

Bessel process, Schramm-Loewner evolution, and Dyson model

– Complex Analysis applied to
Stochastic Processes and Statistical Mechanics –

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Abstract

Bessel process is defined as the radial part of the Brownian motion (BM) in the D -dimensional space, and is considered as a one-parameter family of one-dimensional diffusion processes indexed by D , $\text{BES}^{(D)}$. First we give a brief review of $\text{BES}^{(D)}$, in which D is extended to be a continuous positive parameter. It is well-known that $D_c = 2$ is the critical dimension such that, when $D \geq D_c$ (resp. $D < D_c$), the process is transient (resp. recurrent). *Bessel flow* is a notion such that we regard $\text{BES}^{(D)}$ with a fixed D as a one-parameter family of initial value $x > 0$. There is another critical dimension $\bar{D}_c = 3/2$ and, in the intermediate values of D , $\bar{D}_c < D < D_c$, behavior of Bessel flow is highly nontrivial. The dimension $D = 3$ is special, since in addition to the aspect that $\text{BES}^{(3)}$ is a radial part of the three-dimensional BM, it has another aspect as a *conditional BM to stay positive*.

Two topics in probability theory and statistical mechanics, the *Schramm-Loewner evolution* (SLE) and the *Dyson model* (i.e., Dyson's BM model with parameter $\beta = 2$), are discussed. The $\text{SLE}^{(D)}$ is introduced as a 'complexification' of Bessel flow on the upper-half complex-plane, which is indexed by $D > 1$. It is explained that the existence of two critical dimensions D_c and \bar{D}_c for $\text{BES}^{(D)}$ makes $\text{SLE}^{(D)}$ have three phases; when $D \geq D_c$ the $\text{SLE}^{(D)}$ path is simple, when $\bar{D}_c < D < D_c$ it is self-intersecting but not dense, and when $1 < D \leq \bar{D}_c$ it is space-filling. The Dyson model is introduced as a multivariate extension of $\text{BES}^{(3)}$. By 'inheritance' from $\text{BES}^{(3)}$, the Dyson model has two aspects; (i) as an eigenvalue process of a Hermitian-matrix-valued BM, and (ii) as a system of BMs conditioned never to collide with each other, which we simply call the *noncolliding BM*. The noncolliding BM is constructed as a harmonic transform of absorbing BM in the Weyl chamber of type A, and as a complexification of this construction, the *complex BM representation* is proposed for the Dyson model. Determinantal expressions for spatio-temporal correlation functions with the *asymmetric correlation kernel of Eynard-Mehta type*

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are direct consequence of this representation. In summary, ‘parenthood’ of $\text{BES}^{(D)}$ and $\text{SLE}^{(D)}$, and that of $\text{BES}^{(3)}$ and the Dyson model are clarified.

Other related topics concerning extreme value distributions of noncolliding diffusion processes, statistics of characteristic polynomials of random matrices, and scaling limit of Fomin’s determinant for loop-erased random walks are also given.

We note that the name of Bessel process is due to the *special function* called the modified Bessel function, SLE is a stochastic time-evolution of *complex analytic function* (conformal transformation), and the Weierstrass canonical product representation of *entire functions* plays an important role for the Dyson model. Complex analysis is effectively applied to study stochastic processes of interacting particles and statistical mechanics models exhibiting critical phenomena and fractal structures in equilibrium and nonequilibrium states.

Keywords Complexification, Multivariate extension, Conformal transformation, Random matrices, Entire functions

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1 Family of Bessel processes

1.1 One-dimensional and D -dimensional Brownian motions

We consider motion of a Brownian particle in one dimensional space \mathbb{R} starting from $x \in \mathbb{R}$ at time $t = 0$. At each time $t > 0$ particle position is randomly distributed, and each realization of path is labeled by a parameter ω . Let Ω be the sample path space and $B^x(t, \omega)$ denote the position of the Brownian particle at time $t > 0$, whose path is realized as $\omega \in \Omega$. Let (Ω, \mathcal{F}, P) be a probability space. The *one-dimensional standard Brownian motion* (BM), $\{B^x(t, \omega)\}_{t \in [0, \infty)}$, $x \in \mathbb{R}$, has the following three properties.

1. $B^x(0, \omega) = x$ with probability one (abbr. w.p.1).
2. For any fixed $\omega \in \Omega$, $B^x(t, \omega)$ is a real continuous function of t . In other words, $B^x(t)$ has a continuous path.
3. For any sequence of times, $t_0 \equiv 0 < t_1 < \dots < t_M$, $M \in \mathbb{N} \equiv \{1, 2, 3, \dots\}$, the increments $\{B^x(t_i) - B^x(t_{i-1})\}_{i=1}^M$ are independent, and distribution of each increment is normal with mean $m = 0$ and variance $\sigma^2 = t_i - t_{i-1}$. It means that for any $1 \leq i \leq M$ and $\ell < r$,

$$P(B^x(t_i) - B^x(t_{i-1}) \in [\ell, r]) = \int_{\ell}^r p_{t_i - t_{i-1}}(\delta|0) d\delta,$$

where we define for $a, b \in \mathbb{R}$

$$p_t(b|a) = \begin{cases} \frac{1}{\sqrt{2\pi t}} e^{-(a-b)^2/2t}, & \text{for } t > 0, \\ \delta(a-b), & \text{for } t = 0. \end{cases} \quad (1.1)$$

If we write the conditional probability as $P(\cdot|C)$, where C denotes the condition, the third property given above implies that for any $0 \leq s \leq t$

$$P(B^x(t) \in A | B^x(s) = a) = \int_A p_{t-s}(b|a) db$$

holds $\forall A \subset \mathbb{R}, \forall a \in \mathbb{R}$. Then the probability that the BM is observed in a region $A_i \subset \mathbb{R}$ at time t_i for each $i = 1, 2, \dots, M$ is given by

$$P(B^x(t_i) \in A_i, i = 1, 2, \dots, M) = \int_{A_1} dx_1 \cdots \int_{A_M} dx_M \prod_{i=1}^M p_{t_i - t_{i-1}}(x_i | x_{i-1}), \quad (1.2)$$

where $x_0 \equiv x$. The formula (1.2) means that for any fixed $s \geq 0$, under the condition that $B^x(s)$ is given, $\{B^x(t) : t \leq s\}$ and $\{B^x(t) : t > s\}$ are independent. This independence of the events in the future and those in the past is called *Markov property*¹. The integral

¹A positive random variable τ is called *Markov time* if the event $\{\tau \leq u\}$ is determined by the behavior of the process until time u and independent of that after u . The Brownian motion satisfies the property obtained by changing any deterministic time $s > 0$ into any Markov time τ in the definition of Markov property given here. It is called a *strong Markov property*. A stochastic process which is strong Markov and has a continuous path almost surely is called a *diffusion process*.

kernel $p_t(y|x)$ is called the *transition probability density function* of the BM. As defined by (1.1), $p_t(y|x)$ is nothing but the probability density function of the normal distribution (the Gaussian distribution) with mean $m = x$ and variance $\sigma^2 = t$. It should be noted that $p_t(y|x) = p_t(x|y)$ and $u_t(x) \equiv p_t(y|x)$ is a unique solution of the *heat equation (diffusion equation)*

$$\frac{\partial}{\partial t} u_t(x) = \frac{1}{2} \frac{\partial^2}{\partial x^2} u_t(x), \quad x \in \mathbb{R}, \quad t \in [0, \infty) \quad (1.3)$$

with the initial condition $u_0(x) = \delta(x - y)$. Therefore, $p_t(y|x)$ is also called the *heat kernel*.

Let $D \in \mathbb{N}$ denote the spatial dimension. For $D \geq 2$, the D -dimensional BM in \mathbb{R}^D starting from the position $\mathbf{x} = (x_1, \dots, x_D) \in \mathbb{R}^D$ is defined by the following D -dimensional vector-valued BM,

$$\mathbf{B}^{\mathbf{x}}(t) = (B_1^{x_1}(t), B_2^{x_2}(t), \dots, B_D^{x_D}(t)), \quad t \in [0, \infty), \quad (1.4)$$

where $\{B_i^{x_i}(t)\}_{i=1}^D$ are independent one-dimensional standard BMs.

1.2 D -dimensional Bessel process

For $D \in \mathbb{N}$, the D -dimensional *Bessel process* is defined as the absolute value (*i.e.*, the radial coordinate) of the D -dimensional BM,

$$\begin{aligned} X^{\mathbf{x}}(t) &\equiv |\mathbf{B}^{\mathbf{x}}(t)| \\ &= \sqrt{B_1^{x_1}(t)^2 + \dots + B_D^{x_D}(t)^2}, \quad t \in [0, \infty), \end{aligned} \quad (1.5)$$

where the initial value is given by $X^{\mathbf{x}}(0) = x = |\mathbf{x}| = \sqrt{x_1^2 + \dots + x_D^2} \geq 0$. By definition $X^{\mathbf{x}}(t)$ is nonnegative, $X^{\mathbf{x}}(t) \in \mathbb{R}_+ \equiv \{x \in \mathbb{R} : x \geq 0\}$. See Fig. 1.

By this definition, $X^{\mathbf{x}}(t)$ is a functional of D -tuples of stochastic processes $\{B_i^{x_i}(t)\}_{i=1}^D$. In order to describe the statistics of a function of several random variables, we have to see ‘propagation of error’. For stochastic processes, by *Itô’s formula* we can readily obtain an equation for the stochastic process that is defined as a functional of several stochastic processes. In the present case, we have the following equation,

$$dX^{\mathbf{x}}(t) = dB(t) + \frac{D-1}{2} \frac{dt}{X^{\mathbf{x}}(t)}, \quad t \in [0, \infty), \quad x > 0. \quad (1.6)$$

The first term of the RHS, $dB(t)$, denotes the infinitesimal increment of a one-dimensional standard BM starting from the origin at time $t = 0$, $B(t) = B^0(t)$. It should be noted that $B(t)$ is a different BM from any $B_i^{x_i}(t)$, $1 \leq i \leq D$, which were used to define $X^{\mathbf{x}}(t)$ in Eq.(1.5). Here we assume $X^{\mathbf{x}}(t) > 0$. Then, if $D > 1$, for an infinitesimal increment of time $dt > 0$, the second term in the RHS of (1.6) is positive. It means that there is a drift to increase the value of $X^{\mathbf{x}}(t)$. This drift term is increasing in D and decreasing in $X^{\mathbf{x}}(t)$. Since as $X^{\mathbf{x}}(t) \searrow 0$, the drift term $\nearrow \infty$, it seems that a ‘repulsive force’ is acting to the D -dimensional BM, $\mathbf{B}^{\mathbf{x}}(t)$, $|\mathbf{x}| > 0$ to keep the distance from the origin positive,

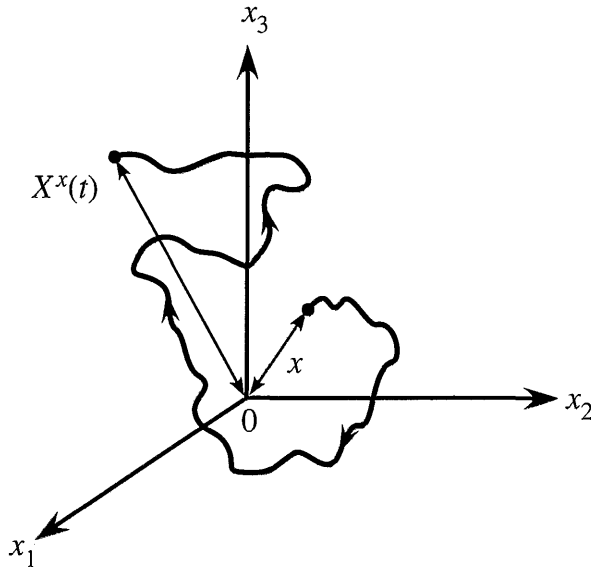


Figure 1: The D -dimensional Bessel process $X^x(t)$ is defined as the radial part of the BM in the D -dimensional space. The initial value x of the Bessel process is the distance between the origin and the position from which the BM is started.

$X^x(t) = |\mathbf{B}^{\mathbf{x}}(t)| > 0$, in order to avoid collision of the Brownian particle at the origin. Such a differential equation as (1.6), which involves random fluctuation term and drift term is called a *stochastic differential equation* (SDE).

What is the origin of the repulsive force between the D -dimensional BM and the origin? Why $\mathbf{B}^{\mathbf{x}}(t)$ starting from a point $\mathbf{x} \neq 0$ does not want to return to the origin? Why the strength of the outward drift is increasing in the dimension $D > 1$?

There is no positive reason for $\mathbf{B}^{\mathbf{x}}(t)$ to avoid visiting the origin, since by definition (1.4) all components $B_i^{x_i}(t)$ enjoy independent BMs. As the dimension of space D increases, however, the possibility *not* to visit the origin (or the fixed special point) increases, since among D directions in the space only one direction is toward the origin (or the fixed special point) and other $D - 1$ directions are orthogonal to it. If one know the second law of thermodynamics, which is also called the *law of increasing entropy*, one will understand that we would like to say here that the repulsive force acting from the origin to the Bessel process is an ‘entropy force’. (Note that the physical dimension of entropy [J/K] is different from that of force [J/m].) Anyway, the important fact is that, while the variance (quadratic variation) of the standard BM is fixed as $(dB(t))^2 = dt$ for a given $dt > 0$, the strength of repulsive drift is increasing in D . Then, the return probability of $X^x(t), x > 0$ to the origin should be a decreasing function of D .

