cases} p & \text{if } \gamma_1(\theta_1^-, \Theta_2^{\pm}(\theta_1^-)) > 0, \text{ i.e., if } \varphi_1 \text{ admits } p \in \mathcal{G}_{\mathcal{R}} \text{ as a pole,} \\ \frac{1}{2}\Theta_2^{\pm}(\theta_1^-) & \text{otherwise.} \end{cases} \quad (24)$$

Note, the choice $\frac{1}{2}\Theta_2^{\pm}(\theta_1^-)$ is arbitrary: any point in $\mathcal{G}_{\mathcal{R}}$ would have been suitable. See Figures 3, 5 and 6.

The function W is built on the function w below (note, w is introduced in [29, Theorem 1]; under the symmetry conditions $\mu_1 = \mu_2$, $\sigma_{11} = \sigma_{22}$, and symmetric reflection vectors in (1), Foschini [27] also obtained an expression for the conformal mapping w , see [27, Figure 3]; see finally [2, Equation (4.6)] for a related formulation of w):

$$w(\theta_2) = T_{\frac{\pi}{\beta}} \left(-\frac{2\theta_2 - (\theta_2^+ + \theta_2^-)}{\theta_2^+ - \theta_2^-} \right),$$

which itself uses the branch points (8), the generalized Chebyshev polynomial (11) and the angle

$$\beta = \arccos -\frac{\sigma_{12}}{\sqrt{\sigma_{11}\sigma_{22}}}. \quad (25)$$

Then we define

$$W(\theta_2) = \frac{w(\theta_2) - w(\Theta_2^{\pm}(\theta_1^-))}{w(\theta_2) - w(q)} = \frac{w(\theta_2) + 1}{w(\theta_2) - w(q)}. \quad (26)$$

The last equality is due to the identity

$$\Theta_2^{\pm}(\theta_1^-) = \frac{\theta_2^+ + \theta_2^-}{2} - \frac{\theta_2^+ - \theta_2^-}{2} \cos \beta.$$

We can see it by a direct computation or using a uniformization of the zero set of the kernel, see [29, Section 5] for more details. Then we have $w(\Theta_2^\pm(\theta_1^-)) = T_{\frac{\pi}{\beta}}(\cos \beta) = \cos \pi = -1$.

Lemma 9. *The function W in (26)*

- (i) *is analytic in $\mathcal{G}_{\mathcal{R}} \setminus \{q\}$, continuous in $\overline{\mathcal{G}_{\mathcal{R}}} \setminus \{q\}$ and bounded at infinity,*
- (ii) *is one-to-one from $\mathcal{G}_{\mathcal{R}} \setminus \{q\}$ onto $\mathbb{C} \setminus [0, 1]$,*
- (iii) *satisfies $W(\theta_2) = W(\overline{\theta_2})$ for all $\theta_2 \in \mathcal{R}$.*

Proof. It can be found in [29, Lemma 6] that w in (10)

- (i') *is analytic in $\mathcal{G}_{\mathcal{R}}$, continuous in $\overline{\mathcal{G}_{\mathcal{R}}}$ and unbounded at infinity (T_a admits an analytic continuation on $\mathbb{C} \setminus (\theta_2^+, \infty)$, and even on \mathbb{C} if a is a non-negative integer: in that case T_a is the classical Chebyshev polynomial of the first kind),*
- (ii') *is one-to-one from $\mathcal{G}_{\mathcal{R}}$ onto $\mathbb{C} \setminus (-\infty, -1]$,*
- (iii') *satisfies $w(\theta_2) = w(\overline{\theta_2})$ for all $\theta_2 \in \mathcal{R}$.*

Here we want to define another conformal gluing function, which glues together the upper part and the lower part of the hyperbola onto the segment $[0, 1]$, and which sends the point q in (24) at infinity, see Figure 6. For this reason we set W as in (26): by construction $W(\Theta_2^\pm(\theta_1^-)) = 0$, $W(\infty) = 1$ and $W(q) = \infty$. The proof of Lemma 9 follows from the above-mentioned properties (i')–(iii') of w together with the definition (26) of W . \square

Remark 10. *The algebraic nature of the mapping w in (10) (or equivalently W in (26)) is directly related to the rationality of $\frac{\beta}{\pi}$. Precisely, as shown in [29, Proposition 13], the following behaviors are possible:*

- *If $\frac{\beta}{\pi} \in \mathbb{Q}$, then w is algebraic;*
- *If in addition $\frac{\pi}{\beta} \in \mathbb{N}$ (and only in this case), then w is a polynomial;*
- *If $\frac{\beta}{\pi} \notin \mathbb{Q}$, then w is non-algebraic.*

3.3. Reduction to a standard BVP. Thanks to the gluing function W in (26), we are able to reformulate the Carleman BVP as a standard BVP for an open contour. See Figure 6 for a compact view of the two complex planes associated to the Carleman's and Riemann's BVPs. Define ψ_1 by

$$\psi_1(t) = \varphi_1 \circ W^{-1}(t), \quad \forall t \in \mathbb{C} \setminus [0, 1] \quad (27)$$

(note, ψ_1 is meromorphic on $\mathbb{C} \setminus [0, 1]$). Equivalently we have $\varphi_1(\theta_2) = \psi_1 \circ W(\theta_2)$ for $\theta_2 \in \mathcal{G}_{\mathcal{R}}$. Obviously W^{-1} is not well defined on $[0, 1]$; however, it does admit upper and lower limits for $t \in [0, 1]$:

$$(W^{-1})^+(t) = \lim_{\substack{u \rightarrow t \\ \Im u > 0}} W^{-1}(u), \quad (W^{-1})^-(t) = \lim_{\substack{u \rightarrow t \\ \Im u < 0}} W^{-1}(u),$$

and similarly for $\psi_1^+(t)$ and $\psi_1^-(t)$. Then for $\theta_2 \in \mathcal{R}$ and $t = W(\theta_2) = W(\overline{\theta_2})$, we have

$$\varphi_1(\theta_2) = \begin{cases} \psi_1^+(t) & \text{if } \Im \theta_2 > 0, \\ \psi_1^-(t) & \text{if } \Im \theta_2 < 0, \end{cases} \quad \varphi_1(\overline{\theta_2}) = \begin{cases} \psi_1^-(t) & \text{if } \Im \theta_2 > 0, \\ \psi_1^+(t) & \text{if } \Im \theta_2 < 0. \end{cases}$$

Define further

$$H(t) = G((W^{-1})^-(t)), \quad \forall t \in [0, 1]. \quad (28)$$

Then Proposition 6 becomes:

Proposition 11 (Riemann BVP). *The function ψ_1 in (27)*

- (1) *is analytic in $\mathbb{C} \setminus [0, 1]$, bounded at infinity if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \leq 0$ and admitting a simple pole at infinity otherwise,*
- (2) *is continuous on $[0, 1]$ from below (with limits ψ_1^-) and above (with limits ψ_1^+), bounded at 0 and 1 (except if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$: in this case it has a pole of order one at 0),*
- (3) *satisfies, with H defined in (28), the boundary condition*

$$\psi_1^+(t) = H(t)\psi_1^-(t), \quad \forall t \in [0, 1]. \quad (29)$$

Proof. Items (1) and (2) directly follow from the corresponding items in Proposition 6. With the above definitions, the boundary equation (21) becomes

$$\begin{cases} \psi_1^-(t) = G(\theta_2)\psi_1^+(t) & \text{if } \Im \theta_2 > 0, \\ \psi_1^+(t) = G(\theta_2)\psi_1^-(t) & \text{if } \Im \theta_2 < 0. \end{cases}$$

Since $\frac{1}{G(\theta_2)} = G(\overline{\theta_2}) = H(t)$ if $\Im \theta_2 > 0$, and $G(\theta_2) = H(t)$ if $\Im \theta_2 < 0$, the last item follows. \square

3.4. Index of the BVP. The resolution of BVPs as in Proposition 11 heavily depends on the *index* χ (see, e.g., [43, Section 5.2]), which is related to the variation of argument of H on $[0, 1]$:

$$\Delta = [\arg H]_0^1, \quad \delta = \arg H(0) \in (-\pi, \pi], \quad \chi = \left\lfloor \frac{\delta + \Delta}{2\pi} \right\rfloor. \quad (30)$$

Δ quantifies the variation of argument of H on $[0, 1]$ and $\arg H(1) = \delta + \Delta$. Since $(W^{-1})^-([0, 1]) = \mathcal{R}^-$, Δ in (30) can be equivalently written $[\arg G]_{\mathcal{R}^-}$ (from the vertex to infinity).

