

GENERALIZED LAGUERRE UNITARY ENSEMBLES AND AN INTERACTING PARTICLES MODEL WITH A WALL

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Abstract

We introduce and study a new interacting particles model with a wall and two kinds of interactions - blocking and pushing - which maintain particles in a certain order. We show that it involves a random matrix model.

1 Interacting particles model

Let us consider k ordered particles evolving in discrete time on the positive real line with interactions that maintain their orderings. The particles are labeled in increasing order from 1 to k . Thus for $t \in \mathbb{N}$, we have

$$0 \leq X_1(t) \leq \dots \leq X_k(t),$$

where $X_i(t)$ is the position of the i^{th} particle at time $t \geq 0$. Particles are initially all at 0. The particles jump at times $n - \frac{1}{2}$ and n , $n \in \mathbb{N}^*$. Let us consider two independent families

$$\left(\xi\left(i, n - \frac{1}{2}\right)\right)_{i=1, \dots, k; n \geq 1}, \quad \text{and} \quad \left(\xi(i, n)\right)_{i=1, \dots, k; n \geq 1},$$

of independent random variables having an exponential law of mean 1. For convenience, we suppose that there is a static particle which always stays at 0. We call it the 0^{th} particle, and denote $X_0(t)$ its position at time $t \geq 0$.

At time $n - 1/2$, for $i = 1, \dots, k$, in that order, the i^{th} particle tries to jump to the left according to a jump of size $\xi(i, n - \frac{1}{2})$ being blocked by the old position of the $(i - 1)^{\text{th}}$ particle. In other words :

- If $X_1(n - 1) - \xi(1, n - \frac{1}{2}) < 0$, then the 1^{st} particle is blocked by 0, else it jumps at $X_1(n - 1) - \xi(1, n - \frac{1}{2})$, i.e.

$$X_1\left(n - \frac{1}{2}\right) = \max\left(0, X_1(n - 1) - \xi\left(1, n - \frac{1}{2}\right)\right).$$

- For $i = 2, \dots, k$, if $X_i(n-1) - \xi(i, n - \frac{1}{2}) < X_{i-1}(n-1)$, then the i^{th} particle is blocked by $X_{i-1}(n-1)$, else it jumps at $X_i(n-1) - \xi(i, n - \frac{1}{2})$, i.e.

$$X_i(n - \frac{1}{2}) = \max(X_{i-1}(n-1), X_i(n-1) - \xi(i, n - \frac{1}{2})).$$

At time n , particles jump successively to the right according to an exponentially distributed jump of mean 1, pushing all the particles to their right. The order in which the particles jump is given by their labels. Thus for $i = 1, \dots, k$, if $X_{i-1}(n) > X_i(n - \frac{1}{2})$ then the i^{th} particle is pushed before jumping, else it jumps from $X_i(n - \frac{1}{2})$ to $X_i(n - \frac{1}{2}) + \xi(i, n)$:

$$X_i(n) = \max(X_{i-1}(n), X_i(n - \frac{1}{2})) + \xi(i, n).$$

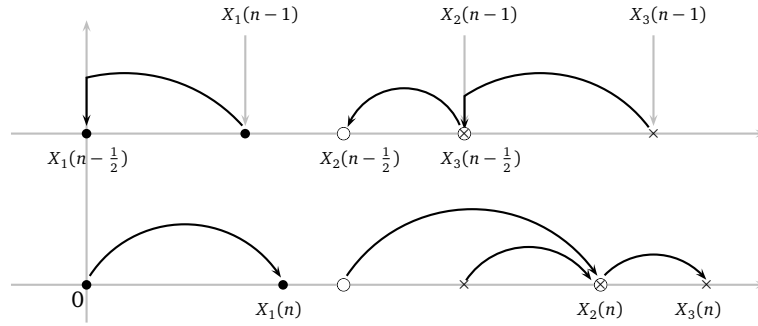


Figure 1: An example of blocking and pushing interactions between times $n - 1$ and n for $k = 3$. Different particles are represented by different kinds of dots.