Let $p_t^{(D)}(y|x)$ be the transition probability density of the D -dimensional Bessel process. We can show that, for any $y \in \mathbb{R}_+$, $u_t^{(D)}(x) \equiv p_t^{(D)}(y|x), x > 0$ solves the following partial differential equation (PDE)

$$\frac{\partial}{\partial t} u_t^{(D)}(x) = \frac{1}{2} \frac{\partial^2}{\partial x^2} u_t^{(D)}(x) + \frac{D-1}{2x} \frac{\partial}{\partial x} u_t^{(D)}(x) \quad (1.7)$$

under the initial condition $u_0^{(D)}(x) = \delta(x - y)$, which is called the *backward Kolmogorov equation* for the D -dimensional Bessel process. We can see clear correspondence between the SDE (1.6) and the PDE (1.7). As shown by (1.3), the BM term, $dB(t)$, in (1.6) is mapped to the diffusion term $(1/2)\partial^2 u_t^{(D)}(x)/\partial x^2$ in (1.7). In (1.7) the drift term is given by using the spatial derivative $\partial/\partial x$ representing the outward drift with the coefficient $(D-1)/2x$ corresponding to the factor $(D-1)/(2X^x(t))$ of the second term in (1.6). The solution is given by

$$p_t^{(D)}(y|x) = \begin{cases} \frac{1}{t} \frac{y^{\nu+1}}{x^\nu} e^{-(x^2+y^2)/2t} I_\nu\left(\frac{xy}{t}\right), & t > 0, x > 0, y \geq 0, \\ \frac{y^{2\nu+1}}{2^\nu t^{\nu+1} \Gamma(\nu+1)} e^{-y^2/2t}, & t > 0, x = 0, y \geq 0 \\ \delta(y-x), & t = 0, x, y \geq 0, \end{cases} \quad (1.8)$$

where $I_\nu(z)$ is the *modified Bessel function* of the first kind defined by

$$I_\nu(z) = \sum_{n=0}^{\infty} \frac{1}{\Gamma(n+1)\Gamma(n+1+\nu)} \left(\frac{z}{2}\right)^{2n+\nu} \quad (1.9)$$

with the gamma function $\Gamma(z) = \int_0^\infty e^{-u} u^{z-1} du$, and the index ν is specified by the dimension D as

$$\nu = \frac{D-2}{2} \iff D = 2(\nu+1). \quad (1.10)$$

This fact that $p_t^{(D)}(y|x)$ is expressed by using $I_\nu(z)$ gives the reason why the process $X^x(t)$ is called the Bessel process.

When $D = 3$, $\nu = 1/2$ by (1.10), and we can use the equality $I_{1/2}(z) = \sqrt{2/\pi z} \sinh z = (e^z - e^{-z})/\sqrt{2\pi z}$. Then (1.8) gives

$$p_t^{(3)}(y|x) = \frac{y}{x} \left\{ p_t(y|x) - p_t(y|-x) \right\} \quad (1.11)$$

for $t > 0, x > 0, y \geq 0$, where $p_t(y|x)$ is the transition probability density (1.1) of BM. If we put $p_t^{\text{abs}}(y|x) = p_t(y|x) - p_t(y|-x)$, we see $p_t^{\text{abs}}(0|x) = 0$ for any $x > 0$, since BM is a symmetric process.

As shown by Fig.2, $p_t^{\text{abs}}(y|x), x, y \geq 0$, gives the transition probability density of the absorbing BM, in which an absorbing wall is put at the origin and, if the Brownian particle starting from $x > 0$ arrives at the origin, it is absorbed there and the motion is stopped. By absorption, the total mass of paths from $x > 0$ to $y > 0$ is then reduced. The factor y/x appearing in (1.11) is for renormalization so that $\int_{\mathbb{R}_+} p_t^{(3)}(y|x) dy = 1, \forall t > 0, \forall x > 0$. We regard this renormalization procedure from p_t^{abs} to $p_t^{(3)}$ as a transformation. Since x is a one-dimensional harmonic function in a rather trivial sense $\Delta^{(1)}x \equiv d^2x/dx^2 = 0$, we say that the three-dimensional Bessel process is an *harmonic transform* (h -transform) of the one-dimensional absorbing BM in the sense of Doob [14]. This implies the equivalence between the three-dimensional Bessel process and ‘the one-dimensional BM conditioned to

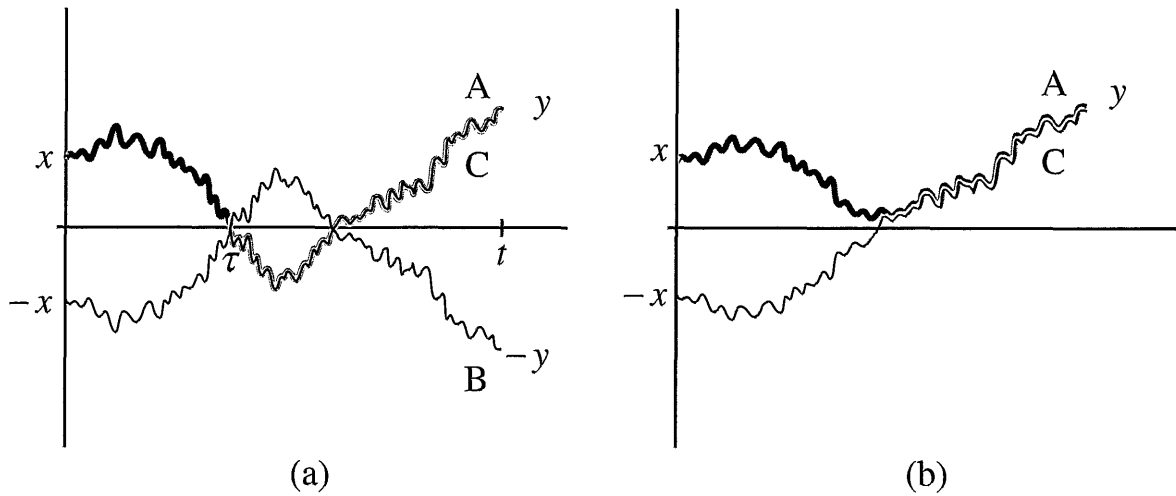


Figure 2: (a) One realization of Brownian path from $x > 0$ to $y > 0$ is drawn (path A), which visits the nonpositive region $\mathbb{R}_- \equiv \{x \in \mathbb{R} : x \leq 0\}$. The path B is a mirror image of path A with respect to the origin $x = 0$, which is running from $-x < 0$ to $-y < 0$. The first time when the path B hits the origin is denoted by τ . The path C is a combination of a part of path B up to time τ and a part of path A after τ such that it runs from $-x < 0$ to $y > 0$. There establishes bijection between path A and path C, which have the same weight as Brownian paths. Since the Brownian path A contributes to $p_t(y|x)$ and the Brownian path B does to $p_t(y|-x)$, such a path from $x > 0$ to $y > 0$ visiting \mathbb{R}_- is cancelled in $p_t^{\text{abs}}(y|x)$. (b) Each path, which does not visit \mathbb{R}_- , gives positive contribution to $p_t^{\text{abs}}(y|x)$, since in this case the weight of Brownian path A contributing to $p_t(y|x)$ is bigger than that of Brownian path C contributing $p_t(y|-x)$. In summary, $p_t^{\text{abs}}(y|x) = p_t(y|x) - p_t(y|-x)$ gives the total weight of Brownian paths, which do not hit the origin and thus are not absorbed at the wall at the origin.

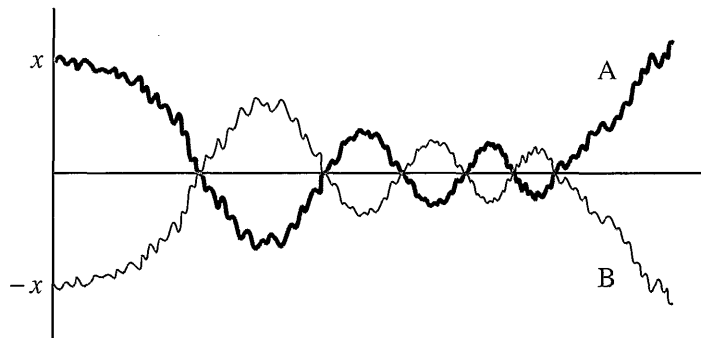


Figure 3: The one-dimensional BM visits the origin frequently. For a Brownian path starting from $x > 0$ (path A), its mirror image with respect to the origin is drawn (path B), which starts from $-x < 0$. If we observe the motion only in the nonnegative region \mathbb{R}_+ , the superposition of Brownian paths A and B gives a path of a reflecting BM, where a reflecting wall is put at the origin.

stay positive'. We will discuss such equivalence of processes in Section 2 more detail. Here we put emphasize the fact that $p_t^{(3)}(0|x) = 0, \forall x > 0$. It means that the three-dimensional Bessel process does not visit the origin. When $D = 3$, the outward drift is strong enough to avoid any visit to the origin. Moreover, we can prove that for any $x > 0$, $X^x(t) \rightarrow \infty$ as $t \rightarrow \infty$ w.p.1 and we say the process is *transient*.

When $D = 1$, $\nu = -1/2$ by (1.10) and we use the equality $I_{-1/2}(z) = \sqrt{2/\pi z} \cosh z = (e^z + e^{-z})/\sqrt{2\pi z}$. In this case (1.8) gives

$$p_t^{(1)}(y|x) = p_t(y|x) + p_t(y|-x) \quad (1.12)$$

for $t > 0, x, y \geq 0$. As shown by Fig.3, (1.12) means the equivalence between the one-dimensional Bessel process and 'the one-dimensional BM with a reflecting wall at the origin'. This is of course a direct consequence of the definition of Bessel process (1.5), since it gives $X^x(t) = |B^x(t)|$ in $D = 1$. The important fact is that the one-dimensional BM starting from $x \neq 0$ visits the origin frequently and we say that the one-dimensional Bessel process is *recurrent*. (Remark that in Eqs. (1.6) and (1.7), the drift terms vanish when $D = 1$. So we have to assume the reflecting boundary condition at the origin when we discuss the one-dimensional Bessel process instead of the one-dimensional BM.)

Now the following question is addressed: At which dimension the Bessel process changes its property from recurrent to transient ?

Before answering this question, here we would like to extend the setting of the question. Originally, the Bessel process was defined by (1.5) for $D \in \mathbb{N}$. We find that, however, the modified Bessel function (1.9) is an analytic function of ν for all values of ν . So we will be able to define the Bessel process for any positive value of dimension $D > 0$ as the diffusion process in \mathbb{R}_+ such that the transition probability density function is given by (1.8), where the index $\nu > -1$ is determined by (1.10) for each value of D . (In the SDE, (1.6), we assume the reflecting boundary condition at the origin for $0 < D < 2$.) Now we introduce an abbreviation $\text{BES}^{(D)}$ for the D -dimensional Bessel process, $D > 0$ ².

For $\text{BES}^{(D)}$ starting from $x > 0$, denote its first visiting time at the origin by

$$T^x = \inf\{t > 0 : X^x(t) = 0\}. \quad (1.13)$$

The answer of the above question is given by the following theorem.

Theorem 1.1 (i) $D \geq 2 \implies T^x = \infty, \forall x > 0, w.p.1.$

²Another characterization of $\text{BES}^{(D)}$ for fractional dimensions D is given by the following. Let $\nu \in \mathbb{R}$ and consider a BM with a constant drift ν , $B^y(t) + \nu t$, which starts from $y \in \mathbb{R}$ at time $t = 0$. The *geometric BM* with drift ν is defined as $\exp(B^y(t) + \nu t), t \geq 0$. For each $t \geq 0$, if we define random time change $t \mapsto A_t$ by

$$A_t = \int_0^t \exp\{2(B_s + \nu s)\} ds,$$

then the following relation is established,

$$X^x(A_t) = \exp(B_t^x + \nu t), \quad t \geq 0,$$

where $X^x(A_t)$ is the $\text{BES}^{(D)}$ with $D = 2(\nu + 1)$ at time A_t starting from $x = e^y$. The above formula is called *Lamperti's relation* [59, 88]

(ii) $D > 2 \implies \lim_{t \rightarrow \infty} X^x(t) = \infty, \forall x > 0, \text{ w.p.1, i.e. the process is transient.}$

(iii) $D = 2 \implies \inf_{t > 0} X^x(t) = 0, \forall x > 0, \text{ w.p.1.}$

That is, BES⁽²⁾ starting from $x > 0$ does not visit the origin, but it can visit any neighbor of the origin.

(iv) $D < 2 \implies T^x < \infty, \forall x > 0, \text{ w.p.1, i.e. the process is recurrent.}$

1.3 Bessel flow and Cardy's formula

In the previous subsection, we have defined the BES^(D) for positive continuous values of dimension $D > 0$ and studied dependence of the probability law of process on D . Theorem 1.1 states that the two-dimension is a *critical dimension*,

$$D_c = 2,$$

for competition between the two effects acting the Bessel process, the ‘random force’ (the martingale term) and the ‘entropy force’ (the drift term) in (1.6) and (1.7): when $D > D_c$, the latter dominates the former and the process becomes transient, and when $D < D_c$, the former is relevant and recurrence to the origin of the process is realized frequently.

Here we show that there is another critical dimension,

$$\overline{D}_c = \frac{3}{2}.$$

In order to characterize the transition at \overline{D}_c , we have to investigate dependence of the behavior of $X^x(t)$ on initial value $x > 0$. We call the one-parameter family $\{X^x(t)\}_{x > 0}$ the *Bessel flow* for each fixed $D > 0$.

For $0 < x < y$, we trace the motions of two BES^(D)'s starting from x and y by solving (1.6) using the *common* BM, $B(t), t \geq 0$,

$$\begin{aligned} X^x(t) &= x + B(t) + \frac{D-1}{2} \int_0^t \frac{ds}{X^x(s)}, \\ X^y(t) &= y + B(t) + \frac{D-1}{2} \int_0^t \frac{ds}{X^x(s)}, \quad t \geq 0. \end{aligned}$$

By considering the coupling of the two processes, we can show that

$$\begin{aligned} x < y &\implies X^x(t) < X^y(t), \quad t < T^x \quad \text{w.p.1} \\ &\implies T^x \leq T^y \quad \text{w.p.1.} \end{aligned}$$

The interesting fact is that in the intermediate fractional dimensions, $\overline{D}_c < D < D_c$, it is possible to see the coincide $T^x = T^y$ even for $x < y$. See Fig.4.

