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RADIAL DUNKL PROCESSES : EXISTENCE AND UNIQUENESS, HITTING TIME, BETA PROCESSES AND RANDOM MATRICES

NIZAR DEMNI¹

ABSTRACT. We begin with the study of some properties of the radial Dunkl process associated to a reduced root system R . It is shown that this diffusion is the unique strong solution for all $t \geq 0$ of a SDE with singular drift. Then, we study T_0 , the first hitting time of the positive Weyl chamber : we prove, via stochastic calculus, a result already obtained by Chybiryakov on the finiteness of T_0 . The second and new part deals with the law of T_0 for which we compute the tail distribution, as well as some insight via stochastic calculus on how root systems are connected with eigenvalues of standard matrix-valued processes. This gives rise to the so-called β -processes. The ultraspherical β -Jacobi case still involves a reduced root system while the general case is closely connected to a non reduced one. This process lives in a convex bounded domain known as principal Weyl alcove and the strong uniqueness result remains valid. The last part deals with the first hitting time of the alcove's boundary and the semi group density which enables us to answer some open questions.

1. PRELIMINARIES

We begin by pointing out some facts on root systems and radial Dunkl processes. We refer to [38] for the Dunkl theory, to both [7] and [27] for a background on root systems and [11], [21] for facts on radial Dunkl processes. Let (V, \langle, \rangle) be a finite real Euclidean space of dimension m . A *reduced* root system R is a finite set of non zero vectors spanning V such that :

$$1 \quad R \cap \mathbb{R}\alpha = \{\alpha, -\alpha\} \text{ for all } \alpha \in R.$$

$$2 \quad \sigma_\alpha(R) = R$$

where σ_α is the reflection with respect to the hyperplane H_α orthogonal to α :

$$\sigma_\alpha(x) = x - 2 \frac{\langle \alpha, x \rangle}{\langle \alpha, \alpha \rangle} \alpha, \quad x \in V$$

A simple system Δ is a basis of V which induces a total ordering in R . A root α is positive if it is a positive linear combination of elements of Δ . The set of positive roots is called a positive subsystem and is denoted by R_+ . Note that the choice of Δ is not unique and that R_+ is uniquely determined by Δ . The reflection group

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W is the one generated by all the reflections σ_α for $\alpha \in R$. Recall that W is finite and the only reflections are of the form σ_α for $\alpha \in R$. Given a root system R with associated positive subsystem R_+ , let C be the *positive Weyl chamber* defined by :

$$C := \{x \in V \mid \langle \alpha, x \rangle > 0 \forall \alpha \in R_+\} = \{x \in V \mid \langle \alpha, x \rangle > 0 \forall \alpha \in \Delta\}$$

and \overline{C} its closure. One of the most important properties is that the convex cone \overline{C} is a fundamental domain, that is each $\lambda \in V$ is conjugate to one and only one $\mu \in \overline{C}$.

The radial Dunkl process is defined as the \overline{C} -valued continuous paths Markov process whose generator is given by :

$$\mathcal{L}u(x) = \frac{1}{2}\Delta u(x) + \sum_{\alpha \in R_+} k(\alpha) \frac{\langle \alpha, \nabla u(x) \rangle}{\langle \alpha, x \rangle}$$

with boundary conditions $\nabla u(x) \cdot \alpha = 0$ for all $x \in H_\alpha$, $\alpha \in R_+$, $k(\alpha) \geq 0$ is the multiplicity function (invariant under the action of W), and $u \in C_c^2(C)$. When $k(\alpha) = 1$ for all $\alpha \in R$, we recover the BM constrained to stay in C , studied by Grabiner ([23]). The semi-group density of X is given by :

$$(1) \quad p_t^k(x, y) = \frac{1}{c_k t^{\gamma+m/2}} e^{-(|x|^2+|y|^2)/2t} D_k^W(x, y) \prod_{\alpha \in R_+} \langle \alpha, y \rangle^{2k(\alpha)}$$

for $x, y \in C$, where $\gamma = \sum_{\alpha \in R_+} k(\alpha)$,

$$D_k^W(x, y) := \sum_{w \in W} D_k\left(\frac{x}{\sqrt{t}}, \frac{wy}{\sqrt{t}}\right)$$

where D_k denotes the Dunkl kernel and c_k is given by the Macdonald-Mehta integral ([38]). Indeed, as $D_k(0, y) = 1$ ([38]), one gets

$$t^{\gamma+m/2} c_k = |W| \int_C e^{-|y|^2/2t} \prod_{\alpha \in R_+} \langle \alpha, y \rangle^{2k(\alpha)} dy = \int_{\mathbb{R}^m} e^{-|y|^2/2t} \prod_{\alpha \in R_+} |\langle \alpha, y \rangle|^{2k(\alpha)} dy$$

since $\mathbb{R}^m = \cup_{w \in W} w\overline{C}$. $D_k^W(x, y)$ is known as *the generalized Bessel function* (up to the constant $|W|$). This process is obtained by projecting the Dunkl process valued in \mathbb{R}^m (which has right-continuous and left-limits paths, see [21]) on \overline{C} . The latter was already introduced by Rösler ([38],[39]) and then studied by Gallardo and Yor ([21],[22]) and Chybiryakov ([11]).

To illustrate all these facts and motivate the reader as well, we will provide some well known examples. We start with the *rank one* case ($m = 1$) for which $R = B_1 = \{\pm 1\}$. Hence $k(\alpha) := k_0 \geq 0$ and X is a *Bessel process* ([37]) of *index* $\nu = k_0 - 1/2$. When $k_0 > 0$, it is the unique strong solution of :

$$dX_t = dB_t + \frac{k_0}{X_t} dt, \quad t \geq 0, X_0 = x > 0.$$

where B is a standard BM. Another well known multivariate example is described by the A_{m-1} -type root system defined as :

$$A_{m-1} = \{\pm(e_i - e_j) \mid 1 \leq i < j \leq m\},$$

with positive and simple systems given by :

$$R_+ = \{e_i - e_j, 1 \leq i < j \leq m\} \quad \Delta = \{e_i - e_{i+1}, 1 \leq i \leq m\}$$

where $(e_i)_{1 \leq i \leq m}$ is the standard basis of \mathbb{R}^m . V is the hyperplane of \mathbb{R}^m consisting of vectors that coordinates sum to zero. Without loss of generality, one can take \mathbb{R}^m instead of V so that $C = \{x \in \mathbb{R}^m, x_1 > \dots > x_m\}$. Besides, there is only one orbit and $k(\alpha) := k_1 \geq 0$. Thus, the corresponding radial Dunkl process satisfies :

$$(2) \quad dX_t^i = d\nu_t^i + k_1 \sum_{j \neq i} \frac{dt}{X_t^i - X_t^j} \quad 1 \leq i \leq m, \quad t < \tau$$

with $X_0^1 > \dots > X_0^m$, where $(\nu^i)_i$ are independent Brownian motions and τ is *the first collision time*. For strictly positive k_1 , this process was deeply studied by Cépa and Lépingle ([8], [9], [10]): it behaves as m -interacting particles on the real line with electrostatic repulsions proportional to the inverse of the distance separating them. Moreover, when $k_1 = 1, 1/2$ respectively, this process evolves like eigenvalues process of Hermitian (Dyson model) and symmetric Brownian motions ([20], [23]). It was shown in [8] that this SDE has a unique strong solution for all $t \geq 0$ and $k_1 > 0$. When reading the proof in [8], one hopes to extend this result for any root system since materials used there are not typical for the A_{m-1} -type. This was the original motivation of this work. Our first result claims that

$$dX_t = dB_t - \nabla \Phi(X_t)dt, \quad X_0 \in C$$

where $\Phi(x) = -\sum_{\alpha \in R_+} k(\alpha) \ln(\langle \alpha, x \rangle)$, $k > 0$, has a unique strong solution for all $t \geq 0$. At the same time and independently, Chybiryakov and Schapira provide two other proofs: both authors used well posed martingale problems associated respectively with the \mathbb{R}^m -valued Dunkl and the radial Heckman-Opdam processes as well as geometric arguments ([11], [41]). The curious reader will wonder what happens if $k(\alpha) = 0$ for some α ? The answer is that the same result holds up to *the first hitting time* of ∂C , say T_0 ([11] p. 37). Next, we are mainly interested in the tail distribution of T_0 . Before proceeding, we reprove via stochastic calculus that $T_0 < \infty$ if $k(\alpha) < 1/2$ for at least one $\alpha \in R_+$ (see [11] for the original proof using local martingales). More precisely, for such an α , we prove that almost surely, $\langle \alpha, X_t \rangle \leq Y_t$ for all $t \geq 0$, where Y is a Bessel process of dimension strictly less than 2. At this level, other proofs exist for the above results. To our knowledge, the contents of the remainder of the paper are new. In [11], the author derived absolute-continuity relations which allow us to write the tail distribution of T_0 when starting from $x \in C$. A W -invariant analytic x -dependent integral, which value at 0 is given by a Selberg integral, is involved. As far as we know, though

$D_k^W(x, y)$ arises as hypergeometric functions for particular root systems (see the end of [2]), forward computations are sophisticated and hard. More precisely, we think that it is possible to use the integral formula given in Corollary 2 in [28] with the integration range $0 < X_t^1 < \dots < X_t^m < 1$, known as the Macdonald's conjecture, then perform limit and sums operations. The matrix cases for which the Jack parameter equals to $= 1, 2$ are more handable with the use of properties of zonal polynomials and Schur functions. However, we think that the approach adopted here is more elegant since on one hand, it disregards the special values of the multiplicity function and on the other hand, does not need long hard formulas. It only relies on some properties picked from Dunkl theory. More precisely, it will be shown that the x -dependent integral is an eigenfunction of some operator which involves the generator \mathcal{L} and the so-called *Euler operator* E_1 . For some particular root systems, this eigenfunction is identified with some hypergeometric series. A surprising fact is that the eigenoperator can be expressed in terms of a Schrödinger operator \mathcal{H} and its minimal eigenvalue E_{min} (minimal energy) (see [38] page 18):

$$\mathcal{L} - E_1 := \mathcal{L} - \sum_{i=1}^m x_i \partial_i = -e^{|x|^2/4} (\mathcal{H} - E_{min}) e^{-|x|^2/4}$$

Moreover, $(X_t)_{t \geq 0}$ specializes for some values of k to eigenvalues processes of self-adjoint matrix processes such as symmetric and Hermitian Brownian motions, Wishart and Laguerre and matrix Jacobi processes. In those cases, computations can be performed using the action of orthogonal and unitary groups. Indeed, Jack polynomials fit zonal polynomials and Schur functions when the Jack parameter equals to $1, 2$ respectively (see [34]). The two first ones are identified as A_{m-1} -type radial Dunkl processes while Wishart and Laguerre processes are related to the B_m root system. The latter goes beyond the radial Dunkl setting: the reduced root system C_m in a particular case (ultraspheric) is involved and more generally, the non reduced system BC_m . This connection was deeply investigated in [3] while identifying special functions associated with root systems with multivariate hypergeometric series. Among them appear multivariate Gauss hypergeometric series and Jacobi polynomials ([32]) and these are eigenfunctions of the β -Jacobi generator. The state space is the so-called *principal Weyl alcove* which is now a bounded convex domain and fundamental for the action of the *affine Weyl group*. Hence, the process evolves like particles in an interval. Then, we extend the strong uniqueness Theorem to the Jacobi context. In the remaining part, we derive some properties: we briefly visit the Brownian motion in the principal Weyl alcove which corresponds to multiplicities all equal to 1. Then, an analogous result on the finiteness of the first hitting time of alcoves walls is obtained using similar computations as those for T_0 . Finally, we derive the semi group density and discuss some open questions left in [16].

2. RADIAL DUNKL PROCESS : EXISTENCE AND UNIQUENESS OF A STRONG SOLUTION

Theorem 1. *Let R be a reduced root system. Let:*

$$\Phi(x) = - \sum_{\alpha \in R_+} k(\alpha) \ln(\langle \alpha, x \rangle) := \sum_{\alpha \in R_+} k(\alpha) \theta(\langle \alpha, x \rangle), \quad x \in C$$

where $k(\alpha) > 0$ for all $\alpha \in R_+$. Then the SDE

$$(3) \quad dX_t = dB_t - \nabla \Phi(X_t) dt, \quad X_0 \in C$$

where X is an adapted continuous process valued in \overline{C} and B is a Brownian motion in \mathbb{R}^m , has a unique strong solution.

Proof: From Theorem 2. 2 in [9], we deduce that :

$$(4) \quad dX_t = dB_t - \nabla \Phi(X_t) dt + n(X_t) dL_t, \quad X_0 \in C$$

where $n(x)$ is a (unitary) inward normal vector to C at x , L is the boundary process defined by:

$$dL_t = \mathbf{1}_{\{X_t \in \partial C\}} dL_t,$$

has a unique strong solution for all $t \geq 0$. Moreover :

$$(5) \quad \mathbb{E} \left[\int_0^T \mathbf{1}_{\{X_t \in \partial C\}} dt \right] = 0$$

$$(6) \quad \mathbb{E} \left[\int_0^T |\nabla \Phi(X_t)| dt \right] < \infty$$

for all $T > 0$. Thus, it remains to prove that the boundary process vanishes. To proceed, we need two Lemmas.

Remark. Both Lemmas below discard the reducedness of R . In fact, this assumption figures in the definition of the Dunkl process and originates from analytic purposes like the commutativity of Dunkl operators ([18]).