There is another description of the dynamic of this model which is equivalent to the previous. At each time $n \in \mathbb{N}^*$, each particle successively attempts to jump first to the left then to the right, according to independent exponentially distributed jumps of mean 1. The order in which the particles jump is given by their labels. At time $n \in \mathbb{N}^*$, the 1^{st} particle jumps to the left being blocked by 0 then immediately to the right, pushing the i^{th} particles, $i = 2, \dots, k$. Then the second particle jumps to the left, being blocked by $\max(X_1(n-1), X_1(n))$, then to the right, pushing the i^{th} particles, $i = 3, \dots, k$, and so on. In other words, for $n \in \mathbb{N}^*$, $i = 1, \dots, k$,

$$X_i(n) = \max(X_{i-1}(n), X_{i-1}(n-1), X_i(n-1) - \xi^-(i, n)) + \xi^+(i, n), \quad (1)$$

where $(\xi^-(i, n))_{i=1, \dots, k; n \in \mathbb{N}}$, and $(\xi^+(i, n))_{i=1, \dots, k; n \in \mathbb{N}}$ are two independent families of independent random variables having an exponential law of mean 1.

2 Results

Let us denote $\mathcal{M}_{k,m}$ the real vector space of $k \times m$ real matrices. We put on it the Euclidean structure defined by the scalar product,

$$\langle M, N \rangle = \text{tr}(MN^*), \quad M, N \in \mathcal{M}_{k,m}.$$

Our choice of the Euclidean structure above defines a notion of standard Gaussian variable on $\mathcal{M}_{k,m}$. Thus a standard Gaussian random variable on $\mathcal{M}_{k,m}$ has a density with respect to the Lebesgue measure on $\mathcal{M}_{k,m}$

$$M \in \mathcal{M}_{k,m} \mapsto \frac{1}{\sqrt{km}2\pi} \exp\left(-\frac{1}{2}\text{tr}(MM^*)\right).$$

We write \mathcal{A}_k for the set $\{M \in \mathcal{M}_{k,k} : M + M^* = 0\}$ of antisymmetric $k \times k$ real matrices, and $i\mathcal{A}_k$ for the set $\{iM : M \in \mathcal{A}_k\}$. Since a matrix in $i\mathcal{A}_k$ is Hermitian, it has real eigenvalues $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_k$. Moreover, antisymmetry implies that $\lambda_{k-i+1} = -\lambda_i$, for $i = 1, \dots, [k/2] + 1$, in particular $\lambda_{[k/2]+1} = 0$ when k is odd.

Our main result is that the positions of the particles of our interacting particles model can be interpreted as eigenvalues of a random walk on the set of matrices $i\mathcal{A}_{k+1}$.

Theorem 2.1. *Let k be a positive integer and $(M(n), n \geq 0)$, be a discrete process on $i\mathcal{A}_{k+1}$ defined by*

$$M(n) = \sum_{l=1}^n Y_l \begin{pmatrix} 0 & i \\ -i & 0 \end{pmatrix} Y_l^*,$$

where the Y_l 's are independent standard Gaussian variables on $\mathcal{M}_{k+1,2}$. For $n \in \mathbb{N}$, let $\Lambda_1(n)$ be the largest eigenvalue of $M(n)$. Then the processes

$$(\Lambda_1(n), n \geq 0) \text{ and } (X_k(n), n \geq 0),$$

have the same distribution.

For a matrix $M \in i\mathcal{A}_{k+1}$ and $m \in \{1, \dots, k+1\}$, the main minor of order m of M is the submatrix

$$\{M_{ij}, 1 \leq i, j \leq m\}.$$

Theorem 2.2. *Let $M(n), n \geq 0$, be a discrete process on $i\mathcal{A}_{k+1}$, defined as in theorem 2.1. For $n \in \mathbb{N}$, $m = 2, \dots, k+1$, we denote $\Lambda_1^{(m)}(n)$, the largest eigenvalue of the main minor of order m of $M(n)$. Then for each fixed $n \in \mathbb{N}^*$, the random vectors*

$$(\Lambda_1^{(2)}(n), \dots, \Lambda_1^{(k+1)}(n)) \quad \text{and} \quad (X_1(n), \dots, X_k(n)),$$

have the same distribution.