Remark 12. *It is important to notice that $\delta \in (-\pi, \pi]$ in (30) corresponds to an arbitrary choice. Any other choice would eventually lead to the same Theorem 1 (though written slightly differently).*

First, we compute δ in (30).

Lemma 13. *We have*

$$\delta = \begin{cases} 0 & \text{if } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \neq 0, \\ \pi & \text{if } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0. \end{cases} \quad (31)$$

The angle $\delta + \Delta \in (-2\pi, 2\pi)$ and we have

$$\tan \frac{\delta + \Delta}{2} = \frac{\det R \cdot \det \Sigma}{\sigma_{12}(r_{11}r_{22} + r_{12}r_{21}) - \sigma_{22}r_{11}r_{12} - \sigma_{11}r_{22}r_{21}}. \quad (32)$$

Note that the denominator of (32) can be negative, zero or positive, depending on the parameters.

Proof of Lemma 13. First of all we show the formula (31). If first $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \neq 0$ then $H(0) = G(\Theta_2^\pm(\theta_1^-)) = 1$, since $\Theta_2^\pm(\theta_1^-) \in \mathbb{R}$ simplifies the quotient (12) and then $\delta = \arg H(0) = 0$. In the other case $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$, and we have

$$\lim_{t \rightarrow 0} H(t) = \lim_{\substack{\theta_2 \rightarrow \Theta_2^\pm(\theta_1^-) \\ \theta_2 \in \mathcal{R}^-}} \frac{\theta_2 - \Theta_2^\pm(\theta_1^-)}{\theta_2 - \Theta_2^\pm(\theta_1^-)} = -1.$$

The last equality is due to the fact that the tangent to \mathcal{R} at $\Theta_2^\pm(\theta_1^-)$ is vertical, see Figure 6. Indeed if we write $\theta_2 - \Theta_2^\pm(\theta_1^-) = a + ib$ when $\theta_2 \rightarrow \Theta_2^\pm(\theta_1^-)$ with $\theta_2 \in \mathcal{R}^-$, the vertical tangent gives $\frac{a}{b} \rightarrow 0$ and then $\lim_{t \rightarrow 0} H(t) = \lim_{\frac{a}{b} \rightarrow 0} \frac{\frac{a}{b} + i}{\frac{a}{b} - i} = -1$. It implies that $\delta = \pi$.

We are now going to show (32). We start by remarking that for $\theta_2 \in \mathcal{R}$, $G(\theta_2) = 1$ if and only if $\theta_2 \in \mathbb{R}$. Accordingly, for $t \in [0, 1]$, $H(t) = 1$ only at $t = 0$. Since $|H| = 1$ on $[0, 1]$, then necessarily $\delta + \Delta \in [-2\pi, 2\pi]$. We now calculate

$$H(1) = \lim_{\substack{\theta_2 \rightarrow \infty \\ \theta_2 \in \mathcal{R}^-}} G(\theta_2).$$

Thanks to Equation (7) we easily compute the following limit

$$\lim_{\substack{\theta_2 \rightarrow \infty \\ \theta_2 \in \mathcal{R}^-}} \frac{\Theta_1^-(\theta_2)}{\theta_2} = \frac{-\sigma_{12} - i\sqrt{\det \Sigma}}{\sigma_{11}}.$$

Using this limit together with the definition of G (see (12))

$$G(\theta_2) = \frac{\left(r_{11} \frac{\Theta_1^-(\theta_2)}{\theta_2} + r_{21}\right) \left(r_{12} \frac{\Theta_1^-(\theta_2)}{\theta_2} + r_{22}\right)}{\left(r_{12} \frac{\Theta_1^-(\theta_2)}{\theta_2} + r_{22}\right) \left(r_{11} \frac{\Theta_1^-(\theta_2)}{\theta_2} + r_{21}\right)},$$

we obtain that

$$\begin{aligned} H(1) &= \frac{\left(r_{11}(-\sigma_{12} - i\sqrt{\det \Sigma}) + r_{21}\sigma_{11}\right) \left(r_{12}(-\sigma_{12} + i\sqrt{\det \Sigma}) + r_{22}\sigma_{11}\right)}{\left(r_{12}(-\sigma_{12} - i\sqrt{\det \Sigma}) + r_{22}\sigma_{11}\right) \left(r_{11}(-\sigma_{12} + i\sqrt{\det \Sigma}) + r_{21}\sigma_{11}\right)} \\ &= \frac{\sigma_{22}r_{11}r_{12} + \sigma_{11}r_{22}r_{21} - \sigma_{12}(r_{11}r_{22} + r_{12}r_{21}) - i \det R \det \Sigma}{\sigma_{22}r_{11}r_{12} + \sigma_{11}r_{22}r_{21} - \sigma_{12}(r_{11}r_{22} + r_{12}r_{21}) + i \det R \det \Sigma} = \exp(i(\delta + \Delta)). \end{aligned}$$

Remembering that $\arg H(1) = \delta + \Delta \in [-2\pi, 2\pi]$, it gives (32) and clearly, $\delta + \Delta$ cannot be equal to $\pm 2\pi$ because $\det R \cdot \det \Sigma \neq 0$. \square

We now prove that

$$\chi = \begin{cases} 0 & \text{if } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \leq 0, \\ -1 & \text{if } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) > 0. \end{cases} \quad (33)$$

In particular, χ is an intrinsically non-continuous function of the parameters, see also Remark 7.

Recall that the function ψ_1 has been defined in (27).

Lemma 14. *The index χ can take only the values 0 and -1 , and we have the dichotomy:*

- $\chi = 0 \iff \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \leq 0 \iff \psi_1$ has no pole at infinity,
- $\chi = -1 \iff \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) > 0 \iff \psi_1$ has a simple pole at infinity.

Note that a simple pole at infinity means that $\psi_1(t) \sim c \cdot t$ for some non-zero constant c .

Proof. We have already seen in Proposition 6 that the sign of $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-))$ determines whether φ_1 has a pole in $\mathcal{G}\mathcal{R}$ or not, and thus if ψ_1 has a pole at infinity by the correspondence of Figure 6. This shows the two equivalences on the right in the statement of Lemma 14. We are thus left with proving the first two equivalent conditions.

First, if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$, $\delta = \pi$ and we have seen that in this case $H(t) \neq 1$ for all $t \in [0, 1]$. By (30), we deduce that $\chi = 0$.