Theorem 1.2 For $0 < x < y < \infty$,

(i) $D \leq 3/2 \implies T^x < T^y \quad \text{w.p.1.}$

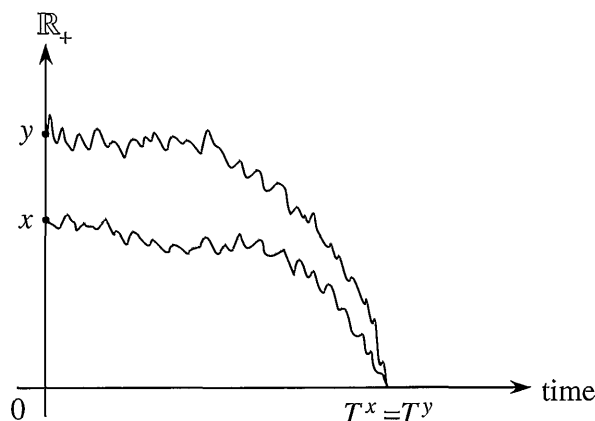


Figure 4: In the intermediate fractional dimensions, $3/2 < D < 2$, there is a positive probability that two Bessel processes starting from different initial positions, $0 < x < y < \infty$, return to the origin simultaneously, $T^x = T^y$.

$$(ii) \quad 3/2 < D < 2 \quad \implies \quad P(T^x = T^y) > 0.$$

Theorem 1.2 (ii) is obtained by proving that, for $0 < x < y < \infty, 3/2 < D < 2$, the event

$$\sup_{t < T^x} \frac{X^y(t) - X^x(t)}{X^x(t)} < \infty \quad (1.14)$$

occurs with positive probability. If (1.14) holds, $\exists c < \infty$ s.t. $X^y(t) \leq (1+c)X^x(t), 0 < t < T^x$ and thus $X^x(t) = 0 \implies X^y(t) = 0 \implies T^x = T^y$. We can confirm that the difference between the event $\{T^x = T^y\}$ and the event (1.14) has probability zero (see Section 1.10 in [60]).

A striking fact is the following exact formula: for $3/2 < D < 2, 0 < x < y < \infty$,

$$P(T^x = T^y) = 1 - \frac{\Gamma(D-1)}{\Gamma(2(D-1))\Gamma(2-D)} \left(\frac{y-x}{y}\right)^{2D-3} \times F\left(2D-3, D-1, 2(D-1); \frac{y-x}{y}\right), \quad (1.15)$$

where $F(\alpha, \beta, \gamma; z)$ is Gauss' hypergeometric function

$$F(\alpha, \beta, \gamma; z) = \sum_{i=0}^{\infty} \frac{(\alpha)_i (\beta)_i}{(\gamma)_i} \frac{z^i}{i!}$$

with the Pochhammer symbol $(c)_i = c(c+1)\cdots(c+i-1)$. If we set $D = 5/3$, (1.15) becomes a version of exact formula for a physical quantity called 'crossing probability' that Cardy derived in a critical percolation model [10, 11]. Cardy's formula has been extended in the context of SLE (see Section 6.7 of [60]), but I think that this exact formula for the Bessel flow can be also called *Cardy's formula*.

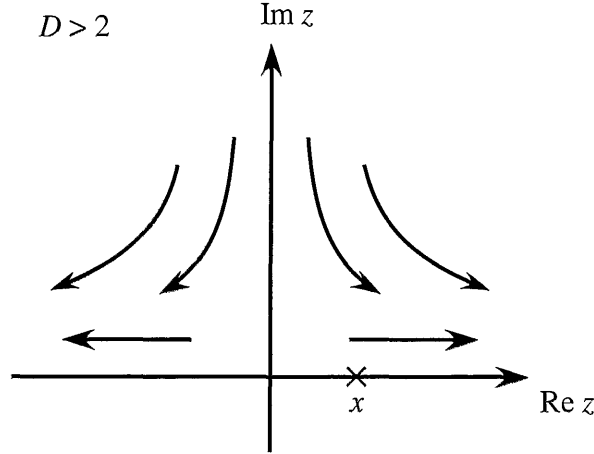


Figure 5: A schematic picture of ‘SLE flow’ on $\overline{\mathbb{H}} \setminus \{0\}$ for $D > 2$.

1.4 Schramm-Loewner evolution (SLE) as complexification of Bessel flow

Now we consider an extension of the Bessel flow $X^x(t)$ defined on \mathbb{R}_+ to flow on the upper-half complex-plane $\mathbb{H} = \{z = x + \sqrt{-1}y : x \in \mathbb{R}, y > 0\}$ and its boundary $\partial\mathbb{H} = \mathbb{R}$. We set $Z^z(t) = X^z(t) + \sqrt{-1}Y^z(t) \in \overline{\mathbb{H}} \setminus \{0\} = \mathbb{H} \cup \mathbb{R} \setminus \{0\}, t \geq 0$ and complexificate (1.6) as

$$dZ^z(t) = dB(t) + \frac{D-1}{2} \frac{dt}{Z^z(t)}, \quad t \in [0, \infty) \quad (1.16)$$

with the initial condition

$$Z^z(0) = z = x + \sqrt{-1}y \in \overline{\mathbb{H}} \setminus \{0\}.$$

The crucial point of this complexification of the Bessel flow is that the BM remains real. Then, there is asymmetry between the real part and the imaginary part of the flow in \mathbb{H} ,

$$dX^z(t) = dB(t) + \frac{D-1}{2} \frac{X^z(t)}{(X^z(t))^2 + (Y^z(t))^2} dt, \quad (1.17)$$

$$dY^z(t) = -\frac{D-1}{2} \frac{Y^z(t)}{(X^z(t))^2 + (Y^z(t))^2} dt. \quad (1.18)$$

Assume $D > 1$. Then as indicated by the minus sign in the RHS of (1.17), the flow is downward in $\overline{\mathbb{H}}$. If the flow goes down and arrives at the real axis, the imaginary part vanishes, $Y^z(t) = 0$, and Eq.(1.17) is reduced to be the same equation as Eq.(1.6) for the $\text{BES}^{(D)}$, which is now considered for $\mathbb{R} \setminus \{0\}$. If $D > 2$, by Theorem 1.1 (ii), the flow on $\mathbb{R} \setminus \{0\}$ is asymptotically outward, $X^x \rightarrow \pm\infty$ as $t \rightarrow \infty$. Therefore, the flow on $\overline{\mathbb{H}}$ will be described as shown by Fig.5. The behavior of flow should be, however, more complicated when $\overline{D}_c = 3/2 < D < D_c$ and $1 < D < \overline{D}_c$.

For $z \in \overline{\mathbb{H}} \setminus \{0\}$, $t \geq 0$, let ³

$$g_t(z) = Z^z(t) - B(t). \quad (1.19)$$

Then, Eq.(1.16) is written as follows:

$$\frac{d}{dt}g_t(z) = \frac{D-1}{2} \frac{1}{g_t(z) + B(t)}, \quad t \geq 0, \quad (1.20)$$

with the initial condition $g_0(z) = z \in \overline{\mathbb{H}} \setminus \{0\}$. For each $z \in \overline{\mathbb{H}} \setminus \{0\}$, set

$$T^z = \inf\{t > 0 : g_t(z) + B(t) = 0\}, \quad (1.21)$$

then the solution of Eq.(1.20) exists up to time T^z . For $t \geq 0$ we put

$$H_t = \{z \in \mathbb{H} : T^z > t\}. \quad (1.22)$$

This ordinary differential equation (1.20) involving the BM is nothing but the celebrated *Schramm-Loewner evolution* (SLE) [77, 60]. For each $t \geq 0$, the solution $g_t(z)$ of (1.20) gives a unique conformal transformation from H_t to \mathbb{H} :

$$g_t(z) : H_t \rightarrow \mathbb{H}, \quad \text{conformal}$$

such that

$$g_t(z) = z + \frac{a(t)}{z} + \mathcal{O}\left(\frac{1}{|z|^2}\right), \quad z \rightarrow \infty$$

with

$$a(t) = \frac{D-1}{2}t.$$

Note that the inverse map g_t^{-1} from \mathbb{H} to H_t , $t \geq 0$, is also conformal. For each $t \geq 0$, there exists a limit

$$\gamma(t) = \lim_{z \rightarrow 0, z \in \mathbb{H}} g_t^{-1}(z - B(t)), \quad (1.23)$$

and using basic properties of BM, we can prove that $\gamma = \gamma[0, \infty) \equiv \{\gamma(t) : t \in [0, \infty)\} \in \overline{\mathbb{H}}$ is a continuous path w.p.1 running from $\gamma(0) = 0$ to $\gamma(\infty) = \infty$ [74]. The path γ obtained from the SLE with the parameter $D > 1$ is called the *SLE^(D) path* ⁴.

The dependence on D of the Bessel flow given by Theorems 1.1 and 1.2 is mapped to the feature of the SLE^(D) paths such that there are three phases.

³Since $B^x(t)$, $X^x(t)$ and $Z^z(t)$ are stochastic processes, they are considered as functions of time $t \geq 0$, where the initial values x and z are put as superscripts ($B(t) \equiv B^0(t)$). On the other hand, as explained below, g_t is considered as a conformal transformation from a domain $H_t \subset \mathbb{H}$ to \mathbb{H} , and thus it is described as a function of $z \in H_t$; $g_t(z)$, where time t is a parameter and put as a subscript.

⁴Usual parameters used for the SLE are $\kappa = 4/(D-1)$ [77] or $a = (D-1)/2 = 2/\kappa$ [60]. If we set $\widehat{g}_t(z) = \sqrt{\kappa}g_t(z)$ and $B(t) = -W(t)$ in (1.20), we have the equation in the form [77],

$$\frac{d}{dt}\widehat{g}_t(z) = \frac{2}{\widehat{g}_t(z) - \sqrt{\kappa}W(t)}.$$

Note that $\sqrt{\kappa}W(t)$ has the same distribution with $B(\kappa t)$, which is a time change $t \mapsto \kappa t$ of the one-dimensional standard BM. In SLE, the Loewner chain for \mathbb{H} is driven by a one-dimensional BM, which is speeded up (or slowed down) by factor κ compared with the standard one. (The parameter κ is regarded as the diffusion constant.)

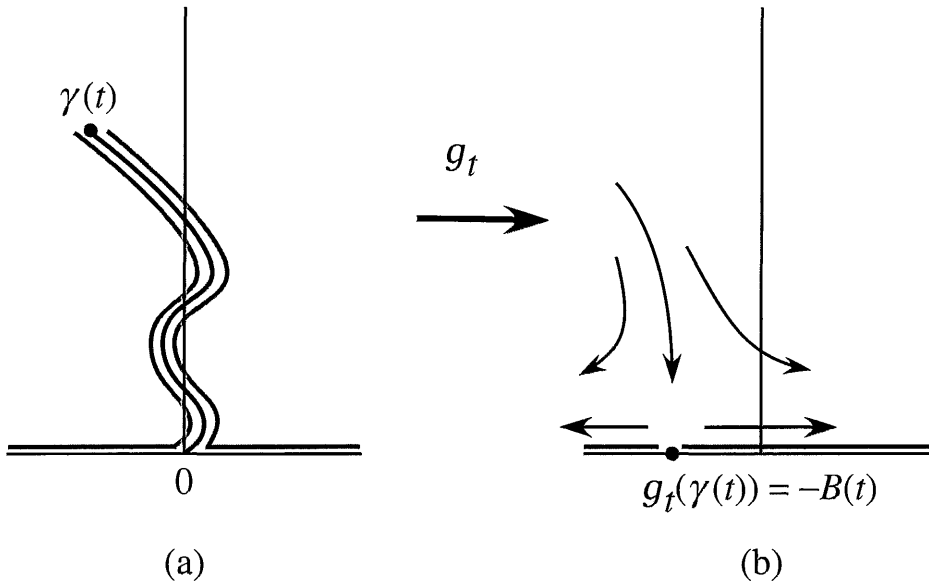


Figure 6: (a) When $D \geq 2$, the $\text{SLE}^{(D)}$ path is simple. (b) By g_t , the $\text{SLE}^{(D)}$ path is erased from \mathbb{H} . The tip of the $\text{SLE}^{(D)}$ path $\gamma(t)$ is mapped to $g_t(\gamma(t)) = -B(t) \in \mathbb{R}$. The flow associated by this conformal transformation is shown by arrows.

[phase 1] When $D \geq D_c = 2$, the $\text{SLE}^{(D)}$ path is a *simple curve*, i.e., $\gamma(s) \neq \gamma(t)$ for any $0 \leq s \neq t < \infty$, and $\gamma(0, \infty) \in \mathbb{H}$ (i.e., $\gamma(0, \infty) \cap \mathbb{R} = \emptyset$). In this phase,

$$H_t = \mathbb{H} \setminus \gamma(0, t], \quad t \geq 0.$$

For each $t \geq 0$, g_t gives a map, which conformally erases a simple curve $\gamma(0, t]$ from \mathbb{H} , and the image of the tip $\gamma(t)$ of the SLE path is $-B(t) \in \mathbb{R} = \partial\mathbb{H}$ as given by (1.23). As shown by Fig.6, it implies that the ‘SLE flow’ in $\overline{\mathbb{H}}$ is downward in the vertical (imaginary-axis) direction and outward from the position $-B(t)$ in the horizontal (real-axis) direction. Since $Z^z(t) = g_t(z) + B(t)$ by (1.19), if we shift this figure by $B(t)$, we will have the similar picture to Fig.5 for the complexified version of Bessel flow for $D > 2$.