Let us notice that there already exists a version of theorem 2.1 with particles jumping by one (see [2], or section 2.3 of [12]) and a continuous version which involves reflected Brownian motions with a wall and Brownian motions conditioned to never collide with each other or the wall (see [1]).

There is a variant of our model with no wall and no left-jumps which has been extensively studied (see for instance Johansson [6], Dieker and Warren [4], or Warren and Windridge [12]). It involves random matrices from the Laguerre Unitary Ensemble. Indeed, in that case, the position of the rightmost particle has the same law as the largest eigenvalue of the process $(M(n), n \geq 0)$ defined by

$$M(n) = \sum_{l=1}^n Z_l Z_l^*, \quad n \geq 0,$$

where the $\sqrt{2}Z_i$'s are independent standard Gaussian variables on \mathbb{C}^k . Let us mention that this model without a wall is equivalent to a maybe more famous one called the TASEP (totally asymmetric simple exclusion process). In this last model, we consider infinitely many ordered particles evolving on \mathbb{Z} as follows. Initially there is one and only one particle on each point of \mathbb{Z}_- . Particles are equipped with independent Poisson clocks of intensity 1 and each of them jumps by one to the right only if its clock rings and the point just to its right is empty. Particles are labeled by \mathbb{N}^* from the right to the left. Then, the time needed for the k^{th} particles to make n jumps is exactly the position of the rightmost particle at time n in the model without a wall and exponential right jumps.

3 Consequences

Thanks to the previous results we can deduce some properties of the interacting particles from the known properties concerning the matrices $M(n)$, $n \geq 0$. For instance, the next proposition follows immediately from theorem 2.1 and theorem 5.2 of [3].

Proposition 3.1. *For $n \geq \lceil \frac{k+1}{2} \rceil$, the distribution function of $X_k(n)$ is given by*

$$\mathbb{P}(X_k(n) \leq t) = \det \left(\int_0^t x^{2j+i+n-\lceil \frac{k+1}{2} \rceil - 3 + 1_{\{k \text{ is even}\}}} e^{-x} dx \right)_{\lceil \frac{k+1}{2} \rceil \times \lceil \frac{k+1}{2} \rceil}, \quad t \in \mathbb{R}_+.$$

Moreover, theorem 2.1 implies proposition 2 of [1]. In fact, letting $X_0(m) = 0$ for every $m \in \mathbb{N}$, identity (1) implies that for $n \in \mathbb{N}$ and $i = 1, \dots, k$,

$$X_i(n) = \max_{0 \leq m \leq n} (X'_{i-1}(m) + \sum_{j=m+1}^n \xi(i, j)), \quad (2)$$

where

$$\begin{aligned} X'_{i-1}(0) &= 0, \\ X'_{i-1}(m) &= \max(X_{i-1}(m), X_{i-1}(m-1)) + \xi^+(i, m), \quad m \geq 1, \\ \xi(i, j) &= \xi^+(i, j) - \xi^-(i, j), \quad j \geq 1. \end{aligned}$$

Identity (2) proves that when n goes to infinity the process

$$\left(\left(\frac{1}{\sqrt{2n}} X_1(\lceil nt \rceil), \dots, \frac{1}{\sqrt{2n}} X_k(\lceil nt \rceil) \right), t \geq 0 \right)$$

converges in distribution to the process

$$((Y_1(t), \dots, Y_k(t)), t \geq 0)$$

defined by

$$\begin{aligned} Y_1(t) &= \sup_{0 \leq s \leq t} (B_1(t) - B_1(s)) \\ Y_j(t) &= \sup_{0 \leq s \leq t} (Y_{j-1}(s) + B_j(t) - B_j(s)), \quad j = 2, \dots, k, \end{aligned}$$

where B_1, \dots, B_k are independent real standard Brownian motion. Besides, Donsker's theorem implies that the process $(\frac{1}{\sqrt{2n}}M([nt]), t \geq 0)$ converges in distribution to a standard Brownian motion on $i\mathcal{A}_{k+1}$ when n goes to infinity. As announced proposition 2 of [1] follows then from theorem 2.1.