Lemma 1. *Set $dG_t := n(X_t) dL_t$. Then, $\forall \alpha \in R_+$,*

$$\mathbf{1}_{\{\langle X_t, \alpha \rangle = 0\}} < dG_t, \alpha \rangle = 0$$

Proof: The proof is roughly a generalization of the one in [8] for $R = A_{m-1}$. In order to convince the reader, we provide an outline. Using the occupation density formula, we may write ($\langle \alpha, X \rangle \geq 0$) :

$$\int_0^\infty L_t^\alpha(\langle \alpha, X \rangle) \theta'(a) da = \langle \alpha, \alpha \rangle \int_0^t \theta'(\langle \alpha, X_s \rangle) ds$$

where $L_t^\alpha(\langle \alpha, X \rangle)$ is the local time of the real continuous semimartingale $\langle \alpha, X \rangle$. On the other hand, the following inequality holds (instead of (2.5) in [8])

for all $a \in C$:

$$\begin{aligned}
\langle \nabla \Phi(x), x - a \rangle &= \sum_{\alpha \in R_+} k(\alpha) \theta'(\langle \alpha, x \rangle) \langle \alpha, x - a \rangle \\
&\stackrel{(1)}{\geq} \sum_{\alpha \in R_+} k(\alpha) [b_\alpha \theta'(\langle \alpha, x \rangle) - c_\alpha \langle \alpha, x - a \rangle - d_\alpha] \\
&\geq \min_{\alpha \in R_+} (b_\alpha k(\alpha)) \sum_{\alpha \in R_+} \theta'(\langle \alpha, x \rangle) - |x - a| \sum_{\alpha \in R_+} k(\alpha) c_\alpha |\alpha| - \sum_{\alpha \in R_+} k(\alpha) d_\alpha \\
&:= A \sum_{\alpha \in R_+} \theta'(\langle \alpha, x \rangle) - B|x - a| - C
\end{aligned}$$

by Cauchy-Schwarz inequality, where in (1), we used eq. (2.1) in [8] : let g be a convex C^1 -function on an open convex set $D \subset \mathbb{R}^m$, then $\forall a \in D$, there exist $b, c, d > 0$ such that for all $x \in D$:

$$\langle \nabla g(x), x - a \rangle \geq b|\nabla g(x)| - c|x - a| - d$$

Note also that $A > 0$ since $b_\alpha k(\alpha) > 0$ for all $\alpha \in R_+$. Then, the continuity of X , (6) and the inequality above yield :

$$\int_0^t \theta'(\langle \alpha, X_s \rangle) ds < \infty$$

which implies that :

$$\int_0^\infty L_t^\alpha(\langle \alpha, X \rangle) \theta'(a) da < \infty$$

Thus, $L_t^0(\langle \alpha, X \rangle) = 0$ since the function $a \mapsto \theta'(a)$ is not integrable at 0. The next step consists in using Tanaka formula to compute $dZ_t := \langle \alpha, X_t \rangle - (\langle \alpha, X_t \rangle)^+$ for $\alpha \in \Delta$:

$$dZ_t = \mathbf{1}_{\{\langle \alpha, X_t \rangle = 0\}} \langle \alpha, dB_t \rangle - \mathbf{1}_{\{\langle \alpha, X_t \rangle = 0\}} \langle \alpha, \nabla \Phi(X_t) \rangle dt + \mathbf{1}_{\{\langle \alpha, X_t \rangle = 0\}} \langle \alpha, dG_t \rangle$$

It is obvious that the second term vanishes. The first vanishes too since it is a continuous local martingale with null bracket (occupation density formula). As $X_t \in \overline{C}$, then $dZ_t = 0$ a.s. which gives the result. \blacksquare

Lemma 2. *Let $x \in \partial C$. Then $\langle n(x), \alpha \rangle \neq 0$ for some $\alpha \in \Delta$ such that $\langle x, \alpha \rangle = 0$.*

Proof : Let us assume that $\langle n(x), \alpha \rangle = 0$ for all $\alpha \in \Delta$ such that $\langle x, \alpha \rangle = 0$. Then, our assumption implies that $\langle x, \alpha \rangle > 0$ for all $\alpha \in \Delta$ such that $\langle n(x), \alpha \rangle \neq 0$. If $\langle n(x), \alpha \rangle < 0$ for these simple roots, then $x - n(x) \in \overline{C}$. By the virtue of the definition of the inward normal $n(x)$ to C at x , i. e.,

$$(7) \quad \langle x - a, n(x) \rangle \leq 0, \quad \forall a \in \overline{C},$$

it follows that $n(x)$ is the null vector which is not possible. Otherwise, choosing :

$$0 < \epsilon < \min_{\alpha / \langle n(x), \alpha \rangle > 0} \frac{\langle x, \alpha \rangle}{\langle n(x), \alpha \rangle}$$

we claim that $a := x - \epsilon n(x) \in \partial C$. Arguing as before, we are done. \blacksquare

Now we proceed to end the proof of the Theorem. Note first that $\partial C = \cup_{\alpha \in \Delta} H_\alpha$ so that :

$$\mathbf{1}_{\{X_t \in \partial C\}} dL_t \leq \sum_{\alpha \in \Delta} \mathbf{1}_{\{\langle X_t, \alpha \rangle = 0\}} dL_t.$$

If $X_t \in H_\alpha$ for one and only one $\alpha \in \Delta$. Then, $n(X_t) = \alpha / \|\alpha\|$ and Lemma 1 gives

$$\mathbf{1}_{\{\langle X_t, \alpha \rangle = 0\}} \langle dG_t, \alpha \rangle = \mathbf{1}_{\{\langle X_t, \alpha \rangle = 0\}} \|\alpha\| dL_t = 0$$

Hence, the boundary process vanishes. More generally, we can use the inequality above and write

$$\begin{aligned} 0 \leq L_t &\leq \sum_{\alpha \in \Delta} \int_0^t \mathbf{1}_{\{\langle X_s, \alpha \rangle = 0\}} \mathbf{1}_{\{\langle n(X_s), \alpha \rangle \neq 0\}} dL_s \\ &= \sum_{\alpha \in \Delta} \int_0^t \mathbf{1}_{\{\langle n(X_s), \alpha \rangle \neq 0\}} \frac{1}{\langle n(X_s), \alpha \rangle} \mathbf{1}_{\{\langle X_s, \alpha \rangle = 0\}} \langle dG_s, \alpha \rangle = 0 \end{aligned}$$

by Lemma 1. \blacksquare

Remark. When $m = 1$, $(X_t)_{t \geq 0}$ is a Bessel process of dimension $\delta = 2k_0 + 1$ and $k_0 > 0 \Leftrightarrow \delta > 1$. It is well known that the local time vanishes (see Ch. XI in [37]) which fits our result.

3. FINITENESS OF THE FIRST HITTING TIME OF THE WEYL CHAMBER

Let $T_0 := \inf\{t > 0, X_t \in \partial C\}$ be the first hitting time of the Weyl chamber. It was shown in [11] (see p. 30) that $T_0 = \infty$ almost surely if $k(\alpha) \geq 1/2$ for all $\alpha \in R_+$. In [9], where $R = A_{m-1}$ and $T_0 = \inf\{t > 0, X_t^i = X_t^j \text{ for some } (i, j)\}$, authors showed that $T_0 < \infty$ a.s. if and only if $0 < k_1 < 1/2$. More generally, the following holds (see [11] p. 75 for the original proof) :

Proposition 1. *Let $\alpha_0 \in \Delta$ and $T_{\alpha_0} := \inf\{t > 0, \langle \alpha_0, X_t \rangle = 0\}$ such that $T_0 = \inf_{\alpha_0 \in \Delta} T_{\alpha_0}$. If $0 < k(\alpha_0) < 1/2$, then $(\langle \alpha_0, X_t \rangle)_{t \geq 0}$ hits almost surely 0. In particular, $T_0 < T_{\alpha_0} < \infty$ a. s.*

Proof : assume $k(\alpha) > 0$ for all $\alpha \in R$ and let $\alpha_0 \in \Delta$. Our scheme is roughly the same as that used in [9], thus we shall show that the process $\langle \alpha_0, X \rangle$ is almost surely less than or equal to a Bessel process with dimension $2k(\alpha_0) + 1$. The result follows from the fact that $2k(\alpha_0) + 1 < 2$ when $k(\alpha) < 1/2$. For this,

we use the SDE (3). For all $t \geq 0$,

$$\begin{aligned} d \langle \alpha_0, X_t \rangle &= \|\alpha_0\| d\gamma_t + \sum_{\alpha \in R_+} k(\alpha) \frac{\langle \alpha, \alpha_0 \rangle}{\langle \alpha, X_t \rangle} dt \\ &= \|\alpha_0\| d\gamma_t + k_0 \frac{\|\alpha_0\|^2}{\langle \alpha_0, X_t \rangle} dt + \sum_{\alpha \in R_+ \setminus \alpha_0} k(\alpha) \frac{\langle \alpha, \alpha_0 \rangle}{\langle \alpha, X_t \rangle} dt \end{aligned}$$

where k_0 is the value of $k(\alpha_0)$ corresponding to the conjugacy class of α_0 . Set

$$R = \cup_{j=1}^p R^j$$

where R^j , $1 \leq j \leq p$ denote the conjugacy classes of R under the W -action, then

$$R_+ = \cup_{i=1}^p R_+^j$$

so that:

$$d \langle \alpha_0, X_t \rangle = \|\alpha_0\| d\gamma_t + k_0 \frac{\|\alpha_0\|^2}{\langle \alpha_0, X_t \rangle} dt + \sum_{j=1}^p k_j \sum_{\alpha \in R_+^j \setminus \alpha_0} \frac{\langle \alpha, \alpha_0 \rangle}{\langle \alpha, X_t \rangle} dt$$

For a conjugacy class R^j and $\alpha \in R^j$, if $\langle \alpha, \alpha_0 \rangle = a(\alpha) > 0$ then, it is easy to check that $\langle \sigma_0(\alpha), \alpha_0 \rangle = -a(\alpha)$ where σ_0 is the reflection with respect to the orthogonal hyperplane H_{α_0} defined by :

$$\sigma_0(x) = x - 2 \frac{\langle x, \alpha_0 \rangle}{\langle \alpha_0, \alpha_0 \rangle} \alpha_0$$

Note that $\sigma_0(\alpha)$ belongs to the same conjugacy class of α and that $\sigma_0(\alpha) \in R_+$ for $\alpha \in R_+ \setminus \alpha_0$. Indeed, $\sigma_0(R_+ \setminus \alpha_0) = R_+ \setminus \alpha_0$ for all $\alpha_0 \in \Delta$ (see Proposition 1. 4 in [27]). Hence,

$$d \langle \alpha_0, X_t \rangle = \|\alpha_0\| d\gamma_t + k_0 \frac{\|\alpha_0\|^2}{\langle \alpha_0, X_t \rangle} dt - \sum_{j=1}^p k_j \sum_{\substack{\alpha \in R_+^j \setminus \alpha_0 \\ a(\alpha) > 0}} \frac{a(\alpha) \langle \alpha - \sigma_0(\alpha), X_t \rangle}{\langle \alpha, X_t \rangle \langle \sigma_0(\alpha), X_t \rangle} dt$$

Furthermore,

$$\alpha - \sigma_0(\alpha) = 2 \frac{\langle \alpha, \alpha_0 \rangle}{\langle \alpha_0, \alpha_0 \rangle} \alpha_0 \quad \Rightarrow \quad \langle \alpha - \sigma_0(\alpha), X_t \rangle = 2a(\alpha) \frac{\langle \alpha_0, X_t \rangle}{\|\alpha_0\|^2}$$

Consequently, one gets :

$$d \langle \alpha_0, X_t \rangle = \|\alpha_0\| d\gamma_t + k_0 \frac{\|\alpha_0\|^2}{\langle \alpha_0, X_t \rangle} dt + F_t dt$$

where $F_t < 0$ on $\{T_{\alpha_0} = \infty\}$. Using the comparison Theorem in [29] (Proposition 2. 18. p. 293 and Exercice 2. 19. p. 294), one claims that $\langle \alpha_0, X_t \rangle \leq Y_{\|\alpha_0\|^2 t}^x$ for all $t \geq 0$ on $\{T_{\alpha_0} = \infty\}$, where Y^x is a Bessel process defined on the same probability space with respect to the same Brownian motion, of dimension $2k_0 + 1$

and starting at $Y_0 = x \geq \alpha_0, X_0 >> 0$. This is not possible since a Bessel process of dimension < 2 hits 0 a. s. ([37] Chap. XI) \blacksquare

Remark. If we remove the assumption $k(\alpha) > 0$ for all $\alpha \in R$, then the SDE (3) can be solved up to time T_0 when starting from $x \in C$ (see [11]).

4. THE LAW OF T_0

Here, we focus on the tail distribution of T_0 deduced from absolute continuity relations derived in ([11]). Recall that (see [11]) the *index* of X is defined by $l(\alpha) := k(\alpha) - 1/2$. The last result asserts that if $-1/2 < l(\alpha) < 0$ for some $\alpha \in \Delta$, then $T_0 < \infty$ almost surely. Besides, if $l(\alpha) \geq 0$ for all $\alpha \in \Delta$ then $T_0 = \infty$ almost surely. Taking into account these statements, two major parts are considered: $l(\alpha) \geq 0$ for all $\alpha \in R$ so that the process with index $-l$ hits 0 almost surely, and $l(\alpha) < 0$ for at least one α . The tail distribution involves a W -invariant x -dependent integral. Our line of thinking relies on showing that it is an eigenfunction of an appropriate differential operator. Then, using uniqueness results for some differential equations, the tail distribution is written in A_{m-1} and B_m cases by means of multivariate hypergeometric functions. In the last case, we recover known results from matrix theory for Wishart and Laguerre processes. However, we find it better to postpone this in the next section where links with eigenvalues of matrix-valued processes are detailed.

4.1. A first formula. Let us denote by P_x^l the law of $(X_t)_{t \geq 0}$ starting from $x \in C$. Let E_x^l be the corresponding expectation. Recall that ([11], Proposition 2.15.c), if $l(\alpha) \geq 0 \forall \alpha \in R_+$, then:

$$\begin{aligned} P_x^{-l}(T_0 > t) &= E_x^l \left[\left(\prod_{\alpha \in R_+} \frac{\langle \alpha, X_t \rangle}{\langle \alpha, x \rangle} \right)^{-2l(\alpha)} \right] \\ &= \prod_{\alpha \in R_+} \langle \alpha, x \rangle^{2l(\alpha)} \frac{e^{-|x|^2/2t}}{c_k t^{\gamma+m/2}} \int_C e^{-|y|^2/2t} D_k^W \left(\frac{x}{\sqrt{t}}, \frac{y}{\sqrt{t}} \right) \prod_{\alpha \in R_+} \langle \alpha, y \rangle dy \\ &= \prod_{\alpha \in R_+} \langle \alpha, x \rangle^{2l(\alpha)} \frac{e^{-|x|^2/2t}}{c_k t^{\gamma'}} \int_C e^{-|y|^2/2t} D_k^W \left(\frac{x}{\sqrt{t}}, y \right) \prod_{\alpha \in R_+} \langle \alpha, y \rangle dy \\ &:= \prod_{\alpha \in R_+} \langle \alpha, x \rangle^{2l(\alpha)} \frac{e^{-|x|^2/2t}}{c_k t^{\gamma'}} g \left(\frac{x}{\sqrt{t}} \right) \end{aligned}$$

where $\gamma = \sum_{\alpha \in R_+} k(\alpha)$ and $\gamma' = \gamma - |R_+|/2$.