[phase 2] When $\overline{D}_c = 3/2 < D < D_c = 2$, the $\text{SLE}^{(D)}$ path can osculate the real axis, $P(\gamma(0, t] \cap \mathbb{R} \neq \emptyset) > 0, \forall t > 0$. Fig.7 (a) illustrates the moment $t > 0$ such that the tip of $\text{SLE}^{(D)}$ path just osculates the real axis. The closed region encircled by the path $\gamma(0, t)$ and the line $[\gamma(t), 0] \in \mathbb{R}$ is called an *SLE hull* at time t and denoted by K_t . In this phase

$$H_t = \mathbb{H} \setminus K_t, \quad t \geq 0.$$

That is, $g_t(z)$ is a map which erases conformally the SLE hull from \mathbb{H} . We can think that by this transformation all the points in K_t are simultaneously mapped to a single point $-B(t) \in \mathbb{R}$, which is the image of the tip $\gamma(t)$. (We say that the hull K_t is *swallowed*. See Fig.7 (b).) By definition (1.22), the moment when $K_t = \mathbb{H} \setminus H_t$ is swallowed is the time T^z at which the equality $Z^z(t) = g_t(z) + B(t) = 0$ holds

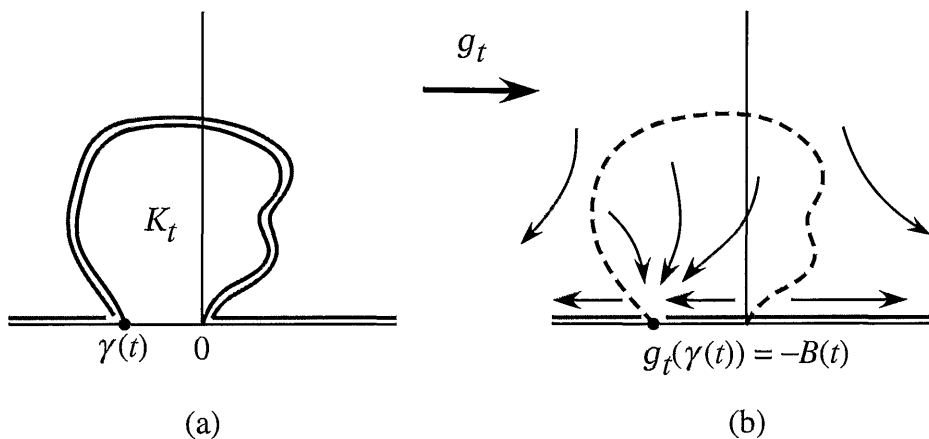


Figure 7: (a) When $3/2 < D < 2$, the $\text{SLE}^{(D)}$ path can osculate the real axis. The SLE hull is denoted by K_t . (b) The SLE hull K_t is swallowed. It means that all the points in K_t are simultaneously mapped to a single point $-B(t) \in \mathbb{R}$, which is the image of the tip of the $\text{SLE}^{(D)}$ path $\gamma(t)$.

$\forall z \in K_t$. (Then the RHS of (1.20) diverges and all the points $z \in K_t$ are lost from the domain of the map g_t .) Theorem 1.2 (ii) states that, when $\overline{D}_c < D < D_c$, two $\text{BES}^{(D)}$ starting from different points $0 < x < y < \infty$ can simultaneously return to the origin. In the complexified version, all $Z^z(t)$ starting from $z \in K_t$ can arrive at the origin simultaneously (*i.e.*, they are all swallowed).

Osculation of the SLE path with \mathbb{R} means that the SLE path has loops. Figure 8 (a) shows the event that the SLE path makes a loop at time $t > 0$. The SLE hull K_t consists of the closed region encircled by the loop and the segment of the SLE path between the origin and the osculating point, and it is completely erased by the conformal transformation g_t to \mathbb{H} as shown by Fig.8(b). Let $0 < s < t$ and consider the map g_s , which is the solution of (1.20) at time s . Assume that $\gamma(s)$ is located on the loop part of $\gamma[0, t]$ as shown by Fig.8(a). The segment $\gamma[0, s]$ of the SLE path is mapped by g_s to a part of \mathbb{R} . Since $\gamma(t)$ osculates a point in $\gamma[0, s]$, its image $g_s(\gamma(t))$ should osculate the real axis \mathbb{R} as shown by Fig.8(c). Since g_s^{-1} is uniquely determined from g_s , the above argument can be reversed. Then equivalence between osculation of the SLE path with \mathbb{R} and self-intersection of the SLE path is concluded. In this intermediate phase $\overline{D}_c < D < D_c$,

$$\text{SLE}^{(D)} \text{ path } \gamma \text{ is self-intersecting, and}$$

$$\bigcup_{t>0} \overline{K_t} = \overline{\mathbb{H}} \quad \text{but} \quad \gamma[0, \infty) \cap \mathbb{H} \neq \mathbb{H} \quad \text{w.p.1.}$$

[phase 3] When $1 < D \leq \overline{D}_c = 3/2$, Theorem 1.2 (i) states for the Bessel flow that the ordering $T^x < T^y$ is conserved for any $0 < x < y$. It implies that in this phase the SLE path should be a *space-filling curve*;

$$\gamma[0, \infty) = \overline{\mathbb{H}}.$$

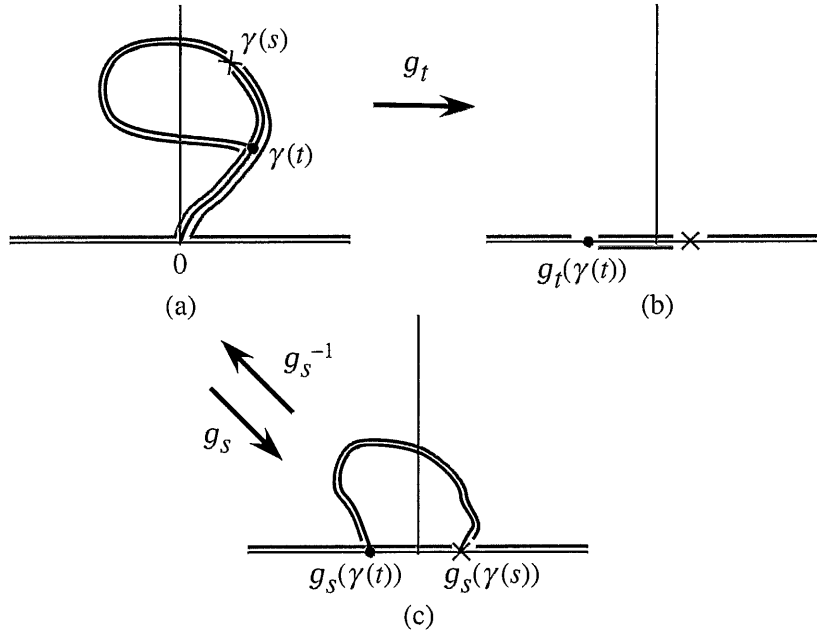


Figure 8: The event that the $SLE^{(D)}$ path osculates \mathbb{R} is equivalent with the event that the $SLE^{(D)}$ path makes a loop.

(Otherwise, swallow of regions occurs, contradicting Theorem 1.2 (i).)

Figure 9 summarizes the three phases of SLE paths.

By complexification of Bessel flow, we can discuss flows on a two-dimensional plane $\overline{\mathbb{H}}$. By this procedure random curves (the SLE paths) are generated in the plane. The SLE paths are fractal curves and their *Hausdorff dimensions* $d_H^{(D)}$ are determined by Beffara [1]. We note that a reciprocity relation is found between D and $d_H^{(D)}$;

$$(D - 1)(d_H^{(D)} - 1) = \frac{1}{2}, \quad D \geq \overline{D}_c = \frac{3}{2}. \quad (1.24)$$

(In the phase 3, $D \leq \overline{D}_c = 3/2$, $d_H^{(D)} \equiv 2$.)

Remark that the SLE map γ_t as well as the SLE path γ are functionals of the BM. Therefore, we have statistical ensembles of random curves $\{\gamma\}$ in the probability space (Ω, \mathcal{F}, P) . The important consequence from the facts that the BM is a strong Markov process with independent increments and g_t gives a conformal transformation is that the statistics of $\{\gamma\}$ has a kind of stationary Markov property (called the *domain Markov property*) and *conformal invariance* with respect to transformation of the domain in which the SLE path γ is defined [60].

The highlight of the theory of SLE would be that, if the value D is properly chosen, the statistics of $\{\gamma\}$ realizes that of the scaling limit of important *statistical mechanics model* exhibiting *critical phenomena* and fractal structures defined on an infinite discrete lattice. The following is a list of the correspondence between the $SLE^{(D)}$ paths with specified values of D and the scaling limits of models studied in statistical mechanics and

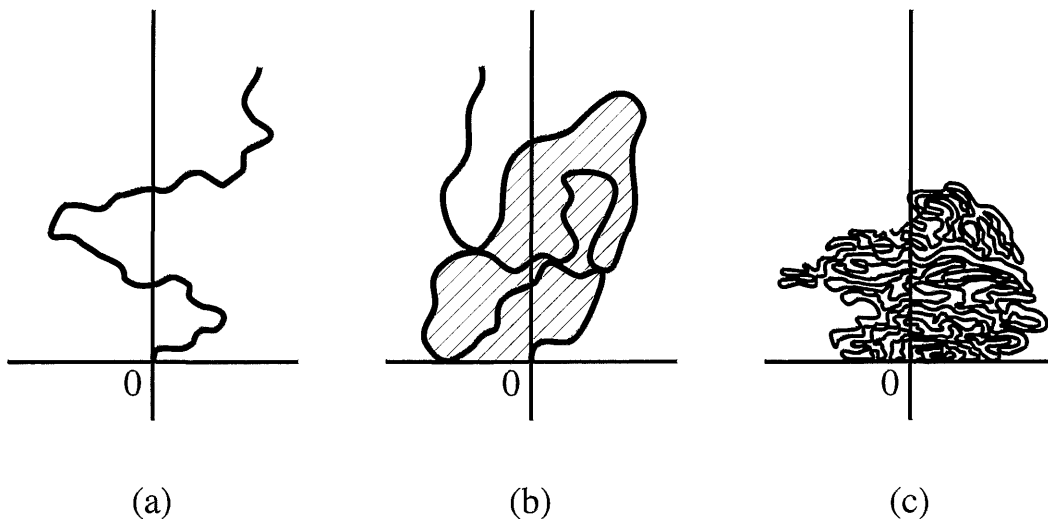


Figure 9: Schematic pictures of $\text{SLE}^{(D)}$ paths in (a) phase 1 ($D \geq D_c = 2$), (b) phase 2 ($\bar{D}_c = 3/2 < D < D_c = 2$), and (c) phase 3 ($1 < D \leq \bar{D}_c = 3/2$).

fractal physics ⁵.

$\text{SLE}^{(3/2)}$	\iff	uniform spanning tree [62]
$\text{SLE}^{(5/3)}$	\iff	critical percolation model (percolation exploration process) [80]
$\text{SLE}^{(2)}$	\iff	Gaussian free surface model (contour line) [78]
$\text{SLE}^{(7/3)}$	\iff	critical Ising model (Ising interface) [81]
$\text{SLE}^{(5/2)}$	\iff	self-avoiding walk [conjecture]
$\text{SLE}^{(3)}$	\iff	loop-erased random walk [62]

1.5 Dyson's BM model as multivariate extension of Bessel process

Here we consider stochastic motion of two particles $(X_1(t), X_2(t))$ in one dimension \mathbb{R} satisfying the following SDEs,

$$\begin{aligned}
 dX_1(t) &= dB_1(t) + \frac{\beta}{2} \frac{dt}{X_1(t) - X_2(t)} \\
 dX_2(t) &= dB_2(t) + \frac{\beta}{2} \frac{dt}{X_2(t) - X_1(t)} \quad t \geq 0,
 \end{aligned} \tag{1.25}$$

⁵ $\text{SLE}^{(D)}$ has a special property called the *restriction property* iff $D = 5/2$ ($\kappa = 8/3$). It is well-known that the self-avoiding walk (SAW) model, which has been studied as a model for polymers, has this property. The conformal invariance of the scaling limit of SAW is, however, not yet proved. If it is proved, the equivalence in probability law between the scaling limit of SAW and the $\text{SLE}^{(5/2)}$ path will be concluded.

with the initial condition $X_1(0) < X_2(0)$, where $B_1(t)$ and $B_2(t)$ are independent one-dimensional standard BMs and $\beta > 0$ is a ‘coupling constant’ of the two particles. The second terms in (1.25) represent the repulsive force acting between two particles, which is proportional to the inverse of distance $X_2(t) - X_1(t)$ of the two particles. Since it is a central force (*i.e.*, depending only on distance, and thus symmetric for two particles), the ‘center of mass’ $X_c(t) \equiv (X_2(t) + X_1(t))/2$ is proportional to a BM; $X_c(t) \stackrel{d}{=} B(t)/\sqrt{2} \stackrel{d}{=} B(t/2), \forall t \geq 0$, where $B(t)$ is a one-dimensional standard BM different from $B_1(t)$ and $B_2(t)$ and the symbol $\stackrel{d}{=}$ denotes the equivalence in distribution. (Note that variance (quadratic variation) $(dX_c(t))^2 = \{(dB_1(t) + dB_2(t))/2\}^2 = dt/2$, since $(dB_1(t))^2 = (dB_2(t))^2 = dt$ and $dB_1(t)dB_2(t) = 0$.) On the other hand, we can see that the relative coordinate defined by $X_r(t) \equiv (X_2(t) - X_1(t))/\sqrt{2}$ satisfies the SDE

$$dX_r(t) = d\tilde{B}(t) + \frac{\beta}{2} \frac{dt}{X_r(t)}, \quad t \geq 0,$$

where $\tilde{B}(t)$ is a BM different from $B_1(t), B_2(t), B(t)$. It is nothing but the SDE for $\text{BES}^{(D)}$ with $D = \beta + 1$.

Dyson [16] introduced N -particle systems of interacting BMs in \mathbb{R} as a solution $\mathbf{X}(t) = (X_1(t), X_2(t), \dots, X_N(t))$ of the following system of SDEs,

$$dX_i(t) = dB_i(t) + \frac{\beta}{2} \sum_{1 \leq j \leq N: j \neq i} \frac{dt}{X_i(t) - X_j(t)}, \quad t \in (0, \tau), \quad i = 1, 2, \dots, N, \quad (1.26)$$

where $\{B_i(t)\}_{i=1}^N$ are independent one-dimensional standard BMs and we define

$$\begin{aligned} \sigma_{ij} &= \inf\{t \geq 0 : X_i(t) \neq X_j(t)\}, \quad 1 \leq i < j \leq N, \\ \tau_{ij} &= \inf\{t > \sigma_{ij} : X_i(t) = X_j(t)\}, \quad 1 \leq i < j \leq N, \\ \tau &= \min_{1 \leq i < j \leq N} \tau_{ij}. \end{aligned}$$

It is called *Dyson’s BM model* with parameter β [65, 22]. As shown above, the $N = 2$ case of Dyson’s BM model is a coordinate transformation of the pair of a (time-change of) BM and $\text{BES}^{(\beta+1)}$. In this sense, Dyson’s BM model can be regarded as a multivariate (multi-dimensional) extension of $\text{BES}^{(\beta+1)}, \beta > 0$ ⁶. In particular, we will characterize Dyson’s BM model with $\beta = 2$ as an extension of the three-dimensional Bessel process, $\text{BES}^{(3)}$, in Section 2.