4 Proofs

Theorem 2.1 is proved by proposition 4.7. Theorem 2.2 follows from propositions 4.8 and 4.9. Let us denote by ϕ the function from \mathbb{R} to \mathbb{R} , defined by $\phi(x) = \frac{1}{2} \exp(-|x|)$, $x \in \mathbb{R}$, and consider the random walk $S_n, n \geq 1$, on \mathbb{R} , starting from 0, whose increments have a density equal to ϕ . The next three lemmas are elementary.

Lemma 4.1. *The process $(|S_n|)_{n \geq 1}$ is a Markov process, with a transition density q given by*

$$q(x, y) = \phi(x + y) + \phi(x - y), \quad x, y \in \mathbb{R}_+.$$

Proof. This is a simple computation which holds for any symmetric random walk. \square

Lemma 4.2. *Let us consider a process $M(n), n \geq 0$, on $i\mathcal{A}_2$, defined by*

$$M(n) = \sum_{l=1}^n Y_l \begin{pmatrix} 0 & i \\ -i & 0 \end{pmatrix} Y_l^*, \quad n \geq 0,$$

where the Y_l 's are independent standard Gaussian variables on $\mathcal{M}_{2,2}$. Then the process of the only positive eigenvalue of $(M(n), n \geq 0)$ is Markovian with transition density given by q .

Proof. The Y_l 's are 2×2 independent real random matrices whose entries are independent standard Gaussian random variables on \mathbb{R} . Let us write $Y_l = \begin{pmatrix} Y_{l,1} & Y_{l,2} \\ Y_{l,3} & Y_{l,4} \end{pmatrix}$. The matrix $Y_l \begin{pmatrix} 0 & i \\ -i & 0 \end{pmatrix} Y_l^*$ is equal to

$$\begin{pmatrix} 0 & i(Y_{l,4}Y_{l,1} - Y_{l,2}Y_{l,3}) \\ i(Y_{l,2}Y_{l,3} - Y_{l,1}Y_{l,4}) & 0 \end{pmatrix}.$$

For $\alpha \in \mathbb{R}$, we have

$$\mathbb{E}(e^{-i\alpha(Y_{l,4}Y_{l,1} - Y_{l,2}Y_{l,3})}) = \frac{1}{1 + \alpha^2}.$$

Thus, the random variables $Y_{l,4}Y_{l,1} - Y_{l,2}Y_{l,3}$, $l = 1, \dots, n$, are independent, with a density equal to ϕ . We conclude using lemma 4.1. \square

Lemma 4.3. *Let r be a real number. Let us consider $(\xi_n^-)_{n \geq 1}$ and $(\xi_n^+)_{n \geq 1}$ two independent families of independent random variables having an exponential law of mean 1. The Markov process $Z(n), n \geq 1$, defined by*

$$Z(n) = \max(Z(n-1) - \xi_n^-, r) + \xi_n^+, \quad n \geq 1,$$

has transition density

$$p_r(x, y) = \phi(x - y) + e^{2r} \phi(x + y), \quad x, y \geq r.$$

Proof. This is a simple computation. \square

Let us notice that when $r = 0$, the Markov process $(Z(n), n \geq 0)$, describes the evolution of the first particle. As its transition kernel p_0 is the same as the transition kernel q defined in lemma 4.1, theorem 2.1 follows when $k = 1$ from lemma 4.2. The general case is more complicated.