Though D_k^W is given by hypergeometric functions in the special cases A_{m-1} and B_m (see the end of [2]), the *Jack polynomials* defining them prevent us from making computations. However, this may be possible when these polynomials fit, for some values of k , *Zonal polynomials* and *Schur functions* (see [34] for definitions). Our

main result does not make these restrictions and uses some properties of the Dunkl kernel D_k :

Theorem 2. *Let T_i be the i -th difference Dunkl operator and $\Delta_k = \sum_{i=1}^m T_i^2$ the Dunkl Laplacian ([38]). Define :*

$$\mathcal{J}_k^x := -\Delta_k^x + \sum_{i=1}^m x_i \partial_i^x := -\Delta_k^x + E_1^x$$

where $E_1^x := \sum_{i=1}^m x_i \partial_i^x$ is the Euler operator and the superscript indicates the derivative action. Then

$$\mathcal{J}_k^x \left[e^{-|y|^2/2} D_k^W(x, y) \right] = E_1^y \left[e^{-|y|^2/2} D_k^W(x, y) \right]$$

Proof: Recall that if f is W -invariant then $T_i^x f = \partial_i^x f$ and that $T_i^x D_k(x, y) = y_i D_k(x, y)$ (see [38]). Then, on one hand :

$$\begin{aligned} \Delta_k^x D_k^W(x, y) &= \sum_{w \in W} \sum_{i=1}^m (wy)_i T_i^x D_k(x, wy) = \sum_{w \in W} \sum_{i=1}^m (wy)_i^2 D_k(x, wy) \\ &= \sum_{i=1}^m y_i^2 \sum_{w \in W} D_k(x, wy) := p_2(y) D_k^W(x, y) \end{aligned}$$

On the other hand :

$$\begin{aligned} E_1^x D_k^W(x, y) &= \sum_{w \in W} \sum_{i=1}^m x_i T_i^x D_k(x, wy) = \sum_{w \in W} \sum_{i=1}^m (x_i) (wy)_i D_k(x, wy) \\ &= \sum_{w \in W} \langle x, wy \rangle D_k(x, wy) = \sum_{w \in W} \langle w^{-1}x, y \rangle D_k(x, wy) \\ &= E_1^y D_k^W(x, y) \end{aligned}$$

where the last equality follows from $D_k(x, wy) = D_k(w^{-1}x, y)$ since $D_k(wx, wy) = D_k(x, y)$ for all $w \in W$. The result follows from an easy computation.

Corollary 1. *g is an eigenfunction of $-\mathcal{J}_k$ corresponding to the eigenvalue $m + |R_+|$.*

Proof: Theorem 1 and an integration by parts give :

$$\begin{aligned}
-\mathcal{J}_k^x g(x) &= -\int_C E_1^y \left[e^{-|y|^2/2} D_k^W(x, y) \right] \prod_{\alpha \in R_+} \langle \alpha, y \rangle dy \\
&= -\sum_{i=1}^m \int_C y_i \prod_{\alpha \in R_+} \langle \alpha, y \rangle \partial_i^y \left[e^{-|y|^2/2} D_k^W(x, y) \right] dy \\
&= \sum_{i=1}^m \int_C e^{-|y|^2/2} D_k^W(x, y) \partial_i \left[y_i \prod_{\alpha \in R_+} \langle \alpha, y \rangle \right] dy \\
&= \int_C e^{-|y|^2/2} D_k^W(x, y) \prod_{\alpha \in R_+} \langle \alpha, y \rangle \sum_{i=1}^m \left[1 + \sum_{\alpha \in R_+} \frac{\alpha_i y_i}{\langle \alpha, y \rangle} \right] dy
\end{aligned}$$

and the proof ends by summing over i . ■

- The A_{m-1} case : as mentioned in the introduction, the A_{m-1} -type root system is characterized by :

$$\begin{aligned}
R &= \{\pm(e_i - e_j), 1 \leq i < j \leq m\} & R_+ &= \{e_i - e_j, 1 \leq i < j \leq m\} \\
\Delta &= \{e_i - e_{i+1}, 1 \leq i \leq m\} & C &= \{y \in \mathbb{R}^m, y_1 > y_2 > \dots > y_m\}
\end{aligned}$$

$W = S_m$ is the permutations group and there is one conjugacy class so that $k = k_1 > 0 \Rightarrow \gamma = k_1 m(m-1)/2$. Moreover, the generalized Bessel function ² is given by ([2] p. 212-214, [12]):

$$\frac{1}{|W|} D_k^W(x, y) = {}_0F_0^{(1/k_1)}(x, y) := \sum_{p=0}^{\infty} \sum_{\tau} \frac{J_{\tau}^{(1/k_1)}(x) J_{\tau}^{(1/k_1)}(y)}{J_{\tau}^{(1/k_1)}(1) p!}$$

where $\tau = (\tau_1, \dots, \tau_m)$ is a partition of weight $|\tau| = p$ and length m , $J_{\tau}^{(1/k_1)}$ is the Jack polynomial of Jack parameter $1/k_1$ ³, (see [2], [34]). Hence, letting V to be the Vandermonde function :

$$P_x^{-l}(T_0 > t) = V(x)^{2k_0-1} \frac{|W| e^{-|x|^2/2t}}{c_k t^{k_0 m(m-1)/2}} \int_C e^{-|y|^2/2} {}_0F_0^{(1/k_0)}\left(\frac{x}{\sqrt{t}}, y\right) V(y) dy$$

Besides, \mathcal{J}_k writes on W -invariant functions

$$-\mathcal{J}_k^x = D_0^x - E_1^x := \sum_{i=1}^m \partial_i^{2,x} + 2k_1 \sum_{i \neq j} \frac{1}{x_i - x_j} \partial_i - \sum_{i=1}^m x_i \partial_i^x$$

²Authors use the change of variable $x \mapsto \sqrt{2}x$, $y \mapsto \sqrt{2}y$ to fit the hypergeometric function obtained when deriving the generating function for Hermite polynomials. This in turn will modify the eigenoperator by a multiplying constant (see p. 183).

³With the same notations in [2], $k_1 = 2/\alpha$. This can be seen either from the eigenoperator below or from the orthogonality weight function involved in the semi group density.

Finally, since g is W -invariant, then

$$(D_0^x - E_1^x)g(x) = m \frac{m+1}{2}g(x),$$

$$g(0) = \int_C e^{-|y|^2/2} V(y) dy = \frac{1}{m!} \int_{\mathbb{R}^m} e^{-|y|^2/2} |V(y)| dy$$

Let us recall that the Gauss hypergeometric function

$${}_2F_1^{(1/k_1)}(e, b, c, z) = \sum_{p=0}^{\infty} \sum_{\tau} \frac{(e)_{\tau} (b)_{\tau} J_{\tau}^{(1/k_1)}(z)}{(c)_{\tau} p!}$$

is the unique symmetric eigenfunction that equals to 1 at 0 of (see [3] p. 585)

$$(8) \quad \sum_{i=1}^m z_i(1-z_i) \partial_i^{2,z} + 2k_1 \sum_{i \neq j} \frac{z_i(1-z_i)}{z_i - z_j} \partial_i^z + \sum_{i=1}^m [c - k_1(m-1) - (e+b+1 - k_1(m-1)) z_i] \partial_i^z$$

associated to the eigenvalue meb . Letting $z = (1/2)(1 - x/\sqrt{b})$, $e = (m+1)/2$ and

$$c = k_1(m-1) + \frac{1}{2}[e+b+1 - k_1(m-1)] = \frac{b}{2} + \frac{k_1}{2}(m-1) + \frac{m+3}{4}$$

the resulting function is an eigenfunction of

$$\sum_{i=1}^m (1 - \frac{x_i^2}{b}) \partial_i^{2,x} + 2k_1 \sum_{i \neq j} \frac{(1 - x_i^2/b)}{x_i - x_j} \partial_i^x - \sum_{i=1}^m (b + \frac{m+3}{2} - k_1(m-1)) \frac{x_i}{b} \partial_i^x$$

and $D_0^x - E_1^x$ is the limiting operator as b tends to infinity. Hence,

Proposition 2. For $k_1 \geq 1/2$,

$$g(x) = g(0)C(m, k_1) \lim_{b \rightarrow \infty} {}_2F_1^{(1/k_1)} \left[m+1, b, \frac{b}{2} + \frac{k_1}{2}(m-1) + \frac{m+3}{2}, \frac{1}{2} \left(1 - \frac{x}{\sqrt{b}} \right) \right]$$

where

$$C(m, k_1)^{-1} = \lim_{b \rightarrow \infty} {}_2F_1^{(1/k_1)} \left(m+1, b, \frac{b}{2} + \frac{k_1}{2}(m-1) + \frac{m+3}{2}, \frac{1}{2} \right)$$

- The B_m case : This root system is defined by

$$R = \{\pm e_i, \pm e_i \pm e_j, 1 \leq i < j \leq m\} \quad R_+ = \{e_i, 1 \leq i \leq m, e_i \pm e_j, 1 \leq i < j \leq m\}$$

$$\Delta = \{e_i - e_{i+1}, 1 \leq i \leq m, e_m\} \quad C = \{y \in \mathbb{R}^m, y_1 > y_2 > \dots > y_m > 0\}$$

The Weyl group is generated by transpositions and "change sign" reflections ($x_i \mapsto -x_i$) and there are two conjugacy classes so that $k = (k_0, k_1) \Rightarrow$

$\gamma = mk_0 + m(m-1)k_1$. The generalized Bessel function⁴ is given by ([2] p. 214) :

$$\frac{1}{|W|} D_k^W(x, y) = {}_0F_1^{(1/k_1)}(k_0 + (m-1)k_1 + \frac{1}{2}, \frac{x^2}{2t}, \frac{y^2}{2t})$$

where

$${}_0F_1^{(1/k_1)}(c, x, y) = \sum_{p=0}^{\infty} \sum_{\tau} (c)_{\tau} \frac{J_{\tau}^{(k_1)}(x) J_{\tau}^{(1/k_1)}(y)}{J_{\tau}^{(1/k_1)}(1) p!}$$

and $(c)_{\tau} := \prod_{i=1}^m (c - k_1(i-1))_{\tau_i}$ is the generalized Pochhammer symbol (see [2]). Then, one has :

$$g(x) = |W| \int_C e^{-|y|^2/2} {}_0F_1^{(1/k_1)}(k_0 + (m-1)k_1 + \frac{1}{2}, \frac{x^2}{2}, \frac{y^2}{2}) \prod_{i=1}^m (y_i) V(y^2) dy$$

The eigenoperator writes on W -invariant functions:

$$\begin{aligned} -\mathcal{J}_k^x &= \sum_{i=1}^m \partial_i^{2,x} + 2k_0 \sum_{i=1}^m \frac{1}{x_i} \partial_i^x + 2k_1 \sum_{i \neq j} \left[\frac{1}{x_i - x_j} + \frac{1}{x_i + x_j} \right] \partial_i^x - E_1^x \\ -\mathcal{J}_k^x g(x) &= m(m+1)g(x), \quad g(0) = \frac{1}{2^m m!} \int_{\mathbb{R}^m} e^{-|y|^2} \prod_{i=1}^m |y_i| |V(y^2)| dy. \end{aligned}$$

A change of variable $x_i = \sqrt{2y_i}$ shows that $u(y) := g(\sqrt{2y})$ satisfies

$$\begin{aligned} -\tilde{\mathcal{J}}_k^y u(y) &= m \frac{(m+1)}{2} u(y), \quad g(0) = u(0) \\ -\tilde{\mathcal{J}}_k^y &= \sum_{i=1}^m y_i \partial_i^{2,y} + 2k_1 \sum_{i \neq j} \frac{y_i}{y_i - y_j} \partial_i^y + \left(k_0 + \frac{1}{2} \right) \sum_{i=1}^m \partial_i^y - E_1^y. \end{aligned}$$

which implies that :

$$u(y) = u(0) {}_1F_1^{(1/k_1)}\left(\frac{m+1}{2}, k_0 + (m-1)k_1 + \frac{1}{2}, y\right)$$

where

$${}_1F_1^{(1/k_1)}(b, c, z) = \sum_{p=0}^{\infty} \sum_{\tau} \frac{(b)_{\tau}}{(c)_{\tau}} \frac{J_{\tau}^{(1/k_1)}(z)}{p!}$$

This can be seen from the differential equation (8) and using ([2]):

$$\lim_{e \rightarrow \infty} {}_2F_1^{(1/k_1)}(e, b, c, \frac{z}{e}) = {}_1F_1^{(1/k_1)}(b, c, z)$$

⁴there is an erroneous sign in one of the arguments in [2]. Moreover, to recover this expression in the B_m case from that given in [2], one should make substitutions $a = k_0 - 1/2$, $k_1 = 1/\alpha$, $q = 1 + (m-1)k_1$. We point to the reader that this is different from the one used in [12] p. 121.

Finally

$$g\left(\frac{x}{\sqrt{t}}\right) = g(0) {}_1F_1^{(1/k_1)}\left(\frac{m+1}{2}, k_0 + (m-1)k_1 + \frac{1}{2}, \frac{x^2}{2t}\right)$$

Hence, the tail distribution is given by :

Proposition 3. For $k_0, k_1 \geq 1/2$,

$$P_x^{-l}(T_0 > t) = C_k \prod_{i=1}^m \left(\frac{x_i^2}{2t}\right)^{k_0-1/2} \left(V\left(\frac{x^2}{2t}\right)\right)^{2k_1-1} e^{-|x|^2/2t} {}_1F_1^{(1/k_1)}\left(\frac{m+1}{2}, k_0+(m-1)k_1+\frac{1}{2}, \frac{x^2}{2t}\right)$$

Remark. 1/Adopting the notations used in [2], one has :

$$-\tilde{\mathcal{J}}_k^y = D_1^y + (a+1)E_0^y - E_1^y \quad (R = B_m, y = x^2),$$

Besides, Theorem 2 was derived there differently for both A_{m-1} and B_m cases when proving a generating function Theorem for generalized Hermite and Laguerre polynomials (page 183 and 192, see also [12]).