2 Two aspects of the Dyson model

In this section, we study the special case of Dyson’s BM model with parameter $\beta = 2$. We call this special case simply *the Dyson model* [48]. As shown above, the case $\beta = 2$ corresponds to the case $D = 3$ of Bessel process. In Sect.1.2, we have shown that $\text{BES}^{(3)}$ has two aspects; (**Aspect 1**) as a radial coordinate of three-dimensional BM, and (**Aspect 2**) as a one-dimensional BM conditioned to stay positive. We show that the Dyson model inherits these two aspects from $\text{BES}^{(3)}$ [47].

⁶We can prove that $\tau < \infty$ for $\beta < 1$ and $\tau = \infty$ for $\beta \geq 1$ [73]. The critical value $\beta_c = 1$ corresponds to $D_c = 2$ of $\text{BES}^{(D)}$.

2.1 The Dyson model as eigenvalue process

Dyson introduced the process (1.26) with $\beta = 1, 2$, and 4 as the eigenvalue processes of matrix-valued BMs in the Gaussian orthogonal ensemble (GOE), the Gaussian unitary ensemble (GUE), and the Gaussian symplectic ensemble (GSE) [16, 65, 22]⁷.

For $\beta = 2$ with given $N \in \mathbb{N}$, we prepare N -tuples of one-dimensional standard BMs $\{B_{ii}^{x_i}\}_{i=1}^N$, each of which starts from $x_i \in \mathbb{R}$, $N(N-1)/2$ -tuples of pairs of BMs $\{B_{ij}, \tilde{B}_{ij}\}_{1 \leq i < j \leq N}$, all of which start from the origin, where the totally $N+2 \times N(N-1)/2 = N^2$ BMs are independent from each other. Then consider an $N \times N$ Hermitian-matrix-valued BM⁸.

$$\mathcal{H}^{\mathbf{x}}(t) = \begin{pmatrix} B_{11}^{x_1}(t) & \frac{B_{12}(t) + \sqrt{-1}\tilde{B}_{12}(t)}{\sqrt{2}} & \cdots & \frac{B_{1N}(t) + \sqrt{-1}\tilde{B}_{1N}(t)}{\sqrt{2}} \\ \frac{B_{12}(t) - \sqrt{-1}\tilde{B}_{12}(t)}{\sqrt{2}} & B_{22}^{x_2}(t) & \cdots & \frac{B_{2N}(t) + \sqrt{-1}\tilde{B}_{2N}(t)}{\sqrt{2}} \\ \cdots & \cdots & \cdots & \cdots \\ \frac{B_{1N}(t) - \sqrt{-1}\tilde{B}_{1N}(t)}{\sqrt{2}} & \frac{B_{2N}(t) - \sqrt{-1}\tilde{B}_{2N}(t)}{\sqrt{2}} & \cdots & B_{NN}^{x_N}(t) \end{pmatrix}. \quad (2.1)$$

By this definition, the initial state of this BM is given by the diagonal matrix

$$\mathcal{H}^{\mathbf{x}}(0) = \text{diag}(x_1, x_2, \dots, x_N). \quad (2.2)$$

We assume $x_1 \leq x_2 \leq \cdots \leq x_N$.

Remember that when we introduced $\text{BES}^{(D)}$ in Sect.1.2, we considered the D -dimensional vector-valued BM, (1.4), by preparing D -tuples of independent one-dimensional standard BMs for its elements. Since the dimension of the space $\mathbb{H}(N)$ of $N \times N$ Hermitian matrices is $\dim \mathbb{H}(N) = N^2$, we need N^2 independent BMs for elements to describe a BM in this space $\mathbb{H}(N)$.

Corresponding to calculating an absolute value of $\mathbf{B}^{\mathbf{x}}(t)$, by which $\text{BES}^{(D)}$ was introduced as (1.5), here we calculate eigenvalues of $\mathcal{H}^{\mathbf{x}}(t)$. For any $t \geq 0$, there is a family of $N \times N$ unitary matrices $\{U(t)\}$ which diagonalize $\mathcal{H}^{\mathbf{x}}(t)$,

$$U(t)^* \mathcal{H}^{\mathbf{x}}(t) U(t) = \text{diag}(\lambda_1(t), \dots, \lambda_N(t)), \quad t \geq 0.$$

Let $\mathbb{W}_N^{\mathbf{A}}$ be the Weyl chamber of type A_{N-1} defined by

$$\mathbb{W}_N^{\mathbf{A}} \equiv \{\mathbf{x} = (x_1, x_2, \dots, x_N) \in \mathbb{R}^N : x_1 < x_2 < \cdots < x_N\}.$$

If we impose the condition $\boldsymbol{\lambda}(t) \equiv (\lambda_1(t), \dots, \lambda_N(t)) \in \mathbb{W}_N^{\mathbf{A}}, t \geq 0$, $U(t)$ is uniquely determined.

⁷Precisely speaking, Dyson considered the Ornstein-Uhlenbeck processes of eigenvalues such that as stationary states they have the eigenvalue distributions of random matrices in GOE, GUE, and GSE. Here we consider matrix-valued BMs, so variances increase in proportion to time $t \geq 0$.

⁸In usual Gaussian random matrix ensembles, mean is assumed to be zero. The corresponding matrix-valued BM are then considered to be started from a zero matrix, *i.e.*, $x_i = 0, 1 \leq i \leq N$ in (2.2). In random matrix theory, general case with non-zero means (*i.e.*, $x_i \neq 0$) is discussed with the terminology ‘random matrices in an external source’ (see, for example, [5]). From the view point of stochastic processes, imposing external sources to break symmetry of the system corresponds to changing initial state.

For each $t \geq 0$, $\lambda_i(t), 1 \leq i \leq N$ are functionals of $\{B_{ii}^{x_i}(t), B_{ij}(t), \tilde{B}_{ij}\}_{1 \leq i < j \leq N}$, and then, again we can apply Itô's formula to take into account propagation of error correctly to derive SDEs for the eigenvalue process $\boldsymbol{\lambda}(t), t \geq 0$ [7, 8, 41]. The result is the following,

$$d\lambda_i(t) = dB_i^{x_i}(t) + \sum_{1 \leq j \leq N, j \neq i} \frac{dt}{\lambda_i(t) - \lambda_j(t)}, \quad t \in (0, \infty), \quad i = 1, 2, \dots, N, \quad (2.3)$$

where $\{B_i^{x_i}\}_{i=1}^N$ are independent BM different from the BMs used to define $\mathcal{H}^{\boldsymbol{x}}(t)$ by (2.1). It is indeed the $\beta = 2$ case of (1.26) as derived by Dyson (originally not by such a stochastic calculus but by applying the perturbation theory in quantum mechanics) [16].

Now the correspondence is summarized as follows.

[Aspect 1]

BES⁽³⁾ \iff radial coordinate of
 $D = 3$ vector-valued BM

the Dyson model \iff eigenvalue process of
with N particles $N \times N$ Hermitian-matrix-valued BM

2.2 The Dyson model as noncolliding BM

Here we try to extend the formula (1.11) for BES⁽³⁾ to multivariate versions.

First we consider a set of two operations, identity ($\sigma_1 = \text{id}$) and reflection ($\sigma_2 = \text{ref}$), such that for $x \in \mathbb{R}$, $\sigma_1(x) = x$ and $\sigma_2(x) = -x$, and signatures are given as $\text{sgn}(\sigma_1) = 1$ and $\text{sgn}(\sigma_2) = -1$, respectively. Then we have

$$p_t^{\text{abs}}(y|x) = \sum_{\sigma \in \{\text{id}, \text{ref}\}} \text{sgn}(\sigma) p_t(y|\sigma(x)), \quad t \geq 0, \quad x, y, \in \mathbb{R}. \quad (2.4)$$

Next we consider a set of all permutations of N indices $\{1, 2, \dots, N\}$, which is denoted by \mathcal{S}_N , and put the following multivariate function

$$\sum_{\sigma \in \mathcal{S}_N} \text{sgn}(\sigma) \prod_{i=1}^N p_t(y_{\sigma(i)}|x_i) = \det_{1 \leq i, j \leq N} [p_t(y_i|x_j)] \quad (2.5)$$

of $\boldsymbol{x} = (x_1, \dots, x_N) \in \mathbb{W}_N^A$ and $\boldsymbol{y} = (y_1, \dots, y_N) \in \mathbb{W}_N^A$ with a parameter $t \geq 0$. Following the argument by Karlin and McGregor [34], we can prove that this determinant gives the transition probability density with duration t from the state \boldsymbol{x} to the state \boldsymbol{y} of N -dimensional absorbing BM, $\boldsymbol{B}^{\boldsymbol{x}}(t) = (B_1^{x_1}(t), \dots, B_N^{x_N}(t))$, in a domain \mathbb{W}_N^A , in which absorbing walls are put at the boundaries of \mathbb{W}_N^A . Since the boundaries of \mathbb{W}_N^A are the hyperplanes $x_i = x_j, 1 \leq i < j \leq N$, the Brownian particle $\boldsymbol{B}^{\boldsymbol{x}}(t)$ is annihilated, when any coincidence of the values of coordinates of $\boldsymbol{B}^{\boldsymbol{x}}(t)$ occurs. The 'survival probability' that the BM is not yet absorbed at the boundary is a monotonically decreasing function of time. If we regard the i -th coordinate $B_i^{x_i}(t)$ as the position of i -th particle on $\mathbb{R}, 1 \leq i \leq N$, the state $\boldsymbol{x} \in \mathbb{W}_N^A$ is considered to represent a configuration of N particles on \mathbb{R} such that a strict ordering $x_1 < x_2 < \dots < x_N$ of positions is maintained, while the state absorbed at

the boundary, $\mathbf{x} \in \partial\mathbb{W}_M^A$, is a configuration in which collision occurs between some pair of neighboring pairs of particles; $1 \leq i \leq N-1$, s.t. $x_i = x_{i+1}$. Since if such collision occurs, the process is totally annihilated, this many-particle system is called *vicious walker model* [20, 57, 32, 38, 33, 12].

As already noted below Eq.(1.11), $h_1(x) = x$ is a harmonic function in $(0, \infty)$ conditioned $\phi(0) = 0$. Similarly, if we consider a harmonic function of N variables

$$\Delta^{(N)}h_N(\mathbf{x}) \equiv \sum_{i=1}^N \frac{\partial^2}{\partial x_i^2} h_N(\mathbf{x}) = 0 \quad (2.6)$$

conditioned $h_N(\widehat{\mathbf{x}}) = 0$, $\forall \widehat{\mathbf{x}} \in \partial\mathbb{W}_N^A$, we will have the following *product of differences*

$$h_N(\mathbf{x}) = \prod_{1 \leq i < j \leq N} (x_j - x_i), \quad (2.7)$$

which is identified with the *Vandermonde determinant* $\det_{1 \leq i, j \leq N} [x_i^{j-1}]$.

Combining above consideration, we put the following function

$$p_t^N(\mathbf{y}|\mathbf{x}) = \frac{h_N(\mathbf{y})}{h_N(\mathbf{x})} \det_{1 \leq i, j \leq N} [p_t(y_i|x_j)], \quad t \geq 0, \quad \mathbf{x}, \mathbf{y} \in \mathbb{W}_N^A. \quad (2.8)$$

We can show that the factor $h_N(\mathbf{y})/h_N(\mathbf{x})$ provides an exact renormalization of the vicious walker model (the absorbing BM in \mathbb{W}_N^A) to compensate any decay of total mass of the process by collision (by absorption at $\partial\mathbb{W}_N^A$), and that (2.8) gives the transition probability density function for the N -particle system of one-dimensional BMs *conditioned never to collide with each other forever* (which we simply call the *noncolliding BM*) [25, 38, 39].

Moreover, by using the harmonicity (2.6), we can confirm that (2.8) satisfies the following partial differential equation,

$$\frac{\partial}{\partial t} p_t^N(\mathbf{y}|\mathbf{x}) = \frac{1}{2} \Delta^{(N)} p_t^N(\mathbf{y}|\mathbf{x}) + \sum_{1 \leq i \neq j \leq N} \frac{1}{x_i - x_j} \frac{\partial}{\partial x_i} p_t^N(\mathbf{y}|\mathbf{x}) \quad (2.9)$$

with the initial condition $p_0^N(\mathbf{y}|\mathbf{x}) = \delta(\mathbf{y} - \mathbf{x}) = \prod_{i=1}^N \delta(y_i - x_i)$ [38]. It can be regarded as the backward Kolmogorov equation of the stochastic process with N particles, $\mathbf{X}(t) = (X_1(t), \dots, X_N(t))$, which solves the system of SDEs,

$$dX_i(t) = dB_i^{x_i}(t) + \sum_{1 \leq j \leq N: j \neq i} \frac{dt}{X_i(t) - X_j(t)}, \quad t \in (0, \infty), \quad i = 1, 2, \dots, N. \quad (2.10)$$

Eq.(2.10) is identified with the $\beta = 2$ case of (1.26). Then the equivalence between the Dyson model and the noncolliding BM is proved.

The result is summarized as follows.

[Aspect 2]

$$\begin{aligned} \text{BES}^{(3)} &\iff h\text{-transform of absorbing BM in } (0, \infty) \\ &\iff \text{BM conditioned to stay positive} \\ \text{the Dyson model} &\iff h\text{-transform of absorbing BM in } \mathbb{W}_N^A \\ &\iff \text{noncolliding BM} \end{aligned}$$

The fact that BES⁽³⁾ and the Dyson model have two aspects implies useful relation between projection from higher dimensional spaces and restriction by imposing conditions. For matrix-valued processes, projection is performed by integration over irrelevant components, and noncolliding conditions are generally expressed by the Karlin-McGregor determinants. The processes discussed here are temporally homogeneous ones, but we can also discuss two aspects of temporally inhomogeneous processes. Actually, we have shown that, from the fact that the temporally inhomogeneous version of noncolliding BM has the two aspects, the Harish-Chandra-Itzykson-Zuber integral formula [26, 30] is derived,

$$\int_{\mathbf{U}(N)} dU \exp \left\{ -\frac{1}{2\sigma^2} \text{Tr}(\Lambda \mathbf{x} - U^* \Lambda \mathbf{y} U)^2 \right\} = \frac{C_N \sigma^{N^2}}{h_N(\mathbf{x}) h_N(\mathbf{y})} \det_{1 \leq i, j \leq N} [p_t(y_i | x_j)], \quad (2.11)$$

$\sigma^2 > 0$, where dU denotes the Haar measure of $\mathbf{U}(N)$ normalized as $\int_{\mathbf{U}(N)} dU = 1$, $\Lambda \mathbf{x} = \text{diag}(x_1, \dots, x_N)$, $\Lambda \mathbf{y} = \text{diag}(y_1, \dots, y_N)$ with $\mathbf{x}, \mathbf{y} \in \mathbb{W}_N^A$, and $C_N = (2\pi)^{N/2} \prod_{i=1}^N \Gamma(i)$ [40]⁹.