If now $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \neq 0$ we notice that

$$\chi = \left\lfloor \frac{\Delta}{2\pi} \right\rfloor = 0 \text{ or } -1.$$

Indeed, we have proved in Lemma 13 that $\Delta \in (-2\pi, 2\pi)$. In particular, the sign of Δ determines χ : if $\text{sgn } \Delta \geq 0$ then $\chi = 0$ and if $\text{sgn } \Delta < 0$, $\chi = -1$. In the rest of the proof, we show that $\text{sgn } \Delta = -\text{sgn } \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-))$. First, Δ can be computed as

$$\Delta = \arg H(1) = [\arg G]_{\mathcal{R}^-} = \left[\arg \frac{\gamma_1 \overline{\gamma_2}}{\gamma_2 \overline{\gamma_1}}(\Theta_1^-(\theta_2), \theta_2) \right]_{\mathcal{R}^-}.$$

Let $\theta_2 = a - ib \in \mathcal{R}^-$ (we must have $b \geq 0$ and $a \geq \Theta_2^\pm(\theta_1^-) > 0$, see Figure 5) and $\theta_1 = \Theta_1^-(\theta_2) \in (-\infty, \theta_1^-]$. Using the expression (6) of γ_1 and γ_2 , we obtain

$$\gamma_1 \overline{\gamma_2}(\theta_1, \theta_2) = \gamma_1(\theta_1, a) \gamma_2(\theta_1, a) + r_{21} r_{22} b^2 + ib \theta_1 \det R,$$

from where we deduce that

$$\arg \frac{\gamma_1 \overline{\gamma_2}}{\gamma_2 \overline{\gamma_1}}(\theta_1, \theta_2) = 2 \arctan \frac{b \cdot \theta_1 \cdot \det R}{\gamma_1(\theta_1, a) \gamma_2(\theta_1, a) + r_{21} r_{22} b^2}. \quad (34)$$

We now look for the sign of (34) when $\theta_2 \rightarrow \Theta_2^\pm(\theta_1^-)$, while remaining in \mathcal{R}^- . This is sufficient to give the sign of Δ , because (34) does not change sign on \mathcal{R}^- due to the fact that $G(\theta_2) = 1$ on \mathcal{R}^- if and only if $\theta_2 = \Theta_2^\pm(\theta_1^-)$.

When $\theta_2 \rightarrow \Theta_2^\pm(\theta_1^-)$ we have $b \rightarrow 0$, $a \rightarrow \Theta_2^\pm(\theta_1^-)$ and $\theta_1 \rightarrow \theta_1^-$. We thus have

$$\begin{aligned} \operatorname{sgn} \arg \frac{\gamma_1 \bar{\gamma}_2}{\gamma_2 \bar{\gamma}_1}(\theta_1, \theta_2) &= \operatorname{sgn} b \cdot \operatorname{sgn} \theta_1 \cdot \operatorname{sgn} \det R \cdot \operatorname{sgn} \gamma_2(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \cdot \operatorname{sgn} \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \\ &= (+1)(-1)(+1)(+1) \operatorname{sgn} \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \\ &= -\operatorname{sgn} \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)), \end{aligned}$$

because $b \geq 0$, $\theta_1 < 0$, $\det R > 0$ by (2), and $\gamma_2(\theta_1^-, \Theta_2^\pm(\theta_1^-)) > 0$ (because $r_{22} > 0$ and $r_{22}\mu_1 - r_{12}\mu_2 < 0$, see Figure 3 to visualize this geometric condition). Then $\operatorname{sgn} \Delta = \operatorname{sgn} \arg \frac{\gamma_1 \bar{\gamma}_2}{\gamma_2 \bar{\gamma}_1}(\theta_1, \theta_2) = -\operatorname{sgn} \gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-))$, concluding the proof. \square

3.5. Resolution of the BVP. We are now in position to conclude the proof of Theorem 1. Reformulating Proposition 11, the function ψ_1 in (27)

- is sectionally analytic in $\mathbb{C} \setminus [0, 1]$,
- is continuous on $[0, 1]$ from below (with limits ψ_1^-) and above (with limits ψ_1^+),
 - is bounded at the vicinities of $[0, 1]$ if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \neq 0$,
 - has a pole of order one at 0 and bounded at 1 if $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$,
- is bounded at infinity if there is no pole before $\Theta_2^\pm(\theta_1^-)$ (then taking the value $\varphi_1(q)$), and with a pole of order one (see Lemma 3) at infinity if not (in short, it has a pole of order $-\chi$ at infinity),
- satisfies $\psi_1^+(t) = H(t)\psi_1^-(t)$ for $t \in [0, 1]$, with index χ given by (33), cf. also Lemma 14.

End of proof of Theorem 1. Our main reference for the resolution of the above so-called homogeneous BVP on an open contour is the book [45] of Muskhelishvili, see in particular [45, §79].

First of all, we prove that there exists a non-zero constant c such that

$$\psi_1(t) = c(t-1)^{-\chi} \exp \Gamma(t), \quad (35)$$

where Γ is the following function, sectionally analytic on $\mathbb{C} \setminus [0, 1]$:

$$\Gamma(t) = \frac{1}{2i\pi} \int_0^1 \frac{\log H(z)}{z-t} dz. \quad (36)$$

To make precise the definition (36), we define $\log H(z)$ by the facts that it should vary continuously over $[0, 1]$, and its initial value is such that $\log H(0) = i\delta$ (i.e., 0 if $H(0) = 1$ and $i\pi$ if $H(0) = -1$, see (31)).

At the vicinities 0 and 1, we have by [45, §29 and §79] that

$$\exp \Gamma(t) = t^{-\frac{\delta}{2\pi}} \Omega_0(t), \quad \exp \Gamma(t) = (t-1)^{\frac{\delta+\Delta}{2\pi}} \Omega_1(t),$$

for some function Ω_0 (resp. Ω_1) analytic in a neighborhood of 0 (resp. 1) and non-zero at 0 (resp. 1). Then we set

$$X(t) = t^{\frac{\delta}{\pi}} (t-1)^{-\chi} \exp \Gamma(t). \quad (37)$$

The function X in (37) is sectionally analytic in $\mathbb{C} \setminus [0, 1]$, and by construction of Γ and the Sokhotski-Plemelj formulas, it satisfies the boundary condition (29) (see [45, §79] for more details). Furthermore it has a pole of order $-\chi + \frac{\delta}{\pi}$ at infinity and is bounded at 0 and 1: indeed, $\frac{\delta}{\pi}$ and $-\chi$ are both equal to 0 or 1, see Lemmas 13 and 14. Then we consider two cases separately.

- First case: $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) \neq 0$. Then $\delta = 0$, and the function X in (37) simplifies into

$$X(t) = (t-1)^{-\chi} \exp \Gamma(t).$$

It satisfies the exact same boundary condition as (29). Looking at the ratio $\frac{\psi_1}{X}$, the boundary condition (29) gives that on $[0, 1]$,

$$\frac{\psi_1^+}{X^+} = \frac{\psi_1^-}{X^-}.$$

The above ratio is then analytic in the entire plane, even at the vicinities 0 and 1. The point 0 is a regular point and 1 is a removable singularity. Indeed, 1 is an isolated singular point, at which

$\frac{\psi_1}{X}$ might be infinite with degree less than unity (namely, $-\chi + \frac{\Delta}{2\pi}$). Moreover, the function $\frac{\psi_1}{X}$ is bounded at infinity, because both X and ψ_1 have a pole of the same order $-\chi$. Thanks to Liouville's theorem, we deduce that $\frac{\psi_1}{X}$ is constant. In conclusion, the formula (35) holds in this case.