4.2. A second formula. In [11] (see Proposition 2.15.b), the author derived another absolute-continuity relation from which we deduce that if $l(\alpha) < 0$ for at least one $\alpha \in R_+$, then

$$\begin{aligned} P_x^l(T_0 > t) &= E_x^0 \left[\prod_{\alpha \in R_+} \left(\frac{\langle \alpha, X_t \rangle}{\langle \alpha, x \rangle}\right)^{l(\alpha)} \exp\left(-\frac{1}{2} \sum_{\alpha, \gamma \in R_+} \int_0^t \frac{\langle \alpha, \gamma \rangle l(\alpha)l(\gamma)}{\langle \alpha, X_s \rangle \langle \gamma, X_s \rangle} ds\right) \right] \\ &= E_x^r \left[\prod_{\alpha \in R_+} \left(\frac{\langle \alpha, X_t \rangle}{\langle \alpha, x \rangle}\right)^{l(\alpha)-r(\alpha)} \exp\left(-\frac{1}{2} \sum_{\alpha, \gamma \in R_+} \int_0^t \frac{\langle \alpha, \gamma \rangle l(\alpha, \gamma)}{\langle \alpha, X_s \rangle \langle \gamma, X_s \rangle} ds\right) \right] \end{aligned}$$

where the last equality follows from part (c) of the same Proposition, $l(\alpha, \gamma) = l(\alpha)l(\gamma) - r(\alpha)r(\gamma)$ and

$$r(\alpha) = \begin{cases} l(\alpha) & \text{if } l(\alpha) \geq 0 \\ -l(\alpha) & \text{if } l(\alpha) < 0 \end{cases}$$

Then $l(\alpha, \gamma) = 0$ if $l(\alpha)l(\gamma) \geq 0$ and $l(\alpha, \gamma) = -2r(\alpha)r(\gamma)$ else. As a result,

$$P_x^l(T_0 > t) = E_x^r \left[\prod_{\substack{\alpha \in R_+ \\ l(\alpha) < 0}} \left(\frac{\langle \alpha, x \rangle}{\langle \alpha, X_t \rangle}\right)^{2r(\alpha)} \exp\left(\sum_{\substack{\alpha, \gamma \in R_+ \\ l(\alpha)l(\gamma) < 0}} \int_0^t \frac{\langle \alpha, \gamma \rangle r(\alpha)r(\gamma)}{\langle \alpha, X_s \rangle \langle \gamma, X_s \rangle} ds\right) \right]$$

Next, note that the exponential functional equals 1 for both root systems A_{m-1} and B_m . For the first, it is obvious since R consists of one orbit so that $\{\alpha, \gamma \in R_+, l(\alpha)l(\gamma) < 0\}$ is empty. This gives the same expression already considered in

the previous subsection. For the second, writing $R_+ = \{e_i, 1 \leq i \leq m\} \cup \{e_j \pm e_k, 1 \leq j < k \leq m\}$ so that $\langle e_i, e_j \pm e_k \rangle = \delta_{ij} \pm \delta_{ik}$ gives

$$\begin{aligned} S &= \sum_{i=1}^m \sum_{i < k} \frac{1}{X_t^i} \left[\frac{1}{X_t^i - X_t^k} + \frac{1}{X_t^i + X_t^k} \right] + \sum_{i=1}^m \sum_{k < i} \frac{1}{X_t^i} \left[\frac{-1}{X_t^k - X_t^i} + \frac{1}{X_t^k + X_t^i} \right] \\ &= \sum_{i=1}^m \sum_{i < k} \frac{2}{(X_t^i)^2 - (X_t^k)^2} - \sum_{i=1}^m \sum_{k < i} \frac{2}{(X_t^k)^2 - (X_t^i)^2} = 0 \end{aligned}$$

where S stands for the sum between brackets. The reader can also check that this holds for C_m and D_m root systems (see the end of the paper for definitions). However, we restrict ourselves to the B_m -case since, for particular values of the multiplicity function, we will recover a known result from matrix theory (see next section). Let us investigate the case $k_0 < 1/2$, $k_1 \geq 1/2$ for which $l_0 < 0$, $l_1 \geq 0$. One writes :

$$g(x) = \int_C e^{-|y|^2/2} D_k^W(x, y) \prod_{i=1}^m (y_i) V^{2k_1}(y^2) dy$$

The machinery used before still applies and gives :

$$- \mathcal{J}_k g = 2m[1 + k_1(m-1)]g$$

Thus

Proposition 4. *In the B_m case and for $k_0 < 1/2$, $k_1 \geq 1/2$, one has:*

$$P_x^l(T_0 > t) = C_k \prod_{i=1}^m \left(\frac{x_i^2}{2t} \right)^{k_0 - 1/2} e^{-|x|^2/2t} {}_1F_1^{(1/k_1)}(1 + k_1(m-1), k_0 + (m-1)k_1 + \frac{1}{2}, \frac{x^2}{2t})$$

In the remaining case $k_0 \geq 1/2$, $k_1 < 1/2$, the tail distribution writes :

$$g(x) = \int_C e^{-|y|^2/2} D_k^W(x, y) \prod_{i=1}^m (y_i)^{2k_0} V(y^2) dy$$

Thus

$$- \mathcal{J}_k g = m[2k_0 + m]g$$

so that

Proposition 5. *For $k_0 \geq 1/2$, $k_1 < 1/2$,*

$$P_x^l(T_0 > t) = C_k V \left(\frac{x^2}{2t} \right)^{2k_1 - 1} e^{-|x|^2/2t} {}_1F_1^{(1/k_1)}(k_0 + \frac{m}{2}, k_0 + (m-1)k_1 + \frac{1}{2}, \frac{x^2}{2t})$$

5. β -PROCESSES AND RANDOM MATRICES

In the sequel, we will see how eigenvalues of some classical matrix-valued processes and radial Dunkl processes are interrelated using SDE. This connection was already checked by physicists throughout eigenvalues probability densities and Fokker-Planck equations for parameter-dependent random matrices ([12]). As we mentioned in the introduction, the A_{m-1} -type is connected to symmetric and Hermitian Brownian motions. Set $k := \beta/2$, $\beta > 0$, then such a process will be called β -Dyson, referring to the Dyson model when $\beta = 2$. This parameter is called the Dyson index. Henceforth, we will adopt new notation for the eigenvalues process, we will write λ instead of X .

5.1. The B_m -type: β -Laguerre processes. The B_m system turns out to be related to eigenvalues of Wishart and Laguerre processes which satisfy the following stochastic differential system (see [5],[13]):

$$d\lambda_i(t) = 2\sqrt{\lambda_i(t)} d\nu_i(t) + \beta \left[\delta + \sum_{k \neq i} \frac{\lambda_i(t) + \lambda_k(t)}{\lambda_i(t) - \lambda_k(t)} \right] dt \quad 1 \leq i \leq m.$$

for $\beta = 1, 2$ and $\delta \geq m + 1, m$ respectively, where $(\nu_i)_i$ are independent Brownian motions and $\lambda_1(0) > \dots > \lambda_m(0)$. Recall that the process remains strictly positive if it is initially strictly positive. This suggests to define the β -Laguerre process as the solution, when it exists, of :

$$d\lambda_i(t) = 2\sqrt{\lambda_i(t)} d\nu_i(t) + \beta \left[\delta + \sum_{k \neq i} \frac{\lambda_i(t) + \lambda_k(t)}{\lambda_i(t) - \lambda_k(t)} \right] dt \quad 1 \leq i \leq m, \quad t < \tau \wedge R_0,$$

where $R_0 = \inf\{t, \lambda_m(t) = 0\}$, $\beta, \delta > 0$ and with $\lambda_1(0) > \dots > \lambda_m(0) > 0$. It is very easy to see that :

$$dD_t := d \left(\prod_{i=1}^m \lambda_i(t) \right) = 2D_t \sqrt{\sum_{i=1}^m \frac{1}{\lambda_i(t)}} d\tilde{B}_t + \beta(\delta - m + 1)D_t \sum_{i=1}^m \frac{1}{\lambda_i(t)} dt$$

for all $t < \tau \wedge R_0$ where \tilde{B} is a standard Brownian motion. It follows that $\forall r \in \mathbb{R}$:

$$\begin{aligned} d(\ln(D_t)) &= 2 \sqrt{\sum_{i=1}^m \frac{1}{\lambda_i(t)}} d\tilde{B}_t + [\beta(\delta - m) + \beta - 2] \sum_{i=1}^m \frac{1}{\lambda_i(t)} dt \\ d(\det(D_t)^r) &= M_t + r[\beta(\delta - m + 1) + 2r - 2] D_t^r \sum_{i=1}^m \frac{1}{\lambda_i(t)} dt \end{aligned}$$

where $M_t = 2r D_t^r \sqrt{\sum_{i=1}^m 1/\lambda_i(t)} d\tilde{B}_t$. From these two SDE, we can argue as in the Wishart and Laguerre cases that $R_0 > \tau$ a.s. when $\delta \geq m - 1 + 2/\beta$ (choose

$2r = 2 - \beta(\delta - m + 1) < 0$ when $\delta > m - 1 + 2/\beta$ then use McKean's argument). Set $r_i := \sqrt{\lambda_i}$, then, for $t < \tau \wedge R_0$:

$$\begin{aligned} dr_i(t) &= d\nu_i(t) + \frac{1}{2r_i(t)} \left[\beta\delta - 1 + \beta \sum_{j \neq i} \frac{r_i^2 + r_j^2}{r_i^2 - r_j^2} \right] dt \\ &= d\nu_i(t) + \frac{\beta(\delta - m + 1) - 1}{2r_i(t)} dt + \frac{\beta}{2} \sum_{j \neq i} \left[\frac{1}{r_i(t) - r_j(t)} + \frac{1}{r_i(t) + r_j(t)} \right] dt \\ &= d\nu_i(t) + \frac{k_0}{r_i(t)} dt + k_1 \sum_{j \neq i} \left[\frac{1}{r_i(t) - r_j(t)} + \frac{1}{r_i(t) + r_j(t)} \right] dt \end{aligned}$$

with $2k_0 = \beta(\delta - m + 1) - 1$, $2k_1 = \beta$. Consequently, the process $r = (r_1, \dots, r_m)$ defined for all $t < \tau \wedge R_0$ is a B_m -radial Dunkl process. Using Theorem 1, one claim that the SDE above has a unique strong solution for all $t \geq 0$ and all β, δ such that $k_0, k_1 > 0$. This strengthen results from matrix theory : in the Wishart setting ($\beta = 1$), the strong uniqueness holds for $\delta > m$ and in the Laguerre case ($\beta = 2$), it holds for $\delta > m - 1/2$. Besides, the generalized Bessel function is given by ([2])⁵ :

$$\frac{1}{|W|} \sum_{w \in W} D_k(x, wy) = {}_0F_1^{(2/\beta)} \left(\frac{\beta\delta}{2}, \frac{x^2}{2}, \frac{y^2}{2} \right) := \sum_{p=0}^{\infty} \sum_{\tau} \binom{\beta\delta}{2}_{\tau} \frac{J_{\tau}^{(2/\beta)}(x^2/2) J_{\tau}^{(2/\beta)}(y^2/2)}{J_{\tau}^{(2/\beta)}(1_m) p!}$$

so that (1) writes

$$(9) \quad p_t^{k_0, k_1}(x, y) = \frac{|W|}{c_k t^{\gamma+m/2}} e^{-(|x|^2+|y|^2)/2t} {}_0F_1^{(2/\beta)} \left(\frac{\beta\delta}{2}, \frac{x^2}{2t}, \frac{y^2}{2t} \right) \prod_{i=1}^m (y_i)^{2k_0} V^{2k_1}(y^2) dy$$

where V stands for the Vandermonde function. Using the variable change $y \mapsto \sqrt{y}$, the semi-group density of the β -Laguerre process writes :

$$q_t^{k_0, k_1}(x, y) = \frac{C_{k_0, k_1}}{t^{\gamma+2k_1+m/2}} e^{-(\sum_{i=1}^m (x_i+y_i)/2t)} {}_0F_1^{(2/\beta)} \left(\frac{\beta\delta}{2}, \frac{x}{2t}, \frac{y}{2t} \right) \prod_{i=1}^m (y_i)^{k_0-1/2} V^{2k_1}(y) dy$$

For $x = 0$ and $t = 1$, we recover the same p.d.f. given in [17] for β -Laguerre ensemble.

Remarks. 1/ Recall that for all $\alpha \in R$, we set $l(\alpha) = k(\alpha) - 1/2$. Hence, in the B_m -case, $l_0 = k_0 - 1/2$, $l_1 = k_1 - 1/2$. For $-l$, all corresponding parameters will be primed. For instance, $-l_0 = k'_0 - 1/2$, $-l_1 = k'_1 - 1/2$. Let us consider a Wishart process of dimension δ' such that $m - 1 \leq \delta' < m + 1$ ([5]), $k'_1 = 1/2$ ($\beta' = 1$) and $k'_0 = (\delta' - m)/2 \Rightarrow -l_1 = 0$, $-l_0 = (\delta' - m - 1)/2 < 0$. Set $\delta' = m + 1 - 2\nu$, $0 < \nu <$

⁵With the same notations used in [2], one has $\beta a'/2 = k_0$, $a = k_0 - 1/2$, $\beta = 2/\alpha \Rightarrow a + q = \beta\delta/2$.