Definition 4.4. We define the function $d_k : \mathbb{R}^{\lfloor \frac{k+1}{2} \rfloor} \rightarrow \mathbb{R}$ by

- when $k = 2p$, $p \in \mathbb{N}^*$,

$$d_k(x) = c_k^{-1} \prod_{1 \leq i < j \leq p} (x_i^2 - x_j^2) \prod_{i=1}^p x_i,$$

- when $k = 2p - 1$, $p \in \mathbb{N}^*$,

$$d_k(x) = \begin{cases} 1 & \text{if } p = 1 \\ c_k^{-1} \prod_{1 \leq i < j \leq p} (x_i^2 - x_j^2) & \text{otherwise,} \end{cases}$$

where

$$c_k = 2^{\lfloor \frac{k}{2} \rfloor} \prod_{1 \leq i < j \leq p} (j - i)(k + 1 - j - i) \prod_{1 \leq i \leq p} (p + \frac{1}{2} - i)^{1_{\{k=2p\}}}.$$

The next proposition gives the transition density of the process of eigenvalues of the process $M(n)$, $n \geq 0$. For the computation, we need a generalized Cauchy-Binet identity (see for instance Johansson [7]). Let (E, \mathcal{B}, m) be a measure space, and let ϕ_i and ψ_j , $1 \leq i, j \leq n$, be measurable functions such that the $\phi_i \psi_j$'s are integrable. Then

$$\det \left(\int_E \phi_i(x) \psi_j(x) dm(x) \right) = \frac{1}{n!} \int_{E^n} \det(\phi_i(x_j)) \det(\psi_i(x_j)) \prod_{k=1}^n dm(x_k). \quad (3)$$

We also need an identity which expresses interlacing conditions with the help of a determinant (see Warren [11]). For $x, y \in \mathbb{R}^n$ we write $x \succeq y$ if x and y are interlaced, i.e.

$$x_1 \geq y_1 \geq x_2 \geq \cdots \geq x_n \geq y_n$$

and we write $x \succ y$ when

$$x_1 > y_1 > x_2 > \cdots > x_n > y_n.$$

When $x \in \mathbb{R}^{n+1}$ and $y \in \mathbb{R}^n$ we add the relation $y_n \geq x_{n+1}$ (resp. $y_n > x_{n+1}$). Let x and y be two vectors in \mathbb{R}^n such that $x_1 > \cdots > x_n$ and $y_1 > \cdots > y_n$. Then

$$1_{x \succ y} = \det(1_{\{x_i > y_j\}})_{n \times n}. \quad (4)$$

For $k \geq 1$, we denote by \mathcal{C}_k the subset of $\mathbb{R}^{\lfloor \frac{k+1}{2} \rfloor}$ defined by

$$\mathcal{C}_k = \{x \in \mathbb{R}^{\lfloor \frac{k+1}{2} \rfloor} : x_1 > \cdots > x_{\lfloor \frac{k+1}{2} \rfloor} > 0\}.$$

Proposition 4.5. Let us consider the process $(M(n), n \geq 0)$, defined as in theorem 2.1. For $n \in \mathbb{N}$, let $\Lambda(n)$ be the first $\lfloor \frac{k+1}{2} \rfloor$ largest eigenvalues of $M(n)$, ordered such that

$$\Lambda_1(n) \geq \cdots \geq \Lambda_{\lfloor \frac{k+1}{2} \rfloor}(n) \geq 0.$$

Then $(\Lambda(n), n \geq 0)$, is a Markov process with a transition density P_k given by

$$P_k(\lambda, \beta) = \frac{d_k(\beta)}{d_k(\lambda)} \det(\phi(\lambda_i - \beta_j) + (-1)^{k+1} \phi(\lambda_i + \beta_j))_{1 \leq i, j \leq \lfloor \frac{k+1}{2} \rfloor}, \quad \lambda, \beta \in \mathcal{C}_k.$$