3 Determinantal processes and entire functions

3.1 Aspect 2 of the Dyson model

As Aspect 2, the Dyson model is constructed as the h -transform of the absorbing BM in \mathbb{W}_N^A . Therefore, at any positive time $t > 0$ the configuration is given as an element of \mathbb{W}_N^A ,

$$\mathbf{X}(t) = (X_1(t), X_2(t), \dots, X_N(t)) \in \mathbb{W}_N^A, \quad t > 0, \quad (3.1)$$

and there is no multiple point at which coincidence of particle positions, $X_i(t) = X_j(t)$, $i \neq j$, occurs. That is, the Dyson model is equivalent with the noncolliding BM. We can consider the Dyson model, however, starting from initial configurations with multiple points. In order to describe configurations with multiple points, we represent each particle configuration by a sum of delta measures in the form

$$\xi(\cdot) = \sum_{i \in \mathbb{I}} \delta_{x_i}(\cdot), \quad t \geq 0 \quad (3.2)$$

with a sequence of points in \mathbb{R} , $\mathbf{x} = (x_i)_{i \in \mathbb{I}}$, where \mathbb{I} is a countable index set. Here for $y \in \mathbb{R}$, $\delta_y(\cdot)$ denotes the delta measure such that $\delta_y(x) = 1$ for $x = y$ and $\delta_y(x) = 0$ otherwise. Then, for (3.2) and $A \subset \mathbb{R}$, $\xi(A) \equiv \int_A \xi(dx) = \sum_{i \in \mathbb{I}: x_i \in A} 1 = \#\{x_i : x_i \in A\}$. If the total number of particles N is finite, $\mathbb{I} = \{1, 2, \dots, N\}$, but we would like to also consider the cases with $N = \infty$. We call measures of the form (3.2) satisfying the condition $\xi(K) < \infty$ for any compact subset $K \subset \mathbb{R}$ nonnegative integer-valued Radon measures on \mathbb{R} and write the space of them as \mathfrak{M} . The set of configurations without multiple point is denoted by $\mathfrak{M}_0 = \{\xi \in \mathfrak{M} : \xi(\{x\}) \leq 1, \forall x \in \mathbb{R}\}$. There is a trivial correspondence between \mathbb{W}_N^A and \mathfrak{M}_0 .

⁹We can apply the present argument also for processes associated with Weyl chambers of other types. See [49, 41].

First we assume $\xi = \sum_{i \in \mathbb{I}} \delta_{x_i} \in \mathfrak{M}_0$, $\xi(\mathbb{R}) = N \in \mathbb{N}$ and consider the Dyson model as an \mathfrak{M}_0 -valued diffusion process,

$$\Xi(t, \cdot) = \sum_{i \in \mathbb{I}} \delta_{X_i(t)}(\cdot), \quad t \geq 0, \quad (3.3)$$

starting from the initial configuration ξ , where $\mathbf{X}(t) = (X_1(t), \dots, X_N(t))$ is the solution of (2.10) under the initial configuration $\mathbf{x} = (x_1, \dots, x_N) \in \mathbb{W}_N^A$. We write the process as $(\Xi(t), \mathbb{P}^\xi)$ and express the expectation with respect to the probability law \mathbb{P}^ξ of the Dyson model by $\mathbb{E}^\xi[\cdot]$. We introduce a filtration $\{\mathcal{F}(t)\}_{t \in [0, \infty)}$ on the space of continuous paths $C([0, \infty) \rightarrow \mathfrak{M})$ defined by $\mathcal{F}(t) = \sigma(\Xi(s), s \in [0, t])$, where σ denotes the sigma field.

Then we introduce a sequence of independent one-dimensional standard BMs, $\mathbf{B}^{\mathbf{x}}(t) = (B_i^{x_i}(t))_{i \in \mathbb{I}}, t \geq 0$ and write the expectation with respect to them as $\mathbb{E}^{\mathbf{x}}[\cdot]$.

Let $\mathbf{1}_{\mathbb{W}_N^A}(\mathbf{x}) = 1$ if $\mathbf{x} \in \mathbb{W}_N^A$ and $\mathbf{1}_{\mathbb{W}_N^A}(\mathbf{x}) = 0$ otherwise. Then Aspect 2 of the Dyson model is expressed by the following equality; for any $0 < t < T < \infty$, any symmetric function g on \mathbb{R}^N ,

$$\mathbb{E}^\xi[g(\mathbf{X}(t))] = \mathbb{E}^{\mathbf{x}} \left[g(\mathbf{B}(t)) \mathbf{1}_{\mathbb{W}_N^A}(\mathbf{B}(t)) \frac{|h_N(\mathbf{B}(t))|}{h_N(\mathbf{x})} \right], \quad (3.4)$$

where we have assumed the relations $\xi = \sum_{i \in \mathbb{I}} \delta_{x_i} \in \mathfrak{M}_0$, $\xi(\mathbb{R}) = N \in \mathbb{N}$, $\mathbf{x} = (x_1, \dots, x_N) \in \mathbb{W}_N^A$ and (3.3). Note that the indicator $\mathbf{1}_{\mathbb{W}_N^A}(\mathbf{B}(t))$ in the RHS annihilates the BM, $\mathbf{B}(t)$, if it hits any boundary of the Weyl chamber, $\partial \mathbb{W}_N^A$, and the factor $|h_N(\mathbf{B}(t))|/h_N(\mathbf{x})$ performs the h -transform of measure for Brownian paths. That is, the RHS of (3.4) indeed gives the expectation of g with respect to the process obtained by the h -transform of the absorbing BM in \mathbb{W}_N^A .

If we apply the Karlin-McGregor determinantal formula (see (2.5)),

$$(\text{RHS}) = \mathbb{E}^{\mathbf{x}} \left[\sum_{\sigma \in \mathcal{S}_N} \text{sgn}(\sigma) \mathbf{1}_{\mathbb{W}_N^A}(\sigma(\mathbf{B}(t))) g(\mathbf{B}(t)) \frac{|h_N(\mathbf{B}(t))|}{h_N(\mathbf{x})} \right],$$

where we recall the definition of determinant and let $\sigma(\mathbf{B}(t)) = (B_{\sigma(1)}, \dots, B_{\sigma(N)})$, $\sigma \in \mathcal{S}_N$. Since $h_N(\mathbf{x})$ is a product of differences, $\sum_{\sigma \in \mathcal{S}_N} \text{sgn}(\sigma) \mathbf{1}_{\mathbb{W}_N^A}(\sigma(\mathbf{B}(t))) |h_N(\mathbf{B}(t))| = \sum_{\sigma \in \mathcal{S}_N} \mathbf{1}_{\mathbb{W}_N^A}(\sigma(\mathbf{B}(t))) h_N(\mathbf{B}(t)) = h_N(\mathbf{B}(t))$, and then (3.4) is simply written as

$$\mathbb{E}^\xi[g(\mathbf{X}(t))] = \mathbb{E}^{\mathbf{x}} \left[g(\mathbf{B}(t)) \frac{h_N(\mathbf{B}(t))}{h_N(\mathbf{x})} \right]. \quad (3.5)$$

In the LHS of (3.5) we should note that the Dyson model is an interacting particle system such that between any pair of particles a long-range repulsive force acts. Strength of the two-body repulsive force is exactly proportional to the inverse of distance between two particles and thus it diverges as the distance goes to zero. By this strong repulsion, any collision of particles is prevented. On the other hand, in the RHS of (3.5), independent BMs are considered. When we calculate the expectation of a function g of them, however, we have to put extra weight $h_N(\mathbf{B}(t))/h_N(\mathbf{x})$ to their paths. Since if $|B_j(t) - B_i(t)| \rightarrow 0$ for any $i \neq j$, this weight becomes zero, again any collision of particle is prevented. An important point of the Karlin-McGregor formula is that this weight for paths is *signed*, *i.e.*, it can be positive and negative. Therefore, all particle configurations realized by intersections of paths in the 1+1 spatio-temporal plane are completely cancelled.

3.2 Complex BM representation

Now we consider complexification of the expression (3.5) [48]. For each $B_i^{x_i}(t), i \in \mathbb{I}$, we introduce an independent one-dimensional BM starting from the origin, $\tilde{B}_i(t)$, and define a complex BM as $Z_i(t) = B_i(t) + \sqrt{-1}\tilde{B}_i(t), i \in \mathbb{I}$. If we write the expectation with respect to $\{\tilde{B}_i(t)\}_{i \in \mathbb{I}}$ as $\tilde{\mathbb{E}}[\cdot]$ and define $\mathbf{E}^{\mathbf{x}} = \mathbb{E}^{\mathbf{x}} \otimes \tilde{\mathbb{E}}$, we can confirm that the RHS of (3.5) can be replaced by

$$\mathbf{E}^{\mathbf{x}} \left[g(\mathbf{B}(t)) \frac{h_N(\mathbf{Z}(t))}{h_N(\mathbf{x})} \right], \quad (3.6)$$

where $\mathbf{Z}(t) = (Z_i(t))_{i \in \mathbb{I}}$.

A key lemma of our theory [48] is the following identity; for $\mathbf{x} = (x_1, \dots, x_N) \in \mathbb{W}_N^{\mathbb{A}}$, $\mathbf{z} = (z_1, \dots, z_N) \in \mathbb{C}^N$,

$$\frac{h_N(\mathbf{z})}{h_N(\mathbf{x})} = \det_{1 \leq i, j \leq N} \left[\Phi_{\xi}^{x_i}(z_j) \right],$$

where

$$\Phi_{\xi}^u(z) = \prod_{x \in \text{supp } \xi \cap \{u\}^c} \left(1 - \frac{z - u}{x - u} \right) \quad (3.7)$$

for $\xi = \sum_{i \in \mathbb{I}} \delta_{x_i} \in \mathfrak{M}_0$ and $\text{supp } \xi \equiv \{x \in \mathbb{R} : \xi(\{x\}) > 0\}$. The function $\Phi_{\xi}^u(z)$ has an expression of the *Weierstrass canonical product with genus zero*. Then, it is an *entire function* with zeros at $\text{supp } \xi \cap \{u\}^c$ (see, for example, [64, 69]).

If we apply this identity to (3.6), we have quantities $\Phi_{\xi}^{x_i}(Z_j(t)), i, j \in \mathbb{I}$, which are conformal transforms of independent complex BMs, $Z_j(t), j \in \mathbb{I}$. Since complex BM is conformal invariant, each $\Phi_{\xi}^{x_i}(Z_j(t))$ is a time change of a complex BM, $Z_j(\cdot)$. Then the average is conserved,

$$\mathbf{E}^{\mathbf{x}}[\Phi_{\xi}^u(Z_j(t))] = \mathbf{E}^{\mathbf{x}}[\Phi_{\xi}^u(Z_j(T))], \quad 0 \leq t \leq T < \infty, \quad (3.8)$$

that is, $\{\Phi_{\xi}^u(Z_j(t))\}_{j \in \mathbb{I}}$ are independent *conformal local martingales* (see, for example, Section V.2 of [72]).

Let $0 < t < T < \infty$. Then for any $\mathcal{F}(t)$ -measurable function F in the continuous path space $C([0, \infty) \rightarrow \mathfrak{M})$, we have the equality

$$\mathbb{E}^{\xi}[F(\Xi(\cdot))] = \mathbf{E}^{\mathbf{x}} \left[F \left(\sum_{i \in \mathbb{I}} \delta_{B_i(\cdot)} \right) \det_{i, j \in \mathbb{I}} \left[\Phi_{\xi}^{x_i}(Z_j(T)) \right] \right]. \quad (3.9)$$

Now it is claimed that any observables of the Dyson model is calculated by a system of independent complex BMs, whose paths are weighted by a multivariate complex function $\det_{i, j \in \mathbb{I}}[\Phi_{\xi}^{x_i}(Z_j(T))]$, which is a conformal local martingale. We call (3.9) the ‘complex BM representation’ of the Dyson model [48].

3.3 Determinantal process with an infinite number of particles

For a configuration $\xi \in \mathfrak{M}$, we write the restriction of configuration in $A \subset \mathbb{R}$ as $(\xi \cap A)(\cdot) = \sum_{i \in \mathbb{I}: x_i \in A} \delta_{x_i}(\cdot)$, a shift of configuration by $u \in \mathbb{R}$ as $\tau_u \xi(\cdot) = \sum_{i \in \mathbb{I}} \delta_{x_i+u}(\cdot)$, and a square of configuration as $\xi^{(2)}(\cdot) = \sum_{i \in \mathbb{I}} \delta_{x_i^2}(\cdot)$, respectively. Let $C_0(\mathbb{R})$ be the set of all continuous real-valued functions with compact supports.