1/2, then, $l_1 = 0, l_1 = \nu \Rightarrow k_1 = 1/2 (\beta = 1)$ and $k_0 = \nu + 1/2 (\delta = m + 1 + 2\nu)$. Results of 4.1 writes:

$$P_x^{-l}(T_0 > t) = C_k \prod_{i=1}^m \left(\frac{x_i^2}{2t} \right)^\nu e^{-|x|^2/2t} {}_1F_1^{(2)}\left(\frac{m+1}{2}, \frac{\delta}{2}, \frac{x^2}{2t}\right)$$

which fits the expression already derived in [15]. When $k'_0 = k'_1 = 0$ ($-l_0 = -l_1 = -1/2$), then $k_0 = k_1 = 1 (\beta = 2, \delta = m + 1/2)$ and the Jack polynomials fits the Schur functions (see [34]). In that case, the following representation holds ([24])

$${}_1F_1^{(1)}(a, b, z) = \frac{\det(z_i^{m-j} {}_1\mathcal{F}_1(a-j+1, b-j+1, z_i)_{1 \leq i, j \leq m})}{V(z)}$$

where ${}_1\mathcal{F}_1$ denotes the univariate hypergeometric function. Hence, the tail distribution writes :

$$P_x^{-l}(T_0 > t) = C_k \det \left[\left(\frac{x_i^2}{2t} \right)^{m-j+1/2} e^{-x_i^2/2t} {}_1\mathcal{F}_1\left(\frac{m+1}{2} - j + 1, m - j + \frac{3}{2}, \frac{x_i^2}{2t}\right) \right]_{1 \leq i, j \leq m}$$

The corresponding process is the Brownian motion in the Weyl chamber of B -type. For Laguerre processes ([13]) of dimension δ such that $m-1/2 \leq \delta < m$, one should apply results derived in section 4.2. Take $k'_1 = 1 (\beta' = 2)$ and $k'_0 = \delta' - m + 1/2 \Rightarrow l_1 = 1/2, l_0 = \delta' - m := -\nu$ with $0 < \nu < 1/2$. Thus $r_1 = 1/2, r_0 = \nu \Rightarrow k_1 = 1 (\beta = 2)$ and $k_0 = \nu + 1/2 (\delta = m + \nu)$ so that :

$$P_x^l(T_0 > t) = C_k \prod_{i=1}^m \left(\frac{x_i^2}{2t} \right)^\nu e^{-|x|^2/2t} {}_1F_1^{(1)}(m, \delta, \frac{x^2}{2t})$$

2/Recall that when $\beta = 2, {}_0F_1^{(1)}$ has a determinantal representation (see [24]) yielding to König and O'Connell result on the V -transform of m -independent squared Bessel processes (BESQs) constrained never to collide (or stay in the A_{m-1} -type Weyl chamber, see [31]). Similar results holds for A_{m-1} -type root system with Brownian motions instead of BESQs. Nevertheless, when $k_0 = k_1 = 1 (\beta = 2, \delta = m + 1/2)$, a similar interpretation involving m -independent Brownian motions killed when they reaches 0 holds. However, the Vandermonde function may be replaced by the product over positive roots. In this case, the eigenvalues process is known as the BM in the Weyl chamber of type B_m (see [23] for further details and other root systems). Since $\gamma = m^2$ and from ([24], [13]) :

$${}_rF_s^{(1)}((m+a_i)_{1 \leq i \leq r}, (m+b_j)_{1 \leq j \leq s}, x, y) = \frac{\det[{}_r\mathcal{F}_s((a_i+1)_{1 \leq i \leq r}, (b_j+1)_{1 \leq j \leq s}, x_i y_j)]_{l, f}}{V(x)V(y)}$$

(9) transforms to :

$$\begin{aligned} p_t^{1,1}(x, y) &= C_m \frac{h(y)}{h(x)} \frac{e^{-(|x|^2+|y|^2)/2t}}{t^{m/2}} \prod_{i,j=1}^m \left(\frac{x_i y_j}{t} \right) \det \left[{}_0\mathcal{F}_1 \left(\frac{1}{2} + 1, \frac{(x_i y_j)^2}{4t^2} \right) \right]_{i,j} \\ &= C_m \frac{h(y)}{h(x)} \frac{e^{-(|x|^2+|y|^2)/2t}}{t^{m/2}} \det \left[\frac{x_i y_j}{t} {}_0\mathcal{F}_1 \left(\frac{3}{2}, \frac{(x_i y_j)^2}{4t^2} \right) \right]_{i,j} \end{aligned}$$

where h is the product over positive roots. Besides, the following holds (see [6]) :

$${}_0\mathcal{F}_1\left(\frac{3}{2}, z\right) = \frac{C}{2\sqrt{z}} \sinh(2\sqrt{z}).$$

Thus,

$$\begin{aligned} p_t^{1,1}(x, y) &= C_m \frac{h(y)}{h(x)} \frac{1}{(2\pi t)^{m/2}} e^{-(|x|^2+|y|^2)/2t} \det \left[\sinh \left(\frac{x_i y_j}{t} \right) \right]_{i,j} \\ &= \frac{h(y)}{h(x)} \det [N_t(y_j - x_i) - N_t(y_j + x_i)]_{i,j} \end{aligned}$$

where $N_t(u) = (1/\sqrt{2\pi t})e^{-u^2/2t}$, which fits Grabiner's result ([23] page 186). This is in agreement with the generator since $\Delta h = 0$ ([23]) and

$$\mathcal{L}f = \Delta f + \Gamma(\log h, f) = \Delta f + \sum_{i=1}^m \partial_i(\log h) \partial_i f,$$

where Γ is the so-called "opérateur du carré du champ" (see [37] Chap. VIII). Besides, for $m = 1$, r is a Bessel process of dimension $2\delta = 3$ and the expression inside the determinant in the second line is exactly the semi-group of the Brownian motion killed when it reaches 0 (see [37] p. 87).

5.2. Generalized Bessel function in the D_m case. In the classification of root systems, the A_{m-1} and B_m are known to be "irreducible" and both of them correspond to some matrix processes. Another one, yet with no underlying matrices, is the D_m root system defined by (see [27] p. 42)

$$R = \{\pm e_i \pm e_j, 1 \leq i < j \leq m\}, \quad R_+ = \{e_i \pm e_j, 1 \leq i < j \leq m\}$$

There is one conjugacy class so that $k(\alpha) = k_1$. Grabiner's result reads for the Brownian motion in the Weyl chamber of D_m -type ($k_1 = 1$) :

$$\begin{aligned} p_t^1(x, y) &= \frac{V(y^2) \det[N_t(y_i - x_j) - N_t(y_i + x_j)] + \det[N_t(y_i - x_j) + N_t(y_i + x_j)]}{V(x^2) \cdot 2} \\ &= \frac{C_m}{t^{\gamma+m/2}} e^{-(|x|^2+|y|^2)/2t} \frac{\det [\sinh(x_i y_j/t)] + \det [\cosh(x_i y_j/t)]}{V(x^2/4t^2)V(y^2)} V^2(y^2) \end{aligned}$$

where $\gamma = m(m-1)$. The second term in the sum involves the transition density of a reflected Brownian motion ($|B|$, see [37] p. 81). A natural way to interpret the announced formula is that the Weyl chamber is given by :

$$C = \{x \in \mathbb{R}^m, x_1 > \cdots > x_{m-1} > |x_m|\},$$

so that C fits the B_m -Weyl chamber when $x_m > 0$, otherwise, it is its conjugate with respect to s_{e_m} since this simple reflection acts only on x_m and retains the others. With the help of the determinantal formula used before ([24], [13]), ${}_0\mathcal{F}_1(\frac{3}{2}, z) = C \sinh(2\sqrt{z})/\sqrt{z}$ and ${}_0\mathcal{F}_1(1/2, z) = \cosh(2\sqrt{z})$ ([6]), one writes:

$$p_t^1(x, y) = \frac{e^{-(|x|^2+|y|^2)/2t}}{c_k t^{\gamma+m/2}} \left[\prod_{i=1}^m \left(\frac{x_i y_i}{2t} \right) {}_0F_1^{(1)} \left(m + \frac{1}{2}, \frac{x^2}{2t}, \frac{y^2}{2t} \right) + {}_0F_1^{(1)} \left(m - \frac{1}{2}, \frac{x^2}{2t}, \frac{y^2}{2t} \right) \right] V^2(y^2)$$

With regard to (1) and setting $q = 1 + (m-1)k_1$, it is natural to prove that :

Proposition 6.

$$\frac{1}{|W|} \sum_{w \in W} D_k(x, wy) = \prod_{i=1}^m \left(\frac{x_i y_i}{2} \right) {}_0F_1^{(1/k_1)} \left(q + \frac{1}{2}, \frac{x^2}{2}, \frac{y^2}{2} \right) + {}_0F_1^{(1/k_1)} \left(q - \frac{1}{2}, \frac{x^2}{2}, \frac{y^2}{2} \right).$$

Proof: it relies on the following characterization ([38]): given a reduced root system R with finite reflection group W , $D_k^W(x, \cdot)$ is the unique W -invariant function valued 1 at $x = 0$ satisfying $\Delta_k D_k^W(x, \cdot) = |x|^2 D_k^W(x, \cdot)$. It is easy to see that the function above is W -invariant since W is the semi-direct product of the symmetric group S_m and $(\mathbb{Z}/2\mathbb{Z})^{m-1}$ acting by even sign changes. However, it is not for the finite reflection group associated to the B_m root system due to the term multiplying the first hypergeometric series. In the D_m case, the Dunkl Laplacian writes on W -invariant functions :

$$\Delta_k = \sum_{i=1}^m \partial_i^2 + 2k_1 \sum_{i \neq j} \left[\frac{1}{y_i - y_j} + \frac{1}{y_i + y_j} \right] \partial_i$$

Let :

$$\begin{aligned} f(x, y) &= {}_0F_1^{(1/k_1)} \left(q + \frac{1}{2}, \frac{x^2}{2}, \frac{y^2}{2} \right) \\ g(x, y) &= {}_0F_1^{(1/k_1)} \left(q - \frac{1}{2}, \frac{x^2}{2}, \frac{y^2}{2} \right), \\ d(x, y) &= \prod_{i=1}^m (x_i y_i / 2) \end{aligned}$$

considered as functions of the variable y such that the generalized Bessel function is proportional to $df + g$ and $\Delta_k[df + g] = \Delta_k(df) + \Delta_k(g)$. Recall that:

$$\frac{1}{|W|} \sum_{w \in W} D_k^{(B_m)}(x, wy) = {}_0F_1^{(1/k_1)} \left(k_0 - \frac{1}{2} + q, \frac{x^2}{2}, \frac{y^2}{2} \right)$$

Note also that Δ_k is a particular case of the Dunkl Laplacian considered for the B_m root system when $k_0 = 0$. As a result

$$\Delta_k g(x, y) = \Delta_k^{(B_m)}(k_0 = 0) \left[\frac{1}{|W|} \sum_{w \in W} D_k^{(B_m)}(x, wy) \right] = \frac{1}{|W|} \langle x, x \rangle g(x, y)$$

For the remaining term, note that both d and f are W -invariant. Write $\Delta_k = \sum_{i=1}^m T_i^2$, where T_i is the difference Dunkl operator (see [38] p. 5). Then using (see [38] p. 6), one has the derivation formula $T_i(df) = dT_i(f) + fT_i(d)$. It gives that

$$\Delta_k(df) = d\Delta_k(f) + f\Delta_k(d) + 2 \sum_{i=1}^m (T_i(d))(T_i(f))$$

Moreover, $T_i(d) = \partial_i(d)$ and $T_i(f) = \partial_i(f)$ by W -invariance. Next we compute :

$$\Delta_k(d)(x, y) = 2k_1 d(x, y) \sum_{i \neq j} \frac{1}{y_i} \left[\frac{1}{y_i - y_j} + \frac{1}{y_i + y_j} \right] = 4k_1 d(x, y) \sum_{i \neq j} \frac{1}{y_i^2 - y_j^2} = 0$$

As a result :

$$\begin{aligned} \Delta_k(df)(x, y) &= [d\Delta_k(f)](x, y) + 2 \sum_{i=1}^m [(\partial_i(d))(\partial_i(f))](x, y) \\ &= [d\Delta_k(f)](x, y) + 2 \left[d \sum_{i=1}^m \frac{1}{y_i} (\partial_i(f)) \right](x, y) \\ &= d(x, y) \left[\Delta_k + 2 \sum_{i=1}^m \frac{1}{y_i} \partial_i \right] f(x, y) = d(x, y) \Delta_k^{(B_m)}(k_0 = 1) f(x, y) \end{aligned}$$

When $k_0 = 1$, f fits the generalized Bessel function in the B_m case $\Rightarrow \Delta_k(df)(x, y) = (1/|W|) \langle x, x \rangle df(x, y)$. Finally :

$$\Delta_k \left[\frac{1}{|W|} \sum_{w \in W} D_k(x, wy) \right] = \langle x, x \rangle \left[\frac{1}{|W|} \sum_{w \in W} D_k(x, wy) \right] \quad \blacksquare$$

6. ALCOVE-VALUED PROCESS

6.1. β -Jacobi processes. Recall that the eigenvalues of the real Jacobi matrix process of parameters (p, q) (see [16] for facts on this process) satisfy :

$$d\lambda_i(t) = 2\sqrt{\lambda_i(t)(1 - \lambda_i(t))} d\nu_i(t) + \left[(p - (p + q)\lambda_i(t)) + \sum_{j \neq i} \frac{\lambda_i(t)(1 - \lambda_j(t)) + \lambda_j(t)(1 - \lambda_i(t))}{\lambda_i(t) - \lambda_j(t)} \right] dt$$

for $0 < \lambda_m(0) < \dots < \lambda_1(0) < 1$ and all $t < \inf\{s > 0, \lambda_m(s) = 0 \text{ or } \lambda_1(s) = 1\} \wedge \tau$. The β -Jacobi process is defined as a solution, whenever it exists, of the SDE differing from the one above by a parameter $\beta > 0$ in front of the bracket. It is easy to see that if λ is a β -Jacobi process of parameters (p, q) , then $1 - \lambda$ is a β -Jacobi process of parameters (q, p) . As mentioned in the introductory part, the

connection with root systems is not new in its own ([3]) however we prefer giving some details of this transition. Setting $\lambda_i = \sin^2 \phi_i$ then $\phi_i = \arcsin \sqrt{\lambda_i} := s(\lambda_i)$ and $0 < \phi_m < \dots < \phi_1 < \pi/2$. The first and second derivatives of s are given by:

$$s'(\lambda_i) = \frac{1}{\sin 2\phi_i}, \quad s''(\lambda_i) = \frac{2(2\sin^2 \phi_i - 1)}{\sin^3 2\phi_i} = -\frac{2 \cos 2\phi_i}{\sin^3 2\phi_i}$$

Using

$$\sin^2 \phi_i - \sin^2 \phi_j = 2 \sin(\phi_i + \phi_j) \sin(\phi_i - \phi_j)$$

$$\sin^2 \phi_i \cos^2 \phi_j + \cos^2 \phi_i \sin^2 \phi_j = \sin^2(\phi_i + \phi_j) + \sin^2(\phi_i - \phi_j)$$

then, Ito's formula gives:

$$\begin{aligned} d\phi_i(t) &= d\nu_i(t) + \beta \frac{(p - (p + q) \sin^2 \phi_i)}{\sin 2\phi_i} - \cot 2\phi_i dt \\ &+ \frac{\beta}{2 \sin 2\phi_i(t)} \sum_{j \neq i} \frac{\sin^2(\phi_i(t) + \phi_j(t)) + \sin^2(\phi_i(t) - \phi_j(t))}{\sin(\phi_i(t) + \phi_j(t)) \sin(\phi_i(t) - \phi_j(t))} dt \end{aligned}$$