• Second case: $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$. Then $\delta = \pi$, $\chi = 0$ and $X(t) = t \exp \Gamma(t)$ in (37). Firstly, we notice that the function $t\psi_1$ satisfies the boundary condition (29), is bounded at 0 and 1, and has a pole of order one at infinity. Moreover, the function X has a pole of order 1 at infinity. Considering then the ratio $\frac{t\psi_1}{X}$, the boundary condition (29) implies that on $[0, 1]$,

$$\frac{t\psi_1^+}{X^+} = \frac{t\psi_1^-}{X^-}.$$

The above ratio function is thus analytic in the entire complex plane, including the vicinities 0 and 1. It is indeed bounded at 1, and has a removable singularity at 0: the point 0 is an isolated singular point, at which $\frac{t\psi_1}{X}$ might be infinite with degree less than $\frac{1}{2}$. Using again Liouville's theorem, we deduce that the function $\frac{t\psi_1}{X}$ is a constant. Formula (35) therefore also holds.

We now deduce from (35) the formula (14) stated in Theorem 1. Going from the t -plane back to the θ_2 -plane (see (27) and Figure 6), one has that for some constant c ,

$$\varphi_1(\theta_2) = \psi_1(W(\theta_2)) = c(W(\theta_2) - 1)^{-\chi} \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}^-} \log G(\theta) \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \right\}. \quad (38)$$

Using the equation (26) relating W and w , we easily obtain

$$W(\theta_2) - 1 = \frac{w(q) + 1}{w(\theta_2) - w(q)}$$

as well as

$$\frac{W'(\theta)}{W(\theta) - W(\theta_2)} = \frac{w'(\theta)}{w(\theta) - w(\theta_2)} - \frac{w'(\theta)}{w(\theta) - w(q)}.$$

Remembering that in the case $\chi = -1$ one has $q = p$, see (24), we finally obtain that for some constant c' ,

$$\varphi_1(\theta_2) = c' \left(\frac{1}{w(\theta_2) - w(p)} \right)^{-\chi} \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}^-} \log G(\theta) \frac{w'(\theta)}{w(\theta) - w(\theta_2)} d\theta \right\}.$$

By definition (4) of the Laplace transform we have $\varphi_1(0) = \nu_1(\mathbb{R}_+)$. To find the exact constant c' (and thereby our main result (14)), we simply evaluate the above formula at $\theta_2 = 0$ and use Lemma 15 below. \square

Lemma 15. *One has* $\begin{pmatrix} \varphi_1(0) \\ \varphi_2(0) \end{pmatrix} = \begin{pmatrix} \nu_1(\mathbb{R}_+) \\ \nu_2(\mathbb{R}_+) \end{pmatrix} = -R^{-1}\mu$.

Proof. Equation (4) evaluated at 0 gives $\varphi_i(0) = \nu_i(\mathbb{R}_+)$. We now evaluate (5) at $\theta_2 = 0$, divide by θ_1 and finally evaluate at $\theta_1 = 0$. This yields $-\mu_1 = r_{11}\varphi_1(0) + r_{12}\varphi_2(0)$. In a similar way, we obtain $-\mu_2 = r_{21}\varphi_1(0) + r_{22}\varphi_2(0)$, thereby concluding the proof. \square

Note that Lemma 15 gives, as announced in Theorem 1,

$$\nu_1(\mathbb{R}_+) = \frac{r_{12}\mu_2 - r_{22}\mu_1}{\det R}, \quad \nu_2(\mathbb{R}_+) = \frac{r_{21}\mu_1 - r_{11}\mu_2}{\det R},$$

which by (2) are positive: remind that this positivity condition is necessary for the existence of the stationary distribution [33].

Clearly, the integral expression (14) of φ_1 is meromorphic in the domain $\mathcal{G}_{\mathcal{R}}$. In Section 4 we shall see that it can be meromorphically continued on the much larger domain $\mathbb{C} \setminus [\theta_2^+, \infty)$.

To conclude this part, let us make Remark 12 more precise. In the case $\gamma_1(\theta_1^-, \Theta_2^\pm(\theta_1^-)) = 0$ (i.e., $H(0) = -1$), we have chosen $\delta = \pi$ in (30). (Recall that choosing δ is tantamount to fixing

a determination of the arg (or log) function.) Remarkably, any other choice of δ would have led to the same explicit expression for φ_1 , though written differently. For instance, if we had taken $\delta = -\pi$ instead, the index χ would have been -1 (instead of 0), and the two determinations of the logarithm would differ by $-2i\pi$. With our notation in the proof of Theorem 1, we would have obtained for some constant c'' ,

$$\psi_1(t) = c'' \frac{t-1}{t} \exp \Gamma(t).$$

This actually corresponds to the formula of Theorem 1 associated with $\chi = -1$.

3.6. Generalizations. Our main Theorem 1 is derived under the hypothesis (13) that the coordinates (μ_1, μ_2) of the drift are negative. However, the conditions (2) (equivalent to the existence of a stationary distribution) allow the drift to have one non-negative coordinate. In the next few lines, we assume that $\mu_1 \geq 0$ or $\mu_2 \geq 0$, and we comment on the slight differences which would have arisen in the analytic treatment of the functional equation (5).

In the case of a drift having one non-negative coordinate, the reflected Brownian motion in the quadrant has a pathwise behavior which is quite different (it spends most of its time near the axis towards which the drift is directed). However, from the point of view of our analysis, the differences are only technical. For instance, the continuation formula (16) of Lemma 3 may use the other branch Θ_1^+ , as if $\mu_1 > 0$, the small branch (taking value 0 at 0) is Θ_1^+ and not Θ_1^- anymore, see (7). As a result, the poles of φ_1 may be located at places different from those described in Proposition 6. Despite these few differences, there still exists a Cauchy-type integral expression for the Laplace transform φ_1 , similar to the one stated in Theorem 1.

Let us also very briefly mention here the case of reflected Brownian motion in higher dimension [31, 34, 7, 6, 15]. Compared to its two-dimensional analogue, much less is known. However, an analogue of the functional equation (5) can still be stated (since the BAR exists in any dimension, see [33, 14]). Clearly, our techniques (based on complex analysis) use the dimension 2, and in our opinion, generalizing in higher dimension these BVP techniques is a difficult open problem. In the discrete case too, the case of dimension 3 is less understood. One can mention [12] for some ideas to state a BVP in dimension 3, as well as [4] for more combinatorial techniques.

4. SINGULARITY ANALYSIS AND ASYMPTOTICS OF THE BOUNDARY DISTRIBUTION

The asymptotics (up to a constant) of the boundary measures has been obtained by Dai and Miyazawa in [18], see also [17, 28] for the interior measure. In this section we show that our expression (14) for the Laplace transform φ_1 stated in Theorem 1 is perfectly suited for singularity analysis, and accordingly to study the asymptotics (including the computation of the constant) of the boundary stationary distributions ν_1 and ν_2 . We shall first recall the result of [18] and express it in terms of our notations. Then we will explain how, thanks to Theorem 1, we could obtain a new proof of this result and make the constants explicit.

4.1. Asymptotic results.

Theorem 16 (Theorem 6.1 of [18]). *The following asymptotics holds:*

$$\lim_{x \rightarrow \infty} \frac{\nu_1(x)}{x^\kappa e^{-\tau_2 x}} = b, \tag{39}$$

where $\kappa \in \{-\frac{3}{2}, -\frac{1}{2}, 0, 1\}$ and $\tau_2 \in \{p, p', \theta_2^+\}$ are given in Table 1, and b is some positive constant. See Figures 7 and 8 to visualize geometrically the different cases.

We now propose a series of three remarks on Theorem 16.

- We have already introduced p , see (17), and p' is as follows: it is the (unique, when it exists) non-zero point $p' = \Theta_2^+(r)$, with r defined by $\gamma_2(r, \Theta_2^-(r)) = 0$ and $r \leq \Theta_1^+(\theta_2^+)$. It will be convenient to adopt the following notation: we will write $p > \theta_2^+$ (resp. $p' > \theta_2^+$) when p (resp. p') does not exist.