For any integer $M \in \mathbb{N}$, a sequence of times $\mathbf{t} = \{t_1, t_2, \dots, t_M\}$ with $0 < t_1 < \dots < t_M < T < \infty$, and a sequence of functions $\mathbf{f} = (f_{t_1}, f_{t_2}, \dots, f_{t_M}) \in C_0(\mathbb{R})^M$, the *moment generating function of multitime distribution* of $(\Xi(t), \mathbb{P}^\xi)$ is defined by

$$\Psi_{\mathbf{t}}^\xi[\mathbf{f}] \equiv \mathbb{E}_\xi \left[\exp \left\{ \sum_{m=1}^M \int_{\mathbb{R}} f_{t_m}(x) \Xi(t_m, dx) \right\} \right]. \quad (3.10)$$

From the fact that $\{\Phi_\xi^u(Z_j(t))\}_{j \in \mathbb{I}}$ are independent conformal local martingale with the property (3.8),

$$\begin{aligned} \mathbf{E}^{\mathbf{x}} \left[\Phi_\xi^{x_i}(Z_j(T)) \right] &= \mathbf{E}^{\mathbf{x}} \left[\Phi_\xi^{x_i}(Z_j(0)) \right] \\ &= \Phi_\xi^{x_i}(x_j) = \delta_{ij}, \quad \forall T \geq 0. \end{aligned}$$

Moreover, we can see that the following property holds for the determinantal weight in the complex BM representation (3.9). Let \mathbb{I}' be a subset of the index set \mathbb{I} and assume that a function f depends on $B_i(t), i \in \mathbb{I}'$, but does not on $B_j(t), j \in \mathbb{I} \setminus \mathbb{I}', 0 < t < T < \infty$. Then

$$\mathbf{E}^{\mathbf{x}} \left[f(\{B_i(t)\}_{i \in \mathbb{I}'} \det_{i,j \in \mathbb{I}'} [\Phi_\xi^{x_i}(Z_j(T))]) \right] = \mathbf{E}^{\mathbf{x}'} \left[f(\{B_i(t)\}_{i \in \mathbb{I}'} \det_{i,j \in \mathbb{I}'} [\Phi_\xi^{x_i}(Z_j(t))]) \right], \quad (3.11)$$

where $\mathbf{E}^{\mathbf{x}'}[\cdot]$ denotes the expectation with respect to $\{B_i^{x_i}(t)\}_{i \in \mathbb{I}'}$. Let $\xi \in \mathfrak{M}_0$ and

$$\begin{aligned} \mathbb{K}^\xi(s, x; t, y) &= \int_{\mathbb{R}} \xi(dv) p_s(x|v) \int_{\mathbb{R}} dw p_t(w|0) \Phi_\xi^v(y + \sqrt{-1}w) \\ &\quad - \mathbf{1}(s > t) p_{s-t}(x|y), \quad (s, t) \in (0, \infty)^2, \quad (x, y) \in \mathbb{R}^2, \end{aligned} \quad (3.12)$$

where $\mathbf{1}(\omega)$ is the indicator function of a condition ω ; $\mathbf{1}(\omega) = 1$ if ω is satisfied and $\mathbf{1}(\omega) = 0$ otherwise. Then, using this ‘reducibility’ (3.11), we have proved that (3.10) is given by a *Fredholm determinant*

$$\Psi_{\mathbf{t}}^\xi[\mathbf{f}] = \text{Det}_{\substack{(s,t) \in \mathbf{t}^2, \\ (x,y) \in \mathbb{R}^2}} \left[\delta_{st} \delta(x-y) + \mathbb{K}^\xi(s, x; t, y) \chi_t(y) \right], \quad (3.13)$$

where $\chi_{t_m}(\cdot) = e^{f_{t_m}(\cdot)} - 1, 1 \leq m \leq M$. We call \mathbb{K}^ξ a *correlation kernel*.

For $L > 0, \alpha > 0$ and $\xi \in \mathfrak{M}$ we put

$$M(\xi, L) = \int_{[-L, L] \setminus \{0\}} \frac{\xi(dx)}{x}, \quad M_\alpha(\xi, L) = \left(\int_{[-L, L] \setminus \{0\}} \frac{\xi(dx)}{|x|^\alpha} \right)^{1/\alpha},$$

and $M(\xi) = \lim_{L \rightarrow \infty} M(\xi, L)$, $M_\alpha(\xi) = \lim_{L \rightarrow \infty} M_\alpha(\xi, L)$, if the limits finitely exist. We have introduced the following conditions for initial configurations $\xi \in \mathfrak{M}$ [45]:

(C.1) there exists $C_0 > 0$ such that $|M(\xi, L)| < C_0$, $L > 0$,

(C.2) (i) there exist $\alpha \in (1, 2)$ and $C_1 > 0$ such that $M_\alpha(\xi) \leq C_1$,

(ii) there exist $\beta > 0$ and $C_2 > 0$ such that

$$M_1(\tau_{-a^2} \xi^{(2)}) \leq C_2 (\max\{|a|, 1\})^{-\beta} \quad \forall a \in \text{supp } \xi.$$

It was shown that, if $\xi \in \mathfrak{M}_0$ satisfies the conditions (C.1) and (C.2), then for $a \in \mathbb{R}$ and $z \in \mathbb{C}$, $\Phi_\xi^a(z) \equiv \lim_{L \rightarrow \infty} \Phi_{\xi \cap [a-L, a+L]}^a(z)$ finitely exists, and

$$|\Phi_\xi^a(z)| \leq C \exp \left\{ c(|a|^\theta + |z|^\theta) \right\} \left| \frac{z}{a} \right|^{\xi(\{0\})} \left| \frac{a}{a-z} \right|, \quad a \in \text{supp } \xi, \quad z \in \mathbb{C},$$

for some $c, C > 0$ and $\theta \in (\max\{\alpha, (2 - \beta)\}, 2)$, which are determined by the constants C_0, C_1, C_2 and the indices α, β in the conditions [45]. Then even if $\xi(\mathbb{R}) = \infty$, under the conditions (C.1) and (C.2), \mathbb{K}^ξ given by (3.12) is well-defined as a correlation kernel and dynamics of the Dyson model with an infinite number of particles $(\Xi(t), \mathbb{P}^\xi)$ exists [45]. We note that in the case that $\xi \in \mathfrak{M}_0$ satisfies the conditions (C.1) and (C.2) with constants C_0, C_1, C_2 and indices α and β , then $\xi \cap [-L, L], \forall L > 0$ does as well. Then we can obtain the convergence of moment generating functions $\Psi_{\mathbf{t}}^{\xi \cap [-L, L]}[\mathbf{f}] \rightarrow \Psi_{\mathbf{t}}^\xi[\mathbf{f}]$ as $L \rightarrow \infty$, which implies the convergence of the probability measures $\mathbb{P}^{\xi \cap [-L, L]} \rightarrow \mathbb{P}^\xi$ in $L \rightarrow \infty$ in the sense of finite dimensional distributions [45].

By definition of Fredholm determinant, the moment generating function (3.13) can be expanded with respect to $\chi_{t_m}(\cdot), 1 \leq m \leq M$, as

$$\Psi_{\mathbf{t}}^\xi[\mathbf{f}] = \sum_{\substack{N_m \geq 0, \\ 1 \leq m \leq M}} \int_{\prod_{m=1}^M \mathbb{W}_{N_m}^A} \prod_{m=1}^M \left\{ d\mathbf{x}_{N_m}^{(m)} \prod_{i=1}^{N_m} \chi_{t_m}(x_i^{(m)}) \right\} \rho^\xi(t_1, \mathbf{x}_{N_1}^{(1)}; \dots; t_M, \mathbf{x}_{N_M}^{(M)}),$$

with

$$\rho^\xi(t_1, \mathbf{x}_{N_1}^{(1)}; \dots; t_M, \mathbf{x}_{N_M}^{(M)}) = \det_{\substack{1 \leq i \leq N_m, 1 \leq j \leq N_n, \\ 1 \leq m, n \leq M}} \left[\mathbb{K}^\xi(t_m, x_i^{(m)}; t_n, x_j^{(n)}) \right],$$

where $\mathbf{x}_{N_m}^{(m)}$ denotes $(x_1^{(m)}, \dots, x_{N_m}^{(m)})$ and $d\mathbf{x}_{N_m}^{(m)} = \prod_{i=1}^{N_m} dx_i^{(m)}$, $1 \leq m \leq M$. The functions ρ^ξ 's are called *multitime correlation functions*, and $\Psi_{\mathbf{t}}^\xi[\mathbf{f}]$ can be regarded as a generating function of them. In general, when the moment generating function for the multitime distribution is given by a Fredholm determinant, all the spatio-temporal correlation functions ρ^ξ are given by determinants of matrices, whose entries are special values of a continuous function \mathbb{K}^ξ , and then the process is said to be *determinantal* [68, 37, 43]. (Therefore, \mathbb{K}^ξ is called a correlation kernel.) The results by Eynard and Mehta reported in [17] for a multi-layer matrix model can be regarded as the theorem that the Dyson model is determinantal for the special initial configuration $\xi = N\delta_0$, *i.e.*, all particles are put at the origin, for any $N \in \mathbb{N}$. The correlation kernel is expressed in this case by using the

Hermite orthogonal polynomials [67]. The present author and H. Tanemura proved that, for any fixed initial configuration $\xi \in \mathfrak{M}$ with $\xi(\mathbb{R}) \in \mathbb{N}$, the Dyson model $(\Xi(t), \mathbb{P}^\xi)$ is determinantal, in which the correlation kernel is given by

$$\begin{aligned} \mathbb{K}^\xi(s, x; t, y) &= \frac{1}{2\pi\sqrt{-1}} \oint_{\Gamma(\xi)} dz p_s(x|z) \int_{\mathbb{R}} dw p_t(-\sqrt{-1}y|w) \\ &\times \frac{1}{\sqrt{-1}w - z} \prod_{x' \in \text{supp } \xi} \left(1 - \frac{\sqrt{-1}w - z}{x' - z} \right) - \mathbf{1}(s > t) p_{s-t}(x|y), \end{aligned} \quad (3.14)$$

where $\Gamma(\xi)$ is a closed contour on the complex plane \mathbb{C} encircling the points in $\text{supp } \xi$ on the real line \mathbb{R} once in the positive direction [45]. When $\xi \in \mathfrak{M}_0$, (3.14) becomes (3.12) by performing the Cauchy integrals.

We note that $\Xi_{\mathbf{t}} \equiv \sum_{t \in \mathbf{t}} \delta_t \otimes \Xi(t)$ is a *determinantal point process* (or *Fermion point process*) on the spatio-temporal field $\mathbf{t} \times \mathbb{R}$ with an operator \mathcal{K} given by $\mathcal{K}f(s, x) = \sum_{t \in \mathbf{t}} \int_{\mathbb{R}} dy \mathbb{K}(s, x; t, y) f(t, y)$ for $f(t, \cdot) \in C_0(\mathbb{R})$, $t \in \mathbf{t}$. When \mathcal{K} is symmetric, Soshnikov [82] and Shirai and Takahashi [79] gave sufficient conditions for \mathbb{K}^ξ to be a correlation kernel of a determinantal point process (see also [27]). Such conditions are not known for asymmetric cases. The present correlation kernels (3.12) and (3.14) are asymmetric cases, $\mathbb{K}^\xi(s, x; t, y) \neq \mathbb{K}^\xi(t, y; s, x)$ by the second terms $-\mathbf{1}(s > t) p_{s-t}(x|y)$. Such form of asymmetric correlation kernels is said to be of the Eynard-Mehta type [4, 48].

From the view point of statistical physics, such asymmetry is useful to describe *nonequilibrium systems* developing in time. In order to demonstrate it, we have studied the following *relaxation phenomenon* of the Dyson model with an infinite number of particles.

We consider the configuration in which every point of the integers \mathbb{Z} is occupied by one particle,

$$\xi^{\mathbb{Z}}(\cdot) = \sum_{i \in \mathbb{Z}} \delta_i(\cdot).$$

See Fig.10. It can be confirmed that $\xi^{\mathbb{Z}}$ satisfies the conditions (C.1) and (C.2) and thus the Dyson model starting from $\xi^{\mathbb{Z}}$, $(\Xi(t), \mathbb{P}^{\xi^{\mathbb{Z}}})$, is well-defined as a determinantal process with an infinite number of particles.

As a matter of fact, we have shown that the correlation kernel is given by [45]

$$\begin{aligned} \mathbb{K}^{\xi^{\mathbb{Z}}}(s, x; t, y) &= \mathbf{K}_{\sin}(s, x; t, y) \\ &+ \frac{1}{2\pi} \int_{|k| \leq \pi} dk e^{k^2(t-s)/2 + \sqrt{-1}k(y-x)} \left\{ \vartheta_3(x - \sqrt{-1}ks, 2\pi\sqrt{-1}s) - 1 \right\} \\ &= \mathbf{K}_{\sin}(s, x; t, y) \\ &+ \sum_{n \in \mathbb{Z} \setminus \{0\}} e^{2\pi\sqrt{-1}xn - 2\pi^2sn^2} \int_0^1 du e^{\pi^2u^2(t-s)/2} \cos \left[\pi u \{ (y-x) - 2\pi\sqrt{-1}sn \} \right], \end{aligned}$$

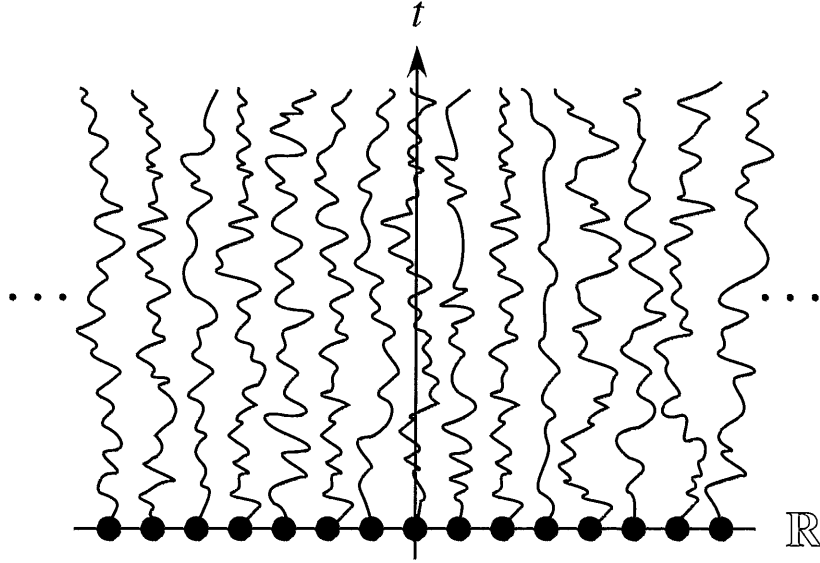


Figure 10: Consider the Dyson model starting from the configuration in which every point of the integers \mathbb{Z} is occupied by one particle. This nonequilibrium determinantal process shows a relaxation phenomenon to the stationary state μ_{sin} .