Writing $\sin^2 \phi_i = (1 - \cos 2\phi_i)/2$, then

$$\begin{aligned} d\phi_i(t) &= d\nu_i(t) + \beta \frac{(p - q)}{2} \frac{dt}{\sin 2\phi_i(t)} + \frac{\beta(p + q) - 2}{2} \cot 2\phi_i(t) dt \\ &+ \frac{\beta}{2 \sin 2\phi_i(t)} \sum_{j \neq i} \frac{\sin^2(\phi_i(t) + \phi_j(t)) + \sin^2(\phi_i(t) - \phi_j(t))}{\sin(\phi_i(t) + \phi_j(t)) \sin(\phi_i(t) - \phi_j(t))} dt \end{aligned}$$

where $0 < \phi_m(0) < \dots < \phi_1(0) < \pi/2$. Moreover,

$$\sin 2\phi_i = [\cot(\phi_i + \phi_j) + \cot(\phi_i - \phi_j)] \sin(\phi_i + \phi_j) \sin(\phi_i - \phi_j)$$

which gives

$$\begin{aligned} d\phi_i(t) &= d\nu_i(t) + \beta \frac{(p - q)}{2} \frac{dt}{\sin 2\phi_i(t)} + \frac{\beta(p + q) - 2}{2} \cot 2\phi_i(t) dt \\ &+ \frac{\beta}{2} \sum_{j \neq i} \frac{[1/\sin^2(\phi_i(t) + \phi_j(t))] + [1/\sin^2(\phi_i(t) - \phi_j(t))]}{\cot(\phi_i(t) + \phi_j(t)) + \cot(\phi_i(t) - \phi_j(t))} dt \end{aligned}$$

Using $1 + \cot^2 z = 1/\sin^2 z$, then

$$\begin{aligned}
d\phi_i(t) &= d\nu_i(t) + \beta \frac{(p-q)}{2} \frac{dt}{\sin 2\phi_i(t)} + \frac{\beta(p+q)-2}{2} \cot 2\phi_i(t) dt \\
&+ \frac{\beta}{2} \sum_{j \neq i} \frac{\cot^2(\phi_i(t) + \phi_j(t)) + \cot^2(\phi_i(t) - \phi_j(t)) + 2}{\cot(\phi_i(t) + \phi_j(t)) + \cot(\phi_i(t) - \phi_j(t))} dt \\
&= d\nu_i(t) + \beta \frac{(p-q)}{2} \frac{dt}{\sin 2\phi_i(t)} + \frac{\beta(p+q)-2}{2} \cot 2\phi_i(t) dt + \beta \times \\
&\sum_{i \neq j} \left\{ \frac{1 - \cot(\phi_i(t) + \phi_j(t)) \cot(\phi_i(t) - \phi_j(t))}{\cot(\phi_i(t) + \phi_j(t)) + \cot(\phi_i(t) - \phi_j(t))} + \frac{\cot(\phi_i(t) + \phi_j(t)) + \cot(\phi_i(t) - \phi_j(t))}{2} \right\} dt
\end{aligned}$$

Using

$$-\cot(u+v) = \frac{1 - \cot(u)\cot(v)}{\cot(u) + \cot(v)}.$$

and

$$\frac{1}{\sin 2\phi_i} = \frac{2 \cos^2 \phi_i - \cos 2\phi_i}{2 \sin \phi_i \cos \phi_i} = \cot \phi_i - \cot 2\phi_i$$

we finally obtain

(10)

$$d\phi_i(t) = d\nu_i(t) + \left[k_0 \cot \phi_i + k_1 \cot 2\phi_i(t) dt + k_2 \sum_{i \neq j} [\cot(\phi_i + \phi_j) + \cot(\phi_i - \phi_j)] \right] dt$$

where

$$(11) \quad 2k_0 = \beta(p-q), \quad k_1 = \beta(q - (m-1)) - 1, \quad 2k_2 = \beta.$$

Easy computations show that $\pi/2 - \phi$ satisfies (10) with (p, q) intertwined.

6.2. Eigenfunctions and Heckman-Opdam's functions. Let $k_2 > 0$ and \mathcal{L} be the generator of ϕ , then the eigenfunctions of \mathcal{L} are given by Gauss hypergeometric series: in fact, let \mathcal{A} be the generator of $(\lambda_1, \dots, \lambda_m)$ (see [16] p. 135 for $\beta = 1$):

$$\begin{aligned}
\mathcal{A} &= 2 \sum_{i=1}^m \lambda_i(1-\lambda_i) \partial_i^2 + \beta \sum_{i=1}^m \left[p - (p+q)\lambda_i + \sum_{j \neq i} \frac{\lambda_i(1-\lambda_j) + \lambda_j(1-\lambda_i)}{\lambda_i - \lambda_j} \right] \partial_i \\
&= 2 \sum_{i=1}^m \lambda_i(1-\lambda_i) \partial_i^2 + \beta \sum_{i=1}^m [p - (m-1) - (p+q-2(m-1))\lambda_i] \partial_i + 2\beta \sum_{i \neq j} \frac{\lambda_i(1-\lambda_i)}{\lambda_i - \lambda_j} \partial_i
\end{aligned}$$

From Equation (8) (k_2 plays the role of k_1), one can see that ${}_2F_1^{(1/k_2)}(a, b, c; \lambda)$ is the unique symmetric analytic function u such that $u(0) = 1$ which satisfies

$$\mathcal{A}u(\lambda) = 2mab u(\lambda), \quad 2c = \beta p = 2k_0 + k_1 + 2k_2(m-1) + 1, \quad 2a + 2b + 1 - 2c = k_1.$$

with k_i , $0 \leq i \leq 2$ cited in (11). Setting $\sin^2 \phi := (\sin^2 \phi_1, \dots, \sin^2 \phi_m)$, then \mathcal{A} transforms to \mathcal{L} . Hence:

$$\mathcal{L}[u(\sin^2 \phi)] = 2mab [u(\sin^2 \phi)]$$

In the same spirit, one can also interpret \mathcal{L} as the “radial part” of the trigonometric version Dunkl-Cherednik Laplacian (with \cot replacing \coth , [3],[36]). By radial part, we mean the restriction on W -invariant functions. Besides, this Laplacian arises, as for Dunkl and Cherednik-Dunkl ones, from differential-difference first-order operators. However, this comes beyond the spirit of this work and will not be done here.

6.3. Existence and uniqueness of a strong solution. The involved root system is the non reduced BC_m defined by

$$\begin{aligned} R &= \{\pm e_i, \pm 2e_i, 1 \leq i \leq m, \pm(e_i \pm e_j), 1 \leq i < j \leq m\} \\ R_+ &= \{e_i, 2e_i, 1 \leq i \leq m, (e_i \pm e_j), 1 \leq i < j \leq m\} \\ \Delta &= \{e_i - e_{i+1}, 1 \leq i \leq m-1, e_m\} \end{aligned}$$

When $k_0 = 0(p = q)$, it reduces to the reduced C_m system

$$\begin{aligned} R &= \{\pm e_i \pm e_j, 1 \leq i < j \leq m, \pm 2e_i, 1 \leq i \leq m\} \\ R_+ &= \{e_i \pm e_j, 1 \leq i < j \leq m, 2e_i, 1 \leq i \leq m\} \\ \Delta &= \{e_i - e_{i+1}, 1 \leq i \leq m-1, 2e_m\} \end{aligned}$$

and it is known as the ultraspheric case. The Weyl group action on \mathbb{R}^m gives rise to three orbits so that the multiplicity function is given by $k = (k_0, k_1, k_2)$. Setting $\tilde{\phi}_i := \phi_i/\pi$, then the process is valued in the *positive Weyl alcove* (see [27]) defined by :

$$\tilde{A} = \{\tilde{\phi} \in \mathbb{R}^m, \langle \alpha, \tilde{\phi} \rangle > 0 \forall \alpha \in \Delta, \langle \tilde{\alpha}, \tilde{\phi} \rangle < 1\}$$

where $\tilde{\alpha} = 2e_1$ is the highest positive root (that is $\tilde{\alpha} - \alpha \in R_+ \forall \alpha \in R$, see [27]). The associated affine Weyl group W_a is the semi-direct product of W and the translation group corresponding to the coroot lattice (\mathbb{Z} -span of $\{2\alpha/||\alpha||^2, \alpha \in R\}$). The generator writes in this case:

$$\mathcal{L}g(\phi) := \frac{1}{2}\Delta g(\phi) - \langle \nabla g(\phi), \nabla \Phi(\phi) \rangle, \quad \Phi(\phi) = - \sum_{\alpha \in R_+} k(\alpha) \log \sin(\langle \alpha, \phi \rangle)$$

Thus, with minor modifications, Theorem 1 states that (10) has a unique strong solution for all $t > 0$ subject to $k_0 > 0, k_1 > 0, k_2 > 0 \Leftrightarrow \beta > 0, p > q > (m-1)+1/\beta$. Applying this to $\pi/2 - \phi$, this holds for $\beta > 0, q > p > (m-1)+1/\beta$. Since the ultraspheric case still involves a root system, then (10) has a unique strong solution for $p \wedge q > (m-1) + 1/\beta$ which simplifies to $p \wedge q > m$ in the real case $\beta = 1$ and $p \wedge q > m - 1/2$ in the complex one $\beta = 2$. Theorem 1 is modified as follows: $\partial \tilde{A} = \cup_{\alpha \in \Delta} H_\alpha \cup H_{\alpha,1}$ where

$$H_{\alpha,1} = \{\tilde{\phi}, \langle \tilde{\alpha}, \tilde{\phi} \rangle = 1\} = \{\phi, \pi - \langle \tilde{\alpha}, \phi \rangle = 0\}$$

Compared with (4), the convex function $x \mapsto -\ln(\langle \alpha, x \rangle)$ should be substituted by $\phi \mapsto -\ln(\sin(\langle \alpha, \phi \rangle))$ and one has to deal with an additional term in the expression of the boundary process $(L_t)_{t \geq 0} : \mathbf{1}_{\{\pi - \langle \tilde{\alpha}, \phi \rangle = 0\}}$. Then the occupation density formula writes:

$$\begin{aligned} \int_0^{\pi/2} L_t^a(\pi - \langle \tilde{\alpha}, \phi \rangle) |\theta'(a)| da &= \langle \tilde{\alpha}, \tilde{\alpha} \rangle \int_0^t |\theta'(\pi - \langle \tilde{\alpha}, X_s \rangle)| ds \\ &= \langle \tilde{\alpha}, \tilde{\alpha} \rangle \int_0^t |\theta'(\langle \tilde{\alpha}, X_s \rangle)| ds \end{aligned}$$

since $\cot(\pi - z) = -\cot(z)$. Hence, the same proof applies and Lemma 1 remains valid for $\tilde{\alpha} \in R_+$. Besides, either it will exist $\alpha \in \Delta$ such that $\langle \alpha, x \rangle = 0$ and Lemma 2 applies, or we will need to prove that $\langle n(x), \tilde{\alpha} \rangle \neq 0$ if x belongs only to $H_{\tilde{\alpha}, 1}$. Let us first recall that the highest root is the unique positive root such that $\tilde{\alpha} - \alpha \in R_+$ for all $\alpha \in R_+$. Thus it may be written as $\tilde{\alpha} = \sum_{\alpha \in \Delta} a_\alpha \alpha$ where $a_\alpha \geq 1$. Else, if there exists $\alpha_0 \in \Delta$ such that $a_{\alpha_0} < 1$ and since $\tilde{\alpha}$ must be greater than all simple roots (in particular greater than α_0) then

$$\tilde{\alpha} - \alpha_0 = (a_{\alpha_0} - 1)\alpha_0 + \sum_{\alpha_0 \neq \alpha \in \Delta} a_\alpha \alpha = c_{\alpha_0} \alpha_0 + \sum_{\alpha_0 \neq \alpha \in \Delta} c_\alpha \alpha$$

for some $c_\alpha \geq 0$. Our claim follows from the fact that Δ is a basis. Next, it is not difficult to see from the definition of $n(x)$ and the fact that $\langle \alpha, x \rangle > 0$ for all $\alpha \in \Delta$ that $n(x)$ is colinear to $-\sum_{\alpha \in \Delta} \alpha$. It follows that

$$\langle n(x), \tilde{\alpha} \rangle = -c \sum_{\alpha \in \Delta} \langle \alpha, \tilde{\alpha} \rangle = -c \sum_{\alpha \in \Delta} \sum_{\theta \in \Delta} a_\alpha \langle \alpha, \theta \rangle$$

If $\langle n(x), \tilde{\alpha} \rangle = 0$, then

$$\left\| \sum_{\alpha \in \Delta} \alpha \right\|^2 = \sum_{\alpha \in \Delta} \sum_{\theta \in \Delta} \langle \alpha, \theta \rangle \leq \sum_{\alpha \in \Delta} \sum_{\theta \in \Delta} a_\alpha \langle \alpha, \theta \rangle = 0$$

which implies that $n(x) = 0$. ■

6.4. Brownian motion in Weyl alcoves. Let

$$h_1(\phi) := \prod_{\alpha \in R_+} \sin(\langle \alpha, \phi \rangle)$$

Then, h_1 is strictly positive on \tilde{A} and vanishes for $\phi \in \partial A$. One can also show that $(1/2)\Delta h_1 = ch$ for some strictly negative constant c . Let $P_t^{h_1}$ denote the semi group given by :

$$P_t^{h_1} f(\phi) := e^{-ct} \frac{P_t(h_1 f)(\phi)}{h_1(\phi)},$$

where P_t denotes the semi group of the process consisting of m -independent BMs in A killed when it first reaches ∂A . The corresponding generator writes :

$$\mathcal{L}^{h_1}(f) = \frac{1}{h_1} \left[\frac{1}{2} \Delta - c \right] (h_1 f) = \frac{1}{2} \Delta + \sum_{i=1}^m (\partial_i \log h_1) \partial_i f$$

which fits our generator for $k_2 = 1$ ($\beta = 2$), $k_1 = 2$ ($q = m + 1/2$), $k_0 = 1$ ($p = q + 1 = m + 3/2$). In the ultraspherical case, this becomes $\beta = 2$, $p = q = m + 1/2$. In both cases, these parameters correspond to the process consisting of m BMs constrained to stay in the BC_m and C_m - Weyl alcoves respectively. Note that p, q are not integers which means that these processes BM can not be realized as eigenvalues processes of complex matrix Jacobi processes which is also the case for the BM in the B_m -Weyl chamber since $\delta = m + 1/2$.