• Theorem 6.1 of [18] deals with the asymptotics of $\nu_1(x, \infty)$, and not with that of $\nu_1(x)$, as stated in Theorem 16 below. However, the two statements are equivalent: when the asymptotics of the density has the form $bx^\kappa e^{-\tau_2 x}$ as in (39), the corresponding tail probability is given by the exact same asymptotics (with another constant b), see [17, Lemma D.5].

• Before stating Table 1, which gives the values of κ and τ_2 of Theorem 16, we briefly recall Dai and Miyazawa's notations: $\theta^{(2,\max)}$ is the (unique) point of the ellipse $\gamma = 0$ such that $\theta_2^{(2,\max)} = \theta_2^+$, $\tau_2 = \sup\{\theta_2 > 0 : \exists \theta_1 \in \mathbb{R}, \varphi(\theta_1, \theta_2) < \infty\}$ and $\theta^{(2,r)}$ is the intersection point of the straight line $\gamma_1 = 0$ and the ellipse $\gamma = 0$. Notice that the definition of τ_2 does not rely on an analytic continuation of φ .

Cases		Dai and Miyazawa's categories		κ	τ_2
p or $p' \in (0, \theta_2^+)$	1.a	$p < \Theta_2^\pm(\theta_1^-)$	$\tau_2 < \theta_2^{(2,\max)}$	Categories I or II	
	1.b	$\Theta_2^\pm(\theta_1^-) \leq p < p'$		Category I	
	1.c	$\Theta_2^\pm(\theta_1^-) \leq p' < p$		Category III, $\tau_2 \neq \theta_2^{(2,r)}$	
	1.d	$\Theta_2^\pm(\theta_1^-) \leq p = p'$		Category III, $\tau_2 = \theta_2^{(2,r)}$	
p and $p' \geq \theta_2^+$	2.a	p and $p' > \theta_2^+$	$\tau_2 = \theta_2^{(2,\max)}$	Category I, $\theta^{(2,r)} \neq \theta^{(2,\max)}$	
	2.b	$\theta_2^+ = p$		Category I, $\theta^{(2,r)} = \theta^{(2,\max)}$	
	2.c	$\theta_2^+ = p'$		Category II, $\theta^{(2,r)} \neq \theta^{(2,\max)}$	
	2.d	$\theta_2^+ = p = p'$		Category II, $\theta^{(2,r)} = \theta^{(2,\max)}$	

TABLE 1. Dai and Miyazawa's categories expressed with our notations

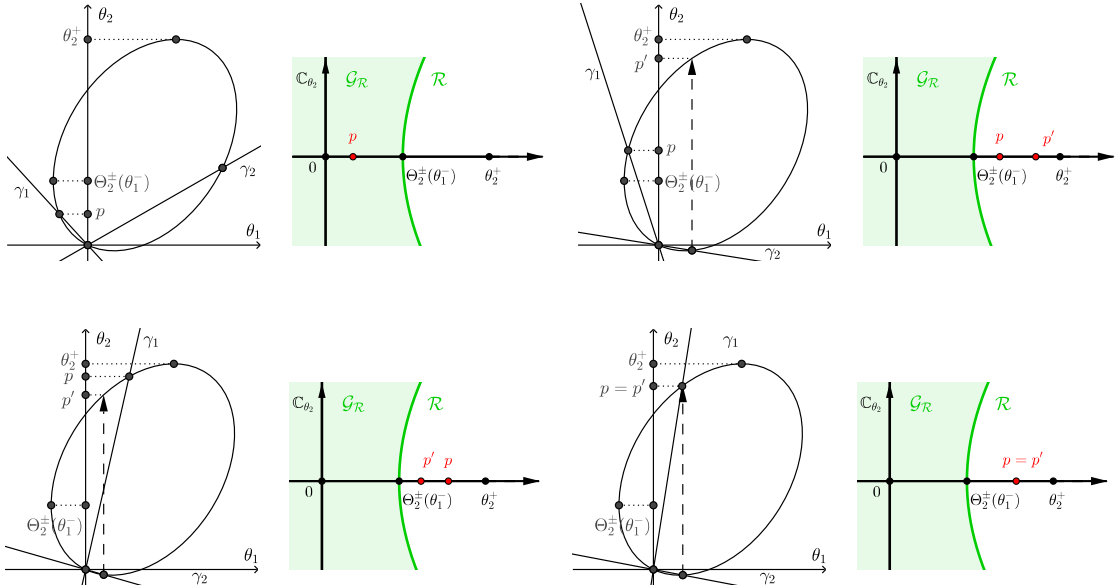


FIGURE 7. Ellipses and curves in the cases 1.a, 1.b, 1.c and 1.d of Table 1

4.2. Sketch of the proof of Theorem 16. The idea is to study the singularities of φ_1 and to use transfer theorems (such as [20, Theorem 37.1]) relating the asymptotics of a function and the singularities of its Laplace transform. In our case, the singularity closest to 0 will determine the asymptotics. Our aim here is not to propose a complete proof of Theorem 16 (for this we refer

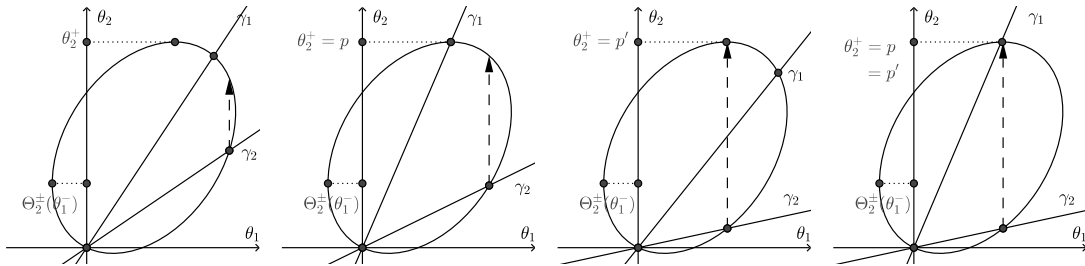


FIGURE 8. Ellipses in the cases 2.a, 2.b, 2.c and 2.d of Table 1

to [18]), but rather to illustrate that Theorem 1 indeed implies Theorem 16, with the additional feature of providing an exact expression for the constant b in (39). We give details for the case 1.a of Table 1, and only sketch the difficulties that would arise in the remaining cases.

Proof of Theorem 16 and computation of b in case 1.a. In that case, the index χ equals -1 and the formula (14) gives

$$\varphi_1(\theta_2) = \nu_1(\mathbb{R}_+) \left(\frac{w(0) - w(p)}{w(\theta_2) - w(p)} \right) \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}^-} \log G(\theta) \left[\frac{w'(\theta)}{w(\theta) - w(\theta_2)} - \frac{w'(\theta)}{w(\theta) - w(0)} \right] d\theta \right\}.$$

Recall that w is analytic, one-to-one on $\mathcal{G}_{\mathcal{R}}$ and further satisfies $w(\theta_2) \neq w(\theta)$, for all $\theta_2 \in \mathcal{G}_{\mathcal{R}}$ and $\theta \in \mathcal{R}^-$. As a first result, the integral part (thus also its exponential) is analytic in the domain $\mathcal{G}_{\mathcal{R}}$. A second consequence is that $\frac{1}{w(\theta_2) - w(p)}$ has a simple pole at p : $\varphi_1(\theta_2) = \frac{b+o(1)}{\theta_2 - p}$, with

$$b = \nu_1(\mathbb{R}_+) \left(\frac{w(0) - w(p)}{w'(p)} \right) \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}^-} \log G(\theta) \left[\frac{w'(\theta)}{w(\theta) - w(p)} - \frac{w'(\theta)}{w(\theta) - w(0)} \right] d\theta \right\}.$$

Theorem 37.1 of [20] gives the announced asymptotics $\nu_1(x) = e^{-px}(b + o(1))$ as $x \rightarrow \infty$. \square

In the other cases, the singularities are not in $\mathcal{G}_{\mathcal{R}}$ and we thus need to extend meromorphically φ_1 in a larger domain:

Proposition 17 (Theorem 11 of [28]). *The Laplace transform φ_1 can be meromorphically continued on the domain $\mathbb{C} \setminus [\theta_2^+, \infty)$.*

Proposition 17 has already been proved in [28], see Theorem 11 there. Note that the formula (14) of Theorem 1 provides an alternative, direct, analytic proof, which can be sketched as follows: the equation (14) is valid a priori only for θ_2 in $\mathcal{G}_{\mathcal{R}}$. However, considering a Hankel contour similar to that of Figure 9, surrounding $[\theta_2^+, \infty)$, we could write φ_1 as an integral over the cut $[\theta_2^+, \infty)$. The study of the so-obtained formula would give the singularities p, p' and θ_2^+ as in Table 1, and would lead to the precise asymptotics (we could even obtain the full asymptotic development) and the computation of b .