Proof. The Markov property follows from the fact that the matrices

$$Y_l \begin{pmatrix} 0 & i \\ -i & 0 \end{pmatrix} Y_l^*, \quad l \in \mathbb{N}^*,$$

are independent and have an invariant distribution for the action of the orthogonal group by conjugacy. Proposition 4.8 of [3] ensures that the transition density with respect to the Lebesgue measure of the positive eigenvalues of $M(n)$, $n \geq 1$, is given by

- when $k = 2p$,

$$P_k(\lambda, \beta) = \frac{d_k(\beta)}{d_k(\lambda)} I_p(\lambda, \beta), \quad \lambda, \beta \in \mathcal{C}_k,$$

- when $k = 2p - 1$,

$$P_k(\lambda, \beta) = \frac{d_k(\beta)}{d_k(\lambda)} \frac{1}{2} (e^{-|\lambda_p - \beta_p|} + e^{-(\lambda_p + \beta_p)}) I_{p-1}(\lambda, \beta), \quad \lambda, \beta \in \mathcal{C}_k,$$

where

$$I_p(\lambda, \beta) = \begin{cases} 1 & \text{if } p = 0 \\ \int_{\mathbb{R}_+^p} 1_{\{\lambda, \beta \succ z\}} e^{-\sum_{i=1}^p (\lambda_i + \beta_i - 2z_i)} dz & \text{otherwise.} \end{cases}$$

When k is even, using identity (4), we write $1_{\{\lambda, \beta \succ z\}} e^{-\sum_{i=1}^p (\lambda_i + \beta_i - 2z_i)}$ as

$$\det(1_{z_i < \lambda_j} e^{-(\lambda_j - z_i)})_{p \times p} \det(1_{z_i < \beta_j} e^{-(\beta_j - z_i)})_{p \times p},$$

and use the Cauchy-Binet identity to get the proposition.

When k is odd, we introduce the measure μ on \mathbb{R} , defined by $\mu = \delta_0 + m$, where δ_0 is the Dirac measure at 0 and m is the Lebesgue measure on \mathbb{R} . We have the identity

$$\phi(x - y) + \phi(x + y) = \int_{\mathbb{R}} 1_{[0, x \wedge y]} e^{-(x+y-2z)} d\mu(z).$$

Thus using the Cauchy-Binet identity with the measure μ we get that the determinant of the proposition is equal to

$$\frac{1}{p!} \int_{\mathbb{R}^p} \det(1_{[0, \lambda_j]}(z_i) e^{-(\lambda_j - z_i)}) \det(1_{[0, \beta_j]}(z_i) e^{-(\beta_j - z_i)}) \prod_{m=1}^p d\mu(z_m).$$

Using identity (4), we obtain that it is equal to

$$\int_{\mathbb{R}^p} 1_{\lambda, \beta \succ z} e^{-\sum_{i=1}^p (\lambda_i + \beta_i - 2z_i)} \prod_{m=1}^p d\mu(z_m).$$

We integrate over z_p in the last integral and use the fact that the coordinates of β and λ are strictly positive to get the proposition. \square

The next proposition gives the transition density of the Markov process $(X(n), n \geq 0)$ defined in section 1. For $m \geq 0$, we denote $\phi^{(m)}$ the m^{th} derivative of ϕ , and we define the function $\phi^{(-m)}$, by

$$\phi^{(-m)}(x) = (-1)^m \int_x^{+\infty} \frac{1}{(m-1)!} (t-x)^{m-1} \phi(t) dt, \quad x \in \mathbb{R}.$$

We easily obtain that

$$\begin{aligned} \phi^{(m)}(x) &= \begin{cases} \frac{1}{2}(-1)^m e^{-x} & \text{if } x \geq 0 \\ \frac{1}{2}e^x & \text{otherwise.} \end{cases} \\ \phi^{(-m)}(x) &= \begin{cases} \phi^{(m)}(x) & \text{if } x \geq 0 \\ -\sum_{i=1}^{\lfloor \frac{m+1}{2} \rfloor} x^{m-(2i-1)} + \phi^{(m)}(x) & \text{otherwise.} \end{cases} \end{aligned}$$