6.5. The first hitting time \tilde{T}_0 . We define similarly the first hitting time of alcove's walls by $\tilde{T}_0 = \inf\{t > 0, (\phi(t)/\pi) \in \partial \tilde{A}\} = \tilde{T}_{\tilde{\alpha}} \wedge \inf\{\tilde{T}_{\alpha}, \alpha \in \Delta\}$, where

$$\begin{aligned} \tilde{T}_{\alpha} &:= \inf\{t > 0, \langle \alpha, \phi(t) \rangle = 0\}, \\ \tilde{T}_{\tilde{\alpha}} &:= \inf\{t > 0, \langle \tilde{\alpha}, \phi(t) \rangle = 2\phi_1 = \pi\}, \end{aligned}$$

and ϕ is the unique strong solution for all $t \geq 0$ of ⁶:

$$d\phi(t) = d\nu(t) + \sum_{\alpha \in R_+} k(\alpha) \cot(\langle \alpha, \phi(t) \rangle) \alpha dt, \quad \frac{\phi(0)}{\pi} \in \tilde{A},$$

for the non reduced root system $R = BC_m$ with $k(\alpha) > 0$ for all α and $p \wedge q > (m - 1) + 1/\beta$. Let us focus on \tilde{T}_{α_0} for some $\alpha_0 \in \Delta$. We shall distinguish two cases :

6.5.1. $\alpha_0 = e_i - e_{i+1}$, $1 \leq i \leq m - 1$. The same scheme described in the proof of Proposition 1 applies here since the main ingredients used there are the SDE and the fact that $\sigma_0(\alpha) \in R_+$ if $\alpha \neq \alpha_0$. The second assertion follows from $\sigma_0(2e_j) = 2\sigma_0(e_j) = 2(\delta_{ij}e_{i+1} + \delta_{(i+1)j}e_i + \mathbf{1}_{\{j \neq i, j \neq i+1\}}e_j) \in R_+$. As a result, one writes for all $t \geq 0$:

$$d \langle \alpha_0, \phi(t) \rangle = \|\alpha_0\| d\gamma_t + k_2 \|\alpha_0\|^2 \cot \langle \alpha_0, \phi(t) \rangle dt + \sum_{\substack{\alpha \in R_+ \\ \alpha \neq \alpha_0}} k(\alpha) a(\alpha) \cot \langle \alpha, \phi(t) \rangle dt$$

where $a(\alpha) = \langle \alpha_0, \alpha \rangle$.

$$d \langle \alpha_0, \phi(t) \rangle = \|\alpha_0\| d\gamma_t + k_2 \|\alpha_0\|^2 \cot(\langle \alpha_0, \phi(t) \rangle) dt + F_t$$

where

$$F_t = \sum_{\substack{\alpha \in R_+ \setminus \alpha_0 \\ a(\alpha) > 0}} k(\alpha) a(\alpha) [\cot(\langle \alpha, \phi(t) \rangle) - \cot(\langle \sigma_0(\alpha), \phi(t) \rangle)],$$

⁶ $k(2e_i) = k_1/2$ for all $1 \leq i \leq m$.

where $\sigma_0 = \sigma_{\alpha_0}$. This drift is strictly negative on $\{\tilde{T}_{\alpha_0} = \infty\}$ since $\phi \mapsto \cot \phi$ is a decreasing function, $\langle \alpha_0, \phi(t) \rangle > 0$ and since :

$$\langle \alpha - \sigma_0(\alpha), \phi(t) \rangle = 2 \frac{a(\alpha)}{\|\alpha\|^2} \langle \alpha_0, \phi(t) \rangle > 0.$$

This implies that $\mathbb{P}_x(\forall t \geq 0, \langle \alpha_0, \phi(t) \rangle \leq Z_t) = 1$ where $\phi(0) = x$ and :

$$dZ_t = \|\alpha_0\| d\gamma_t + \|\alpha_0\|^2 k_2 \cot(Z_t) dt, \quad Z_0 = \langle \alpha_0, \phi(0) \rangle = x$$

on the same probability space. Using (10) with $\beta = 1, m = 1$, one can easily see that $(Z_t)_{t \geq 0} = (\arcsin \sqrt{J_{\|\alpha_0\|^2 t}})_{t \geq 0}$ where J is a one dimensional Jacobi process of parameters $d = 2k_2 + 1, d' = 1$ (see [43]) : that is :

$$dJ_t = 2\sqrt{J_t(1-J_t)} d\gamma_t + (d - (d+1)J_t) dt, \quad 0 < k_2 < 1/2 \Leftrightarrow 0 < d < 2.$$

As J hits 0 almost surely when $0 < d < 2$ (use the skew product in [43] and properties of squared Bessel processes), then so does Z and by the way $\langle \alpha_0, \phi \rangle$ for $k_2 < 1/2 \Rightarrow \tilde{T}_{\alpha_0} < \infty$ a. s.

6.5.2. $\alpha_0 = e_m$. Compared with the previous case, the difference arises from the fact that $\sigma_0(\alpha) \in R_+$ if $\alpha \in R_+ \setminus \{e_m, 2e_m\}$ and the latter is easily checked since for $\alpha = e_i \pm e_j$ this amounts to consider the reduced root system B_m , else for $\alpha = e_i$ with $i \neq m, \sigma_0(e_i) = e_i$. According to this, one gets :

$$d \langle \alpha_0, \phi(t) \rangle = d\phi_m(t) = d\gamma_t + k_0 \cot(\phi_m(t)) dt + k_1 \cot(2\phi_m) + F_t$$

where

$$F_t = \sum_{\substack{\alpha \in R_+ \setminus \{e_m, 2e_m\} \\ a(\alpha) > 0}} k(\alpha) a(\alpha) [\cot(\langle \alpha, \phi(t) \rangle) - \cot(\langle \sigma_0(\alpha), \phi(t) \rangle)]$$

where $R_+^1 = \{e_i - e_j, 1 \leq i < j \leq m\}$. Using once again (10), we shall compare this process with $(\arcsin \sqrt{J_t})_{t \geq 0}$ where

$$dJ_t = 2\sqrt{J_t(1-J_t)} d\gamma_t + (d - (d+d')J_t) dt, \quad d' = k_1 + 1, d = 2k_0 + k_1 + 1.$$

Hence, $\tilde{T}_{e_m} < \infty$ a.s. if $0 < 2k_0 + k_1 < 1/2 \Leftrightarrow \beta p - (\beta(m-1)) < 2$. This agrees with the case $m = 1$ for which $p < 2$ (use the skew product in [43]). Finally, note that since $a(\alpha) = 0$ for $\alpha \in \{e_i, 2e_i, 1 \leq i \leq m-1\}$, F only involves $k_2 = \beta$ which is independent from p, q . Keeping in mind that $\pi/2 - \phi$ is still a β -Jacobi process with (p, q) intertwined which has no effect on the strict negativity of F by the above remark, we conclude that $\tilde{T}_{\tilde{\alpha}} < \infty$ for $0 < \beta q - \beta(m-1) < 2$. \blacksquare

6.6. Semi-group density. We end this paper by giving the semi group density of the β -Jacobi process. Before proceeding, we briefly consider two cases for which we can write down the semi-group density: the univariate case and the complex Hermitian one ($\beta = 2$). Let $P_n^{r,s}$ denote the Jacobi polynomial of degree n defined by ([1]):

$$P_n^{r,s}(\lambda) := \frac{(r+1)_n}{n!} {}_2F_1 \left(-n, n+r+s+1, r+1; \frac{1-\lambda}{2} \right),$$

for $\lambda \in [-1, 1]$, $r, s > -1$, where ${}_2F_1$ is the univariate Gauss hypergeometric function. These polynomials are orthogonal with respect to the measure $Z^{r,s}(\lambda)d\lambda := (1-\lambda)^r(1+\lambda)^s d\lambda$ and the associated inner product in $L^2([-1, 1])$ given by

$$\langle f, g \rangle_{L^2([-1,1])} := \int_{[-1,1]} f(\lambda)g(\lambda)Z^{r,s}(\lambda)d\lambda$$

Moreover, $(P_n^{r,s})_{n \geq 0}$ form a complete set of this Hilbert space and satisfy

$$\left\{ \sqrt{1-\lambda^2} \partial_\lambda^2 + [(s-r) - (s+r+2)\lambda] \partial_\lambda \right\} P_n^{r,s}(\lambda) = -n(n+r+s+1)P_n^{r,s}(\lambda)$$

The above eigenoperator defines a diffusion which is related to the one we considered with $m = 1$ via the map $\lambda \mapsto (1-\lambda)/2$ and a deterministic time change ($t \mapsto t/2$). The semi group density w.r.t Lebesgue measure is written (see [44],[14])

$$p_t^{r,s}(\theta, \lambda) = \sum_{n=0}^{\infty} e^{-2r_n t} P_n^{r,s}(\theta) P_n^{r,s}(\lambda) W^{r,s}(\lambda)$$

where r_n denotes the eigenvalues above, $(P_n^{r,s})_n$ are orthonormal polynomials, $p = 2(r+1)$, $q = 2(s+1)$ and $W^{r,s}(\lambda)d\lambda$ is the probability measure corresponding to the measure $Z^{r,s}(\lambda)d\lambda$. No closed forms seems to be known for this density, nonetheless an attempt to get a handier expression was tried in [14]. Multivariate analogs appeared in literature ([3], [28], [32] for instance) and are obtained by applying the Gram-Schmidt orthogonalization to the symmetric Jack polynomials w.r.t. measure

$$Z_m^{r,s,\beta}(\lambda)d\lambda := \prod_{i=1}^m \lambda_i^r (1-\lambda_i)^s \prod_{1 \leq i < j \leq m} |\lambda_i - \lambda_j|^\beta d\lambda_1 \dots d\lambda_m$$

We shall denote them⁷ by $P_\tau^{r,s,\beta}$ for a given partition τ (instead of $G_\tau^{\alpha,\beta}$ used in literature) and stress that some of the properties cited above extend to the higher dimensional case ([32]): an expansion in terms of ${}_2F_1^{(2/\beta)}(-l, b, c, \lambda)$ exists for $\tau = (l^m)$ with m components all equal to l ; $(P_\tau^{r,s,\beta})$, where τ is a partition of length $\leq m$, form a basis of the Hilbert space $L^2([0, 1]^m, W_m^{r,s,\beta}(\lambda)d\lambda)$ where $W_m^{r,s,\beta}(\lambda)d\lambda$ is the normalization $Z_m^{r,s,\beta}(\lambda)d\lambda$ in order to be a probability measure ([32]). The normalizing constant is given by a McDonald-Selberg integral computed in [28].

⁷The normalization is different from the one used in both [3] and [32].

Moreover, $(P_\tau^{r,s,\beta})_\tau$ are the unique symmetric polynomial eigenfunctions of the Laplace Beltrami operator $-\mathcal{L}$ (thus defined on $[0, 1]^m$) with $\beta(p - (m - 1)) = 2(r + 1)$, $\beta(q - (m - 1)) = 2(s + 1)$, associated with the eigenvalues

$$(12) \quad 2r_{n,\tau}^\beta := 2 \left[\sum_{i=1}^m \tau_i(\tau_i - 1 - \beta(i - 1)) + |\tau|(r + s + \beta(m - 1) + 2) \right], \quad |\tau| = n.$$

However, with regard to the strong uniqueness for all $t \geq 0$ previously derived, we shall restrict ourselves to $p \wedge q > (m - 1) + 1/\beta$. $\beta(q - (m - 1)) > 1$ is equivalent to $s > -1/2$ and $\beta(p - (m - 1)) > 1$ is equivalent to $r > -1/2$. As a result, $r, s > -1/2$.

It is known that the eigenvalues process of the complex Hermitian Jacobi process (or 2-Jacobi process) is the h -transform (in the Doob sense) for $h = V$ of a process whose components are real Jacobi processes of parameters $2(p - (m - 1)) = 2(r + 1)$, $2(q - (m - 1)) = 2(s + 1)$ constrained to never collide (or to stay in the A_{m-1} -type Weyl chamber). Here, V denotes as usual the Vandermonde function. More precisely, V is an eigenfunction of the generator of the one dimensional Jacobi process of parameters (p, q) (see appendix in [16]), say \mathbf{L} , that is

$$\mathbf{L}V = cV = -m(m - 1) \left(\frac{2(m - 2)}{3} + \frac{p + q}{2} \right) V$$

Noting that the parameters r, s are the same both in the univariate and in the multivariate cases, it follows by Karlin-McGregor formula ([29]) that the semi group density writes on $\{0 < \lambda_m < \dots < \lambda_1 < 1\}$

$$\begin{aligned} K_t^{r,s,2}(\theta, \lambda) &:= e^{-ct} \frac{V(\lambda)}{V(\theta)} \det \left(\sum_{n=0}^{\infty} e^{-2n(n+r+s+1)t} P_n^{r,s}(\theta_i) P_n^{r,s}(\lambda_j) W^{r,s}(\lambda_j) \right)_{i,j} \\ &= e^{-ct} \det \left(\sum_{n=0}^{\infty} e^{-2n(n+r+s+1)t} P_n^{r,s}(\theta_i) P_n^{r,s}(\lambda_j) \right)_{i,j} \frac{W_m^{r,s,2}(\lambda)}{V(\theta)V(\lambda)} \\ &= e^{-ct} \left[\sum_{\sigma_1 \in S_m} \epsilon(\sigma_1) \sum_{n_1, \dots, n_m \geq 0} e^{-2 \sum_{i=1}^m n_i(n_i+r+s+1)t} \prod_{i=1}^m P_{n_i}^{r,s}(\theta_i) P_{n_i}^{r,s}(\lambda_{\sigma_1(i)}) \right] \frac{W_m^{r,s,2}(\lambda)}{V(\theta)V(\lambda)} \\ &= e^{-ct} \left[\sum_{\sigma_1, \sigma_2 \in S_m} \epsilon(\sigma_1) \sum_{n_1 \geq \dots \geq n_m \geq 0} e^{-2 \sum_{i=1}^m n_{\sigma_2(i)}(n_{\sigma_2(i)}+r+s+1)t} \prod_{i=1}^m P_{n_{\sigma_2(i)}}^{r,s}(\theta_i) P_{n_{\sigma_2(i)}}^{r,s}(\lambda_{\sigma_1(i)}) \right] \frac{W_m^{r,s,2}(\lambda)}{V(\theta)V(\lambda)} \end{aligned}$$