Remark 18. *We have already commented on the fact that within a single formula, Theorem 1 actually captures two different expressions, depending on the value of χ in (33). From an asymptotic viewpoint, Section 4 shows that different cases exist as well, depending on various parameters. It should be noted that these cases are all different, i.e., δ and χ do not govern the asymptotic behavior of the boundary measures.*

5. ALGEBRAIC NATURE AND SIMPLIFICATION OF THE LAPLACE TRANSFORMS

Motivations. In this section we are interested in the following question: in which extent is it possible to simplify the expressions of the Laplace transforms given in Theorem 1? For instance, is this possible that these functions be algebraic or even rational?

This is of paramount importance: first, simplified expressions would lead to an easier analysis, in particular for asymptotic analysis or for taking inverse Laplace transforms; second, understanding the parameters (Σ, μ, R) for which the Laplace transforms are rational should reveal intrinsic structure of the model. In the particular case of the identity covariance matrix Σ , some attempts of simplifications may be found in [26, Chapter 4].

In the literature, this question has received much interest in the discrete setting. One can first think at the famous Jackson's networks [36] and their product form solutions. In a closer context, Latouche and Miyazawa [40], Chen, Boucherie and Goseling [10], obtain geometric necessary and sufficient conditions for the stationary distribution of random walks in the quarter plane to be sums of geometric terms. Such criteria can be applied, e.g., to derive an approximation scheme to error bounds for performance measures of random walks in the quarter plane [11].

In our context of reflected diffusions in the quadrant, analogues of these results are obtained by Dieker and Moriarty [19]: a simple condition (involving the single angle (40)) for the stationary density to be a sum of exponential terms is derived. We will discuss the links with our results in Section 5.1. In the simplest case the sum of exponential terms consist of one single term: this is the skew-symmetric condition (23) of [34]. In this case, a factorization such as in Remark 8 exists and applying an invariant method we are able to solve the skew-symmetric case, see our Section 5.2, yielding new proofs to already known results.

Another main reason which can lead to simplified expressions comes from the rationality of $\frac{\beta}{\pi}$, with β in (25). Before being more precise, let us mention that this reason is deeply different than the first one: β only depends on the covariance matrix and is therefore independent of the reflection matrix, while Dieker and Moriarty's condition is also dependent on the reflection angles, see (40). In the discrete setting, the rationality of $\frac{\beta}{\pi}$ is rather interpreted as a condition on the finiteness of a certain group of transformations, and this finiteness was shown to have a decisive influence on the D-finiteness (a function is D-finite if it satisfies a linear differential equation with polynomial coefficients on \mathbb{Q}) of the generating functions, see [25, 5, 21]. For reflected Brownian motion in the quadrant with orthogonal reflections, it is shown in [29] that the Laplace transform is algebraic (note, any algebraic function is D-finite) if and only if the group is finite. We present a structural result in Section 5.3.

Finally, we focus in Section 5.4 on the case of orthogonal reflections, and derive a new proof of the main result of [29], as a consequence of Theorem 1.

5.1. Dieker and Moriarty's criterion. For the sake of completeness, let us mention the following result:

Theorem 19 (Theorem 1 in [19]). *The stationary density is a sum of exponentials if and only if*

$$\frac{\varepsilon + \delta - \pi}{\beta} \in -\mathbb{N} = -\{0, 1, 2, \dots\}, \quad (40)$$

with ε and δ in $(0, \pi)$ and

$$\tan \varepsilon = \frac{\sin \beta}{\frac{r_{21}}{r_{11}} \sqrt{\frac{\sigma_{11}}{\sigma_{22}}} + \cos \beta}, \quad \tan \delta = \frac{\sin \beta}{\frac{r_{12}}{r_{22}} \sqrt{\frac{\sigma_{22}}{\sigma_{11}}} + \cos \beta}. \quad (41)$$

Notice that this is not the exact statement of [19, Theorem 1], as in Dieker and Moriarty's paper, the Brownian motion is assumed to have an identity covariance matrix and evolves in a wedge with an arbitrary opening angle, whereas we consider Brownian motion with arbitrary covariance matrix but in the quarter plane. A simple linear transform, which is made explicit in Appendix A, makes both statements equivalent. The expression (41) of the angles ε and δ follows from this transformation.

5.2. Skew-symmetric case. The skew-symmetric case holds when the matrix condition (23) is satisfied [34]. In dimension two, (23) can be reduced to the single equation

$$2\sigma_{12} = \frac{r_{21}}{r_{11}}\sigma_{11} + \frac{r_{12}}{r_{22}}\sigma_{22}. \quad (42)$$

Using the identities in (49) we easily show that the above condition is equivalent to

$$\varepsilon + \delta = \pi,$$

which is the case where the quantity of Dieker and Moriarty's criterion is equal to 0, see (40).

It is known, see [34], that condition (23) is satisfied if and only if the stationary distribution has a product form, i.e., $\pi(x_1, x_2) = \pi_2(x_1)\pi_1(x_2)$, where the π_i 's are the marginal densities of π . Furthermore it implies that the stationary distribution is exponential, meaning that

$$\pi(x_1, x_2) = \alpha_1\alpha_2 e^{-\alpha_1 x_1 - \alpha_2 x_2}, \quad \text{with } \alpha = -2 \cdot \text{diag}(\Sigma)^{-1} \cdot \text{diag}(R) \cdot R^{-1} \cdot \mu. \quad (43)$$

See for example [31, §10], [34] or [18] for more details on these results.

In fact, in the skew-symmetric case, a rational factorization such as in the Remark 8 exists. It is then possible to find again, in another way, some already known results. Indeed if the skew-symmetric condition holds, we shall prove below that for $\theta_2 \in \mathcal{R}$,

$$G(\theta_2) = \frac{F(\theta_2)}{F(\overline{\theta_2})}, \quad \text{with } F(\theta_2) = \alpha_2 - \theta_2. \quad (44)$$

where by (43) α_2 takes the value

$$\alpha_2 = \frac{2r_{22}(r_{21}\mu_1 - r_{11}\mu_2)}{\sigma_{22} \det R}. \quad (45)$$

Using (42) and (18) it is easy to remark that $\alpha_2 = p$, which is the only possible pole of φ_1 in $\mathcal{G}_{\mathcal{R}}$ as we observe after Lemma 3. It follows from Proposition 6 that the function $\varphi_1(\theta_2)(\alpha_2 - \theta_2)$

- is bounded on $\mathcal{G}_{\mathcal{R}}$ and converges at infinity to $\nu_1(0)$ (thanks to the initial value theorem and Lemma 3),
- is continuous on $\overline{\mathcal{G}_{\mathcal{R}}}$,
- satisfies the boundary condition $\varphi_1(\theta_2)(\alpha_2 - \theta_2) = \varphi_1(\overline{\theta_2})(\alpha_2 - \overline{\theta_2})$ for all $\theta_2 \in \mathcal{R}$.