For $k \geq 2$, we denote by \mathcal{D}_k the subset of \mathbb{R}^k defined by

$$\mathcal{D}_k = \{x \in \mathbb{R}^k : 0 < x_1 < x_2 < \dots < x_k\}.$$

Proposition 4.6. *The Markov process $X(n) = (X_1(n), \dots, X_k(n)), n \geq 0$, has a transition density Q_k given by*

$$Q_k(y, y') = \det(a_{i,j}(y_i, y'_j))_{1 \leq i, j \leq k}, \quad y, y' \in \bar{\mathcal{D}}_k,$$

where for $x, x' \in \mathbb{R}$,

$$a_{i,j}(x, x') = (-1)^{i-1} \phi^{(j-i)}(x+x') + (-1)^{i+j} \phi^{(j-i)}(x-x').$$

Proof. Let us show the proposition by induction on k . For $k = 1$, the equality holds by lemma 4.2. Suppose that it is true for $k - 1$. We write C_1, \dots, C_k , and L_1, \dots, L_k for the columns and the rows of the matrix of which we compute the determinant. There are two cases :

- If $y'_{k-1} \geq y_{k-1}$ then for $i = 1, \dots, k-1$, $y_i \leq y'_{k-1} \leq y'_k$. A quick calculation shows that all the components of the column $C_k + e^{y'_{k-1}-y'_k} C_{k-1}$ are equal to zero except the last one, which is equal to $p_{y'_{k-1}}(y'_{k-1} \vee y_k, y'_k)$, where p_r is defined in lemma 4.3 for $r \in \mathbb{R}$.
- If $y'_{k-1} \leq y_{k-1}$ then for $i = 1, \dots, k-1$, $y'_i \leq y_{k-1} \leq y_k$. We replace the last line L_k by the line $L_k - e^{y_{k-1}-y'_k} L_{k-1}$ having all its components equal to zero except the last one, which is equal to $p_{y_{k-1}}(y_k, y'_k)$.

Then we conclude developing the determinant according to its last column in the first case or to its last row in the second one, and using the induction property. \square

We have now all the ingredients needed to prove theorems 2.1 and 2.2. For λ an element of the closure $\bar{\mathcal{C}}_k$ of \mathcal{C}_k , we denote $\text{GT}_k(\lambda)$ the subset of $\mathbb{R}^{\lfloor \frac{k+1}{2} \rfloor \times \lfloor \frac{k+2}{2} \rfloor}$ defined by

$$\text{GT}_k(\lambda) = \{(x^{(2)}, \dots, x^{(k+1)}) : x^{(k+1)} = \lambda, x^{(i)} \in \mathbb{R}_+^{\lfloor \frac{i}{2} \rfloor}, x^{(i)} \succeq x^{(i-1)}, 3 \leq i \leq k+1\}.$$

We let

$$\text{GT}_k = \cup_{\lambda \in \bar{\mathcal{C}}_k} \text{GT}_k(\lambda).$$

GT_k is the set of Gelfand-Tsetlin patterns for the orthogonal group. If $(x^{(2)}, \dots, x^{(k+1)})$ is an element of GT_k , then $(x_1^{(2)}, \dots, x_1^{(k+1)})$ belongs to the closure $\bar{\mathcal{D}}_k$ of \mathcal{D}_k . Thus we define L_k as the Markov kernel on $\bar{\mathcal{C}}_k \times \bar{\mathcal{D}}_k$ such that for $\lambda \in \bar{\mathcal{C}}_k$, the probability measure $L_k(\lambda, \cdot)$ is the image of the uniform probability measure on $\text{GT}_k(\lambda)$ by the projection $p : \text{GT}_k \rightarrow \bar{\mathcal{D}}_k$ defined by