Note that, for a given partition $(n_1 \geq \dots \geq n_m \geq 0)$ and a permutation $\sigma_2 \in S_m$, one has

$$\sum_{i=1}^m n_{\sigma_2(i)}(n_{\sigma_2(i)} + r + s + 1) = \sum_{i=1}^m n_i(n_i + r + s + 1)$$

Thus summing first over σ_1 with the change of variables $\sigma = \sigma_1 \sigma_2$, one gets:

$$\begin{aligned} K_t^{r,s,2}(\theta, \lambda) &= e^{-ct} \sum_{n_1 \geq \dots \geq n_m \geq 0} e^{-2 \sum_{i=1}^m n_i(n_i+r+s+1)t} \frac{\det[P_{n_i}^{r,s}(\theta_j)]_{i,j}}{V(\theta)} \frac{\det[P_{n_i}^{r,s}(\lambda_j)]_{i,j}}{V(\lambda)} W_m^{r,s,2}(\lambda) \\ &= e^{-ct} \sum_{n_1 > \dots > n_m \geq 0} e^{-2 \sum_{i=1}^m n_i(n_i+r+s+1)t} \frac{\det[P_{n_i}^{r,s}(\theta_j)]_{i,j}}{V(\theta)} \frac{\det[P_{n_i}^{r,s}(\lambda_j)]_{i,j}}{V(\lambda)} W_m^{r,s,2}(\lambda) \end{aligned}$$

Set $n_i = \tau_i + m - i$, then $\tau_1 > \dots > \tau_m \geq 0$. Moreover, with regard to (12), one easily check that

$$r_{n,\tau}^2 = \sum_{i=1}^m \tau_i(\tau_i + r + s + 1 + 2(m - i))$$

so that

$$\sum_{i=1}^m n_i(n_i + r + s + 1) = r_{n,\tau}^2 - c/2$$

The final result writes

$$\begin{aligned} K_t^{r,s,2}(\theta, \lambda) &= \sum_{\tau_1 \geq \dots \geq \tau_m \geq 0} e^{-2r_{n,\tau}^2} \frac{\det[P_{\tau_i+m-i}^{r,s}(\theta_j)]_{i,j}}{V(\theta)} \frac{\det[P_{\tau_i+m-i}^{r,s}(\lambda_j)]_{i,j}}{V(\lambda)} W_m^{r,s,2}(\lambda) \\ &= \sum_{\tau_1 \geq \dots \geq \tau_m \geq 0} e^{-2r_{n,\tau}^2} P_\tau^{r,s,2}(\theta) P_\tau^{r,s,2}(\lambda) W_m^{r,s,2}(\lambda) \end{aligned}$$

where we used the determinantal representation of the Jacobi multivariate polynomials in the complex case⁸ (see [32]) :

$$P_\tau^{r,s,2}(\lambda) = \frac{\det[P_{\tau_i+m-i}^{r,s}(\lambda_j)]_{i,j}}{V(\lambda)}$$

From these observations, it is natural to claim that :

Proposition 7. *The semi group density of the β -Jacobi process is given by*

$$(13) \quad K_t^{r,s,\beta}(\theta, \lambda) := \sum_{n=0}^{\infty} \sum_{|\tau|=n} e^{-r_{n,\tau} t} P_\tau^{r,s,\beta}(\theta) P_\tau^{r,s,\beta}(\lambda) W_m^{r,s}(\lambda) \mathbf{1}_{\{0 < \lambda_m < \dots < \lambda_1 < 1\}}$$

with respect to $d\lambda$. As a result, it is positive.

Proof: given a bounded symmetric function f on $[0, 1]^m$, define

$$T_t f(\theta) := \int_{0 < \lambda_m < \dots < \lambda_1 < 1} f(\lambda) \sum_{n=0}^{\infty} \sum_{|\tau|=n} e^{-2r_{n,\tau}^\beta t} P_\tau^{r,s,\beta}(\theta) P_\tau^{r,s,\beta}(\lambda) W_m^{r,s}(\lambda) d\lambda$$

for $\theta = (0 < \theta_1 < \dots < \theta_m < 1)$ and $T_0 f = f$. The above expression makes sense: this uses the boundness of f , the exponential term with strictly positive t and Fubini Theorem. Besides, $T_t \mathbf{1} = 1$ and $\|T_t\|$ is bounded for all $t \geq 0$. The first

⁸We adopt a different normalization since we consider orthonormal polynomials.

claim follows easily from the orthogonality of the Jacobi polynomials and $P_0 = 1$ so that the only non zero term is that corresponding to $n = 0$. The second one is obvious for $t = 0$ and uses the exponential term when $t \geq \epsilon > 0$. One also easily checks that $T_t T_s = T_{t+s}$ and that $\mathcal{L}T_t f(\lambda) = \partial_t T_t f(\lambda)$ using the dominated convergence theorem. Now, let us consider the Cauchy problem associated to \mathcal{L} :

$$\begin{cases} \frac{\partial u_f}{\partial t}(t, \lambda) = \mathcal{L}u_f(t, \lambda) \\ u_f(0, \cdot) = f, \end{cases}$$

where $u_f \in C^{1,2}(\mathbb{R}_+^* \times S := \{0 < \lambda_m < \dots < \lambda_1 < 1\}) \cap C_b(\mathbb{R}^+ \cap S)$ with reflecting boundary condition :

$$\langle \nabla u(t, \lambda), n(\lambda) \rangle = 0, \quad (t, \lambda) \in \mathbb{R}_+^* \times \partial S$$

where $n(\lambda)$ is a unitary inward normal vector at λ . Define $u_t(f)(\lambda) := u_f(t, \lambda)$. It is shown ([42]) that there is a unique solution to the Cauchy problem with initial condition. Consequently, $(T_t)_{t \geq 0}$ is the semi group of the eigenvalues process $(\lambda(t))_{t \geq 0}$ with density given by $K_t^{r,s,\beta}$. ■

Remark. As the reader can check, the computations performed in the complex Hermitian case do not restrict to Jacobi polynomials. We only used the determinantal representation in terms of their univariate counterparts. As a result, one gets similar formulas replacing Jacobi by Hermite and Laguerre polynomials.

Now, we are able to answer some open questions left in [16]. For the real Jacobi matrix ($\beta = 1$), it is known that for $p \wedge q \geq m - 1$ and if the eigenvalues are distinct at time $t = 0$, then they remain distinct forever. It is then natural to wonder if this remains valid when starting from non distinct eigenvalues (see [16] p. 138-139). The Markov property together with the previous result for distinct eigenvalues are sufficient to claim that this is true provided that the eigenvalues semi group has a density which is absolutely continuous with respect to Lebesgue measure on \mathbb{R}^m . By virtue of $K_t^{r,s,1}(\theta, \phi)$, for $p \wedge q > m$,

$$\mathbb{P}_{\lambda(0)}(\forall t \geq 0, \forall i \neq j, \lambda_i(t) \neq \lambda_j(t)) = 1, \quad \lambda_1(0) \geq \dots \geq \lambda_m(0).$$

We argue in the same way to claim that for $p \wedge q \geq m + 1$, the process will never hit the boundaries (0 and 1 for λ or 0 and $\pi/2$ for ϕ) even if it did at time $t = 0$.

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REFERENCES

- [1] *G. E. Andrews, R. Askey, R. Roy.* Special functions. *Cambridge University Press.* 1999.
- [2] *T. H. Baker, P. J. Forrester.* The Calogero-Sutherland model and generalized classical polynomials. *Comm. Math. Phys.* **188.** 1997, 175-216.

- [3] *R. J. Beerends, E. M. Opdam.* Certain hypergeometric series related to the root system BC . *Trans. Amer. Math. Soc.* **339**, no. 2. 1993, 581-607.
- [4] *N. Bourbaki.* Groupes et Algèbres de Lie, chapitres IV-V-VI. *Hermann, Paris*, 1968.
- [5] *M. F. Bru.* Wishart Processes. *J. Theoretical Probability*, **4**, no. 4, 1991, 725 -751.
- [6] *Yu.A.Brychkov, O.I.Marichev, A.P.Prudnikov.* Integrals and Series, vol 3: more special functions. *Gordon and Breach science publishers*.
- [7] *R. Carter.* Lie Algebras of Finite and Affine Type. *Cambridge University Press*. 2005.
- [8] *E. Cépa, D. Lépingle.* Brownian particles with electrostatic repulsion on the circle : Dyson's model for unitary random matrices revisited. *E. S. A. I. M : Probability and Statistics*. **5**. 2001, 203-224.
- [9] *E. Cépa, D. Lépingle.* Diffusing particles with electrostatic repulsions. *P. T. R. F.* **107**. 1997, 429-449.
- [10] *E. Cépa.* Equations différentielles stochastiques multivoques. *Sém. Proba.* **XXIX**. 1995, 86-107.
- [11] *O. Chybiryakov.* Processus de Dunkl et Relation de Lamperti. Ph. D. Thesis, Paris VI Univ. June 2005.
- [12] *I. Cherednik, P. J. Forrester, D. Uglov.* Random matrices, Log-gases and the Calogero-Sutherland model. *Quantum Many-body Problems and Representation Theory*. MSJ Memoirs. 1998, 97-181.
- [13] *N. Demni.* Laguerre process and generalized Hartman-Watson law. To appear in *Bernoulli Journal*.
- [14] *N. Demni, M. Zani.* Large deviations for statistics of Jacobi process. Submitted.
- [15] *C.Donati-Martin, Y.Doumerc, H.Matsumoto, M.Yor.* Some properties of Wishart process and a matrix extension of the Hartman-Watson law. *Pub. R. I. M. S. Kyoto university*, **40**, No 4, 2004, 1385-1412.
- [16] *Y. Doumerc.* Matrix Jacobi Process. *Ph. D. Thesis*. Paul Sabatier Univ. May 2005.
- [17] *I. Dumitriu, A, Edelman.* Eigenvalue statistics for Beta-ensembles. *J. Math. Phys.* **43**, (11). 2002, 5830-5847.
- [18] *C. F. Dunkl.* Differential-difference operators associated to reflections groups. *Trans. Amer. Math. Soc.* **311**, no. 1, 1989,167-183.
- [19] *E. B. Dynkin.* Markov Processes, vol. 1. *Springer-Verlag*, 1965.
- [20] *F. J. Dyson.* A Brownian motion model for the eigenvalues of a random matrix. *J. Math. Phys.* **3**. 1962, 1191-1198.
- [21] *L. Gallardo, M. Yor.* A chaotic representation property of the multidimensional Dunkl process. *Annals of Probability* **34**, no. 4, 2006, 1530-1549.
- [22] *L. Gallardo, M. Yor.* Some remarkable properties of the Dunkl martingales. Séminaire de Probabilités **XXXIX**, dedicated to Paul-André Meyer. *Lecture notes in Math*. Springer, Berlin. 2005.
- [23] *D. J. Grabiner.* Brownian motion in a Weyl chamber, non-colliding particles and random matrices. *Ann. IHP.* **35**, 1999, no. 2, 177-204.
- [24] *K. I. Gross, St. P. Richards.* Special functions of matrix argument, *Bull. Amer. Math. Soc.* **24**, no 2, 1991, 349-355.
- [25] *G. J, Heckman, E. M, Opdam.* Root systems and Hypergeometric functions I. *Comp. Math.* **64**. 1987, 329-352.
- [26] *D. Hobson, W. Werner.* Non-colliding Brownian motions on the circle. *Bull. London Math. Society.* **28**. 1996, 643-650.
- [27] *J. E. Humphreys.* Reflections Groups and Coxeter Groups. *Cambridge University Press.* **29**. 2000.

- [28] *J. Kaneko*. Selberg integrals and hypergeometric functions with Jack polynomials. *SIAM J. Math. Anal.* **24**, 1993, 1086-1110.
- [29] *I. Karatzas, S. E. Shreve*. Brownian Motion and Stochastic Calculus, 2nd edition. *Springer-Verlag, New York*. 1991.
- [30] *M. Katori, H. Tanemura*. Noncolliding Brownian motions and Harish-Chandra formula. *Elec. Comm. Proba.* **8**. 2003, 112-121.
- [31] *W. König, N. O'Connell*. Eigenvalues of the Laguerre process as non-colliding squared Bessel processes. *Elec. Comm. in Proba.* **6**, 2001, 107-114.
- [32] *M. Lassalle*. Polynômes de Jacobi généralisés. *C. R. A. S. Paris.* **312**. Série **I**, 1991, 425-428.
- [33] *N.N.Lebedev*. Special Functions And Their Applications. *Dover Publications, INC*. 1972.
- [34] *L. G. Mcdonald*. Symmetric Functions and Hall Polynomials, 2nd ed. *Clarendon Press, Oxford*, 1995.
- [35] *R . J. Muirhead*. Aspects of Multivariate Statistical Theory. *Wiley Series in Probability*, 1982.
- [36] *E. Opdam*. Lectures on Dunkl Theory. *Available on arXiv*.
- [37] *D. Revuz, M. Yor*. Continuous Martingales and Brownian Motion. 3rd ed. *Springer*, 1999.
- [38] *M. Rösler*. Dunkl operator : theory and applications, Orthogonal polynomials and special functions (Leuven, 2002). *Lecture Notes in Math.* Vol. 1817, Springer, Berlin, 2003, 93-135.
- [39] *M. Rösler, M. Voit*. Markov processes related with Dunkl operators. *Adv. in Appl. Math.* **21**, no. 4, 1998, 575-643.
- [40] *M. Rösler*. Generalized Hermite polynomials and the heat equation for Dunkl operators. *Comm. Math. Phys.* **192**, no. 3, 1998, 519-542.
- [41] *B. Schapira*. The Heckman-Opdam Markov processes. To appear in *P. T. R. F.*
- [42] *D. W. Stroock, S.S. Varadhan*. Diffusion processes with boundary conditions. *Comm. Pure. Appl. Math.*, **XXIV**, 1971, 147-225.
- [43] *J. Warren, M. Yor*. The Brownian Burglar : Conditioning Brownian motion by its local time process. *Sém. Proba.* **XXXII.**, 1998, 328-342.
- [44] *E. Wong*. The construction of a class of stationary Markov. *Proceedings. The 16th Symposium. Applied Math.* AMS. Providence. RI. 1964. 264-276.