$(s, t) \in [0, \infty)^2, (x, y) \in \mathbb{R}^2$, where

$$\begin{aligned} \mathbf{K}_{\text{sin}}(s, x; t, y) &= \frac{1}{2\pi} \int_{|k| \leq \pi} dk e^{k^2(t-s)/2 + \sqrt{-1}k(y-x)} - \mathbf{1}(s > t) p_{s-t}(x|y) \\ &= \begin{cases} \int_0^1 du e^{\pi^2 u^2(t-s)/2} \cos\{\pi u(y-x)\} & \text{if } t > s \\ K_{\text{sin}}(x, y) & \text{if } t = s \\ - \int_1^\infty du e^{\pi^2 u^2(t-s)/2} \cos\{\pi u(y-x)\} & \text{if } t < s, \end{cases} \end{aligned}$$

with

$$K_{\text{sin}}(x, y) = \frac{1}{2\pi} \int_{|k| \leq \pi} dk e^{\sqrt{-1}k(y-x)} = \frac{\sin\{\pi(y-x)\}}{\pi(y-x)}, \quad x, y \in \mathbb{R}, \quad (3.15)$$

and ϑ_3 is a version of the *Jacobi theta function* defined by

$$\vartheta_3(v, \tau) = \sum_{n \in \mathbb{Z}} e^{2\pi\sqrt{-1}vn + \pi\sqrt{-1}\tau n^2}, \quad \Im\tau > 0. \quad (3.16)$$

By this explicit expression, we can see that

$$\lim_{u \rightarrow \infty} \mathbb{K}^{\xi^{\mathbb{Z}}}(u + s, x; u + t, y) = \mathbf{K}_{\text{sin}}(s, x; t, y). \quad (3.17)$$

The correlation kernel $\mathbf{K}_{\text{sin}}(s, x; t, y)$ surviving in the long-term limit is called the *extended sine kernel*. It is symmetric, $\mathbf{K}_{\text{sin}}(s, x; t, y) = \mathbf{K}_{\text{sin}}(t, y; s, x)$, and the determinantal process with the correlation kernel $\mathbf{K}_{\text{sin}}(s, x; t, y)$ is an *equilibrium dynamics*. It is time-reversal with respect to the determinantal point process μ_{sin} , in which any spatial correlation function is given by a determinant with the correlation kernel (3.15) called the *sine kernel*. This stationary state μ_{sin} is a scaling limit (called the bulk-scaling-limit) of the eigenvalue distribution of random matrices in GUE [65, 22]. See [83, 70, 67, 43].

The theory of entire functions discusses the relations between the growth of an entire function and the distribution of its zeros [64, 69]. Here we set distributions of zeros as initial configurations satisfying some conditions and control the behavior of particles at infinity to realize nonequilibrium dynamics of long-rang interacting infinite-particle systems. Systematic study of determinantal processes with infinite number of particles exhibiting relaxation phenomena is now in progress [44, 45, 46, 48]¹⁰.

4 Related Topics

At the end of this manuscript, I briefly introduce related topics, which we are interested in.

4.1 Extreme value distributions of noncolliding diffusion processes

As explained in Section 1.2, BES⁽³⁾ is the conditional BM to stay positive. When we impose an additional condition such that it starts from the origin at time $t = 0$ and return to the origin at time $t = 1$, the process is called the *three-dimensional Bessel bridge* with duration 1, which is here denoted by $Y(t), 0 \leq t \leq 1$.

Here we consider the maximum of $Y(t)$,

$$H_1 = \max_{0 < t < 1} Y(t).$$

See Fig.11 We can show that the distribution of H_1 is described as [52]

$$\mathbb{P}(H_1 \leq h) = -\frac{1}{2} \sum_{n \in \mathbb{Z}} H_2(\sqrt{2}nh) e^{-2h^2n^2}, \quad (4.1)$$

where $H_i(x)$ is the i -th Hermite polynomial

$$H_i(x) = i! \sum_{j=0}^{\lfloor i/2 \rfloor} \frac{(-1)^j (2x)^{i-2j}}{j!(i-2j)!}, \quad i \in \{0, 1, 2, \dots\} \quad (4.2)$$

¹⁰It will be interesting to discuss intrinsic relations between the above mentioned relaxation phenomenon to μ_{sin} and the *interpolation/sampling theorem* of Whittaker and others [9], which represents a function f in the cardinal sampling series

$$f(z) = \frac{\sin \pi z}{\pi} \sum_{n \in \mathbb{Z}} f(n) \frac{(-1)^n}{z - n}.$$

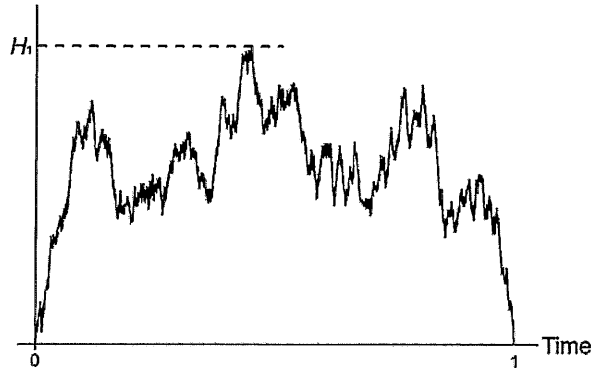


Figure 11: Let H_1 be the maximum of the three-dimensional Bessel bridge with duration 1. The expectation $E[H_1^s]$ is expressed by using the value of Riemann zeta function at $s \in \mathbb{C}$.

with $[a]$ = the greatest integer that is not greater than $a \in \mathbb{R}$. By the equation $P(H_1 \leq h) = \int_0^h p_1(u) du$, the probability density function $p_1(u)$ of H_1 is defined, and s -th moment of the random variable H_1 is calculated by

$$E[H_1^s] = \int_0^\infty dh h^s p_1(s).$$

As discussed by Biane, Pitman, and Yor [2], if we set $E[H_1^s] = 2(\pi/2)^{s/2} \xi(s)$, the following equality is established,

$$\begin{aligned} \xi(s) &= \frac{1}{2} s(s-1) \pi^{-s/2} \Gamma(s/2) \zeta(s) \\ &= \frac{1}{2} + \frac{1}{4} s(s-1) \int_1^\infty du (u^{s/2-1} + u^{(1-s)/2-1}) (\vartheta_3(0, \sqrt{-1}u) - 1), \end{aligned} \quad (4.3)$$

where $\Gamma(z)$ is the gamma function, $\zeta(z)$ is the *Riemann zeta function*,

$$\zeta(z) = \sum_{n=1}^{\infty} \frac{1}{n^z}, \quad \Re z > 1, \quad (4.4)$$

and $\vartheta_3(v, \tau)$ is given by (3.16). See also Chapter 11 in [87].

Then we consider an N -tuples of three-dimensional Bessel bridges conditioned never to collide with each other, $\mathbf{Y}^{(N)}(t) = (Y_1^{(N)}(t), \dots, Y_N^{(N)}(t))$,

$$0 < Y_1^{(N)} < Y_2^{(N)} < \dots < Y_N^{(N)}, \quad 0 < t < 1,$$

with $\mathbf{Y}^{(N)}(0) = \mathbf{Y}^{(N)}(1) = \mathbf{0}$. It is found that (4.1) is generalized for

$$H_N = \max_{0 < t < 1} Y_N^{(N)}$$

as [76, 52]

$$P(H_N \leq h) = \frac{(-1)^N}{2^{N^2} \prod_{i=1}^N \Gamma(2i)} \det_{1 \leq i, j \leq N} \left[\sum_{n \in \mathbb{Z}} H_{2(i+j-1)}(\sqrt{2nh}) e^{-2n^2 h^2} \right]. \quad (4.5)$$

Extensive study of extreme value distributions of noncolliding processes has been reported in [24, 36, 18, 76, 19, 3, 66, 71, 31]. See also [53, 86, 42, 58].

Recently, Forrester, Majumdar, and Schehr clarified an equality between the distribution function $P(H_N \leq h)$ and the partition function of the $2d$ *Yang-Mills theory* on a sphere with a gauge group $\text{Sp}(2N)$ [23]. Moreover, by using this equivalence, they proved that in a scaling limit of $P(H_N \leq h)$, the *Tracy-Widom distribution* of GOE of random matrices [84, 85] is derived.

4.2 Characteristic polynomials of random matrices

Let \mathcal{H} be an $N \times N$ Hermitian random matrix in the GUE. The characteristic polynomial of a variable α is then given by

$$P(\alpha) = \det(\alpha I - \mathcal{H}),$$

where I is the $N \times N$ unit matrix. In the connection with the Riemann zeta function (4.4), statistical property of $P(\alpha)$ has been studied [50, 51, 28, 6]. Here we consider the ensemble average of m -product of characteristic polynomials

$$M_{\text{GUE}}(m, \boldsymbol{\alpha}; N, \sigma^2) = \left\langle \prod_{n=1}^m P(\alpha_n) \right\rangle_{\text{GUE}(N, \sigma^2)} = \left\langle \prod_{n=1}^m \prod_{i=1}^N (\alpha_n - \lambda_i) \right\rangle_{\text{GUE}(N, \sigma^2)}, \quad (4.6)$$

where $\langle \cdot \rangle_{\text{GUE}(N, \sigma^2)}$ denotes the ensemble average in the GUE with variance σ^2 of $N \times N$ Hermitian matrices $\{\mathcal{H}\}$, whose eigenvalues are $\{(\lambda_i)_{i=1}^N\}$. The probability density function of eigenvalues of \mathcal{H} in GUE with variance σ^2 is given by

$$\mu_{N, \sigma^2}(\boldsymbol{\xi}) = \frac{\sigma^{-N^2}}{C_N} \exp\left(-\frac{|\boldsymbol{x}|^2}{2\sigma^2}\right) h_N(\boldsymbol{x})^2, \quad (4.7)$$

$\boldsymbol{\xi} = \sum_{i=1}^N \delta_{x_i} \in \mathfrak{M}$, $x_1 \leq x_2 \leq \dots \leq x_N$, where $C_N = (2\pi)^{N/2} \prod_{i=1}^N \Gamma(i)$, $|\boldsymbol{x}|^2 = \sum_{i=1}^N x_i^2$, and $h_N(\boldsymbol{x})$ is given by (2.7). Then if the expectation of a measurable function F of a random variable $\Xi \in \mathfrak{M}$ with respect to (4.7) is written as

$$\mathbf{E}_{N, \sigma^2}[F(\Xi)] = \int_{\mathbb{W}_N^\Lambda} F(\boldsymbol{\xi}) \mu_{N, \sigma^2}(\boldsymbol{\xi}) d\boldsymbol{x}$$

with setting $\boldsymbol{\xi} = \sum_{i=1}^N \delta_{x_i}$, $\boldsymbol{x} = (x_1, \dots, x_N)$, where $d\boldsymbol{x} = \prod_{i=1}^N dx_i$, (4.6) is written as

$$M_{\text{GUE}}(m, \boldsymbol{\alpha}; N, \sigma^2) = \mathbf{E}_{N, \sigma^2} \left[\prod_{n=1}^m \prod_{X \in \Xi} (\alpha_n - X) \right].$$

In Section 3.3, we showed that the Dyson model $(\Xi(t), \mathbb{P}^\xi)$ is a determinantal process with the correlation kernel (3.12) for any fixed initial configuration $\xi \in \mathfrak{M}_0$ if $\xi(\mathbb{R}) = N \in \mathbb{N}$. Here we consider the situation such that the initial configuration ξ is distributed according to (4.7). Note that by the term $h_N(\boldsymbol{x})^2$ in (4.7), the GUE eigenvalue distribution is in \mathfrak{M}_0 w.p.1.

By using Aspect 2 of the Dyson model, we can prove that this process, denoted by $(\Xi(t), \mathbb{P}^{\mu_{N,\sigma^2}})$, is equivalent with the time shift $t \rightarrow t + \sigma^2$ of the Dyson model starting from the configuration $N\delta_0$ (i.e., all N particles are put at the origin). That is, the equality

$$(\Xi(t), \mathbb{P}^{\mu_{N,\sigma^2}}) = (\Xi(t + \sigma^2), \mathbb{P}^{N\delta_0}) \quad (4.8)$$

holds for arbitrary $\sigma^2 > 0$ in the sense of finite dimensional distribution [35].

This equivalence is highly nontrivial, since even from its special consequence, the following determinantal expression is derived for the ensemble average of product of characteristic polynomials; For any $N, n \in \mathbb{N}$, $\alpha = (\alpha_1, \alpha_2, \dots, \alpha_{2n}) \in \mathbb{C}^{2n}$, $\sigma^2 > 0$,

$$M_{\text{GUE}}(2n, \alpha; N, \sigma^2) = \frac{\gamma_{N,2n} \sigma^{n(2N+n)}}{h_n(\alpha_1, \dots, \alpha_n) h_n(\alpha_{n+1}, \dots, \alpha_{2n})} \\ \times \det_{1 \leq i, j \leq n} \left[\begin{array}{c|cc} 1 & H_{N+n}(\alpha_i/\sqrt{2\sigma^2}) & H_{N+n}(\alpha_{n+j}/\sqrt{2\sigma^2}) \\ \alpha_i - \alpha_{n+j} & H_{N+n-1}(\alpha_i/\sqrt{2\sigma^2}) & H_{N+n-1}(\alpha_{n+j}/\sqrt{2\sigma^2}) \end{array} \right], \quad (4.9)$$

where

$$\gamma_{N,2n} = 2^{-n(2N+2n-1)/2} \prod_{i=2}^n \frac{(N+n-i)!}{(N+n-1)!},$$

and $H_i(x)$ is the i -th Hermite polynomial given by (4.2).

Moreover, in order to simplify the above expression, we can use the following identity, which was recently given by Ishikawa *et al.* [29] as a generalization of the Cauchy determinant

$$\det_{1 \leq i, j \leq n} \left(\frac{1}{x_i + y_j} \right) = \frac{h_n(\mathbf{x}) h_n(\mathbf{y})}{\prod_{i=1}^n \prod_{j=1}^n (x_i + y_j)}.$$

For $n \geq 2$, $\mathbf{x} = (x_1, \dots, x_n)$, $\mathbf{y} = (y_1, \dots, y_n)$, $\mathbf{a} = (a_1, \dots, a_n)$, $\mathbf{b} = (b_1, \dots, b_n) \in \mathbb{C}^n$,

$$= \frac{\det_{1 \leq i, j \leq n} \left[\begin{array}{c|cc} 1 & a_i