Then, using an invariant lemma (see Lemma 2 in [43, Section 10.2]), we conclude that for some constant C ,

$$\varphi_1(\theta_2) = \frac{C}{\alpha_2 - \theta_2}.$$

Evaluating the above equation at 0 and using Lemma 15 gives $C = \frac{\sigma_{11}}{2r_{11}}\alpha_1\alpha_2$. Inverting the Laplace transform implies that the stationary distribution is exponential. Then using the functional equation (5), we find that $\frac{\sigma_{11}\sigma_{22}}{4r_{11}r_{22}}\alpha_1\alpha_2\varphi(\theta_1, \theta_2) = \varphi_1(\theta_2)\varphi_2(\theta_1)$, which means that π has a product form.

Proof of Equation (44). Let $\theta_2 \in \mathcal{R}$ and note $\theta_1 = \Theta_1^-(\theta_2)$. Elementary computations give

$$\gamma_1(\theta_1, \theta_2)\gamma_2(\theta_1, \overline{\theta_2}) = \theta_1 \left(\frac{2r_{22}}{\sigma_{22}}(r_{21}\mu_1 - r_{11}\mu_2) - \theta_2 \det R \right).$$

To find this, we just have to use the skew-symmetric condition (42) and to remark that for $\theta_2 \in \mathcal{R}$ Vieta's formulas give $\theta_2\overline{\theta_2} = \frac{1}{\sigma_{22}}(\sigma_{11}\theta_1^2 + 2\mu_1\theta_1)$ and $\overline{\theta_2} = -\theta_2 - \frac{1}{\sigma_{22}}(2\sigma_{12}\theta_1 + 2\mu_2)$. Then, using the expression (45) of α_2 we find $G(\theta_2) = \frac{\alpha_2 - \theta_2}{\alpha_2 - \theta_2}$ for $\theta_2 \in \mathcal{R}$. \square

5.3. Structural form of the Laplace transforms. In the case $\frac{\beta}{\pi} \in \mathbb{Q}$, and in this case only, the function W in (26) is algebraic (as the generalized Chebyshev polynomial (10) is, see Remark 10), yielding the following structural result:

Proposition 20. *If $\frac{\beta}{\pi} \in \mathbb{Q}$, the Laplace transform φ_1 of Theorem 1 is the product of an algebraic function by the exponential of a D-finite function.*

Proof. This easily follows from the fact that the Cauchy integral of a D-finite function is D-finite, see, e.g., [51]. \square

However, it is not true in general that the exponential of a D-finite function is still D-finite.

5.4. Orthogonal reflections. Here we consider the case of orthogonal reflections, which is equivalent for the reflection matrix R in (1) to be the identity matrix; see also Figure 2. By developing the theory of Tutte's invariants (introduced by Tutte in [53] for the enumeration of properly colored triangulations, and used in [3] for the enumeration of quadrant walks) for the Brownian motion, we proved in [29] the following result:

Theorem 21 (Theorem 1 in [29]). *Let R be the identity matrix in (1). The Laplace transform φ_1 is equal to*

$$\varphi_1(\theta_2) = \frac{-\mu_1 w'(0)}{w(\theta_2) - w(0)} \theta_2. \quad (46)$$

In this section we derive a new proof of this result, as a consequence of Theorem 1. More generally, the proof below would work for any parameters (Σ, μ, R) such that $G(\theta_2) = \frac{F(\theta_2)}{F(\theta_2)}$ (cf. Remark 8), yielding a rational expression of $\varphi_1(\theta_2)$ in terms of $w(\theta_2)$ and θ_2 .

Proof. Let us first notice that the index $\chi = 0$. Indeed, since $G(\theta_2) = \frac{\bar{\theta}_2}{\theta_2}$, we have $\arg G(\theta_2) = -2 \arg \theta_2 > 0$ for $\theta_2 \in \mathcal{R}$, and thus $\Delta > 0$, see the proof of Lemma 14.

Starting from the formula (38), we have for some constant C

$$\varphi_1(\theta_2) = C \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}^-} (\log \bar{\theta} - \log \theta) \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \right\} = C \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}} \log \theta \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \right\}.$$

To compute the above integral, we first integrate on the contour represented on Figure 9. The residue theorem gives

$$\frac{1}{2i\pi} \left\{ \int_{\mathcal{R}_R} + \int_{\mathcal{C}_R} + \int_{\mathcal{C}_\varepsilon} + \int_{-R+i\varepsilon}^{i\varepsilon} + \int_{-i\varepsilon}^{-R-i\varepsilon} \right\} \log \theta \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta = \log \theta_2 - \log q. \quad (47)$$

It is easy to see that

$$\lim_{R \rightarrow \infty} \frac{1}{2i\pi} \int_{\mathcal{R}_R} \log \theta \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta = \frac{1}{2i\pi} \int_{\mathcal{R}} \log \theta \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta$$

and that in the limits when $\varepsilon \rightarrow 0$ and $R \rightarrow \infty$, the contributions on \mathcal{C}_ε and \mathcal{C}_R both converge to 0, because W is analytic at 0 and ∞ , respectively.

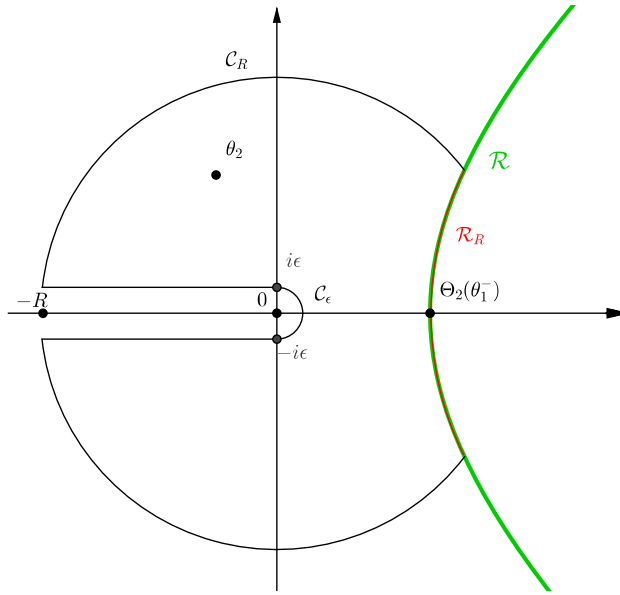


FIGURE 9. Integration contour used in the proof of Theorem 21

Furthermore,

$$\begin{aligned}
\lim_{\varepsilon \rightarrow 0} \lim_{R \rightarrow \infty} \frac{1}{2i\pi} \left\{ \int_{-R+i\varepsilon}^{i\varepsilon} + \int_{-i\varepsilon}^{-R-i\varepsilon} \right\} \log \theta \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \\
&= \lim_{\varepsilon \rightarrow 0} \frac{1}{2i\pi} \int_{-\infty}^0 (\log(t+i\varepsilon) - \log(t-i\varepsilon)) \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \\
&= \frac{1}{2i\pi} \int_{-\infty}^0 2i\pi \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \\
&= \log \frac{W(0) - W(\theta_2)}{W(\infty) - W(\theta_2)} \\
&= \log \frac{w(\theta_2) - w(0)}{w(q) - w(0)}.
\end{aligned}$$

Above we have used the dominated convergence and the fact that the principal determination of the logarithm gives us $\lim_{\varepsilon \rightarrow 0} (\log(t+i\varepsilon) - \log(t-i\varepsilon)) = 2i\pi$. Letting $R \rightarrow \infty$ and then $\varepsilon \rightarrow 0$ in (47), we have, for some constants C and C' ,

$$\varphi_1(\theta_2) = C \exp \left\{ \frac{1}{2i\pi} \int_{\mathcal{R}} \log \theta \frac{W'(\theta)}{W(\theta) - W(\theta_2)} d\theta \right\} = C' \frac{\theta_2}{w(\theta_2) - w(0)}.$$

Since by Lemma 15 one has $\varphi_1(0) = -\mu_1$, this eventually gives the right constant in (46). \square

APPENDIX A. EQUIVALENCE BETWEEN BROWNIAN MOTION IN WEDGES AND BROWNIAN MOTION IN THE QUARTER PLANE

We use the notation of Section 1. Up to an isomorphism, studying Brownian motion in the quarter plane with arbitrary covariance matrix Σ is equivalent to studying Brownian motion in a cone of angle $\beta = \arccos -\frac{\sigma_{12}}{\sqrt{\sigma_{11}\sigma_{22}}}$, with covariance identity. See for example [1, Equation (23)] and [49, Lemma 3.23]. In this short section we relate the key parameters (angles of the reflection vectors and drift) before and after the linear transformation.

Let us define the linear transforms

$$T = \begin{pmatrix} \frac{1}{\sin \beta} & \cot \beta \\ 0 & 1 \end{pmatrix} \begin{pmatrix} \frac{1}{\sqrt{\sigma_{11}}} & 0 \\ 0 & \frac{1}{\sqrt{\sigma_{22}}} \end{pmatrix}, \quad T^{-1} = \begin{pmatrix} \sqrt{\sigma_{11}} & 0 \\ 0 & \sqrt{\sigma_{22}} \end{pmatrix} \begin{pmatrix} \sin \beta & -\cos \beta \\ 0 & 1 \end{pmatrix}. \quad (48)$$

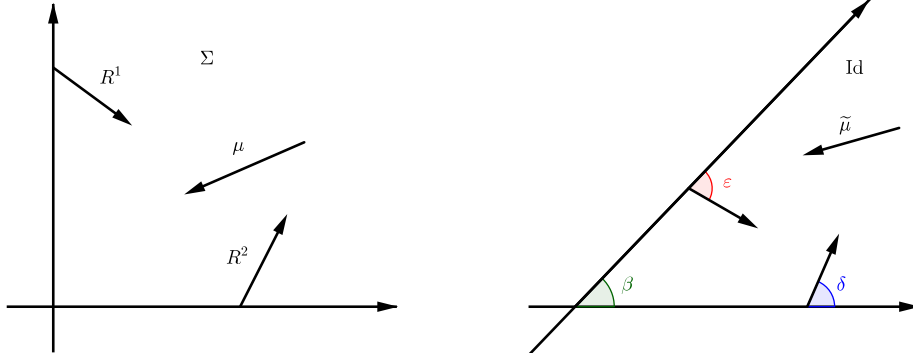


FIGURE 10. The linear transformation T in (48) from the quadrant to the wedge of opening angle β

Obviously the reflected Brownian motion associated to (Σ, μ, R) becomes a Brownian motion (with covariance identity) in a wedge of angle β and with parameters $(\text{Id}, T\mu, TR)$. The new angles of reflection are δ and ε (cf. Figure 10), such that

$$\begin{cases} \tan \delta = \frac{\sin \beta}{a + \cos \beta}, & \cos \delta = \frac{a + \cos \beta}{\sqrt{a^2 + 2a \cos \beta + 1}}, & \sin \delta = \frac{\sin \beta}{\sqrt{a^2 + 2a \cos \beta + 1}}, \\ \tan \varepsilon = \frac{\sin \beta}{b + \cos \beta}, & \cos \varepsilon = \frac{b + \cos \beta}{\sqrt{b^2 + 2b \cos \beta + 1}}, & \sin \varepsilon = \frac{\sin \beta}{\sqrt{b^2 + 2b \cos \beta + 1}}, \end{cases} \quad (49)$$

where $a = \frac{r_{12}}{r_{22}} \sqrt{\frac{\sigma_{22}}{\sigma_{11}}}$ and $b = \frac{r_{21}}{r_{11}} \sqrt{\frac{\sigma_{11}}{\sigma_{22}}}$. The new drift is $\tilde{\mu} = T\mu$, where

$$\tilde{\mu}_1 = \frac{\mu_1}{\sqrt{\sigma_{11}}} \frac{1}{\sin \beta} + \frac{\mu_2}{\sqrt{\sigma_{22}}} \cot \beta \quad \text{and} \quad \tilde{\mu}_2 = \frac{\mu_2}{\sqrt{\sigma_{22}}}.$$

To conclude the appendix, we prove Corollary 2, which gives an explicit formula for the Laplace transform of the stationary distribution of a reflected Brownian motion in a wedge.

Proof of Corollary 2. Let us arbitrarily choose $\sigma_{11} = \sigma_{22} = 1$ and $\sigma_{12} = -\cos \beta$. It implies that the linear transform T in (48) is the same as T_1 in the statement of Corollary 2. We consider the process $Z = T^{-1}\tilde{Z}$, which is a reflected Brownian motion in the quadrant of parameters (Σ, μ, R) , with

$$\Sigma = T^{-1}\tilde{\Sigma}(T^{-1})^\top, \quad \mu = T^{-1}\tilde{\mu} \quad \text{and} \quad R = T^{-1}\tilde{R}.$$

Let Π (resp. $\tilde{\Pi}$) be the invariant measure of Z (resp. \tilde{Z}) and π (resp. $\tilde{\pi}$) its density. It is easy to notice that

$$\Pi = \tilde{\Pi} \circ T \quad \text{and} \quad \pi = |\det T| \tilde{\pi} \circ T.$$

Indeed, by a fundamental property of the invariant measure, for all $x \in \mathbb{R}_+^2$ and all measurable set A in \mathbb{R}_+^2 we have the following limits

$$\mathbb{P}_x[Z_t \in A] \xrightarrow{t \rightarrow \infty} \Pi(A), \quad \mathbb{P}_x[Z_t \in A] = \mathbb{P}_x[TZ_t \in TA] = \mathbb{P}_{Tx}[\tilde{Z}_t \in TA] \xrightarrow{t \rightarrow \infty} \tilde{\Pi}(TA).$$

It yields $\Pi = \tilde{\Pi} \circ T$. Furthermore by a simple change of variable we have

$$\tilde{\Pi} \circ T(A) = \int_{TA} \tilde{\pi}(\tilde{x}) d\tilde{x} = |\det T| \int_A \tilde{\pi}(Tx) dx,$$

therefore $|\det T|\tilde{\pi} \circ T$ is the density of $\Pi \circ T$ and is then equal to π . Theorem 1 gives the value of φ , the Laplace transform of π . Lastly, a simple change of variable $\tilde{x} = Tx$ in Equation (3) leads to

$$\tilde{\varphi}(\tilde{\theta}) = \iint_{\mathcal{C}_\beta} \exp(\tilde{\theta} \cdot \tilde{x}) \tilde{\pi}(\tilde{x}) d\tilde{x} = \iint_{T^{-1}\mathcal{C}_\beta} \exp(\tilde{\theta} \cdot Tx) \tilde{\pi}(Tx) |\det T| dx = \iint_{\mathbb{R}_+^2} \exp(T^\top \tilde{\theta} \cdot x) \pi(x) dx = \varphi(T^\top \tilde{\theta}),$$

where \mathcal{C}_β is the wedge of angle β where the process evolves. □

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