$$p((x^{(2)}, \dots, x^{(k+1)})) = (x_1^{(2)}, \dots, x_1^{(k+1)}),$$

where $(x^{(2)}, \dots, x^{(k+1)}) \in \text{GT}_k$. For $\lambda \in \bar{\mathcal{C}}_k$, the volume of $\text{GT}_k(\lambda)$ is given by $d_k(\lambda)$. Thus $L_k(\lambda, \cdot)$ has a density with respect to the Lebesgue measure on GT_k given by

$$\frac{1}{d_k(\lambda)} \mathbf{1}_{x \in \text{GT}_k(\lambda)}, \quad x \in \mathbb{R}^{\lfloor \frac{k+1}{2} \rfloor \lfloor \frac{k+2}{2} \rfloor}.$$

Rogers and Pitman proved in [9] (see also lemma 4 of [1]) that it is sufficient to show that the intertwining (5) holds, to get the equality in law of the processes $(\Lambda_1(n), n \geq 0)$ and $(X_k(n), n \geq 0)$. So theorem 2.1 follows from proposition 4.7.

Proposition 4.7.

$$L_k Q_k = P_k L_k \tag{5}$$

Proof. The proof is the same as the one of proposition 6 in [1]. We use the determinantal expressions for Q_k and P_k to show that both sides of equality (5) are equal to the same determinant. For this we use that the coefficients $a_{i,j}$'s given in proposition 4.6 satisfy for $x, x' \in \mathbb{R}_+$,

$$\begin{aligned} a_{i,j}(x, x') &= \int_x^{+\infty} a_{i-1,j}(u, x') du, \\ a_{i,j}(x, x') &= - \int_{x'}^{+\infty} a_{i,j+1}(u, x') du, \\ a_{2i,2j}(x, 0) &= 0, \quad a_{2i,2i-1}(0, x) = 1, \quad a_{2i,j}(0, x) = 0, \quad 2i \leq j \end{aligned}$$

The computation of the left hand side of (5) rests on the first identity. The computation of the right hand side rests on the others. \square

The measure $L_k(0, \cdot)$ is the Dirac measure at the null vector of $\text{GT}_k(0)$. Thus, the following proposition is an immediate consequence of proposition 4.7.

Proposition 4.8.

$$Q_k^n(0, \cdot) = P_k^n L_k(0, \cdot)$$

Keeping the same notations as in theorem 2.2, we have the following proposition, from which theorem 2.2 follows.

Proposition 4.9. $P_k^n L_k(0, \cdot)$ is the law of the random variable

$$(\Lambda_1^{(2)}(n), \dots, \Lambda_1^{(k+1)}(n)).$$

Proof. The density of the positive eigenvalues $\Lambda(n)$ of $M(n)$ is given by $P_k^n(0, \cdot)$. Then the proposition follows immediately from theorem 3.4 of [3]. \square

5 Concluding remarks

As recalled, the model with a wall that we have introduced is a variant of another one with no wall and no left-jumps. For this last model, the proofs of the analogue results as those of theorems 2.1 and 2.2 rest on some combinatorial properties of Young tableaux. Indeed, Young tableaux are used to describe the irreducible representations of the unitary group. The matrices from the Laguerre Unitary Ensemble belong to the set of Hermitian matrices which is, up to a multiplication by the complex i , the Lie algebra of the Unitary group. Their laws are invariant for the action of the unitary group by conjugacy. It is a general result that the law of their eigenvalues can be deduced from some combinatorial properties of the irreducible representations of the unitary group.

In our case, the distribution of the eigenvalues of $(M(n), n \geq 0)$ can be deduced from combinatorial properties of the irreducible representations of the orthogonal group (see [3] for details). Many combinatorial approaches have been developed to describe these representations. Among them we can mention the orthogonal tableaux and the analogue of the Robinson Schensted algorithm for the orthogonal group (see Sundaram [10]), or more recently those based on the very general theory of crystal graphs of Kashiwara [8]. None of them seems to lead to the interacting particles model with a wall that we have introduced. It would be interesting to find what kind of tableau involved in the description of irreducible representations of the orthogonal group would lead to this model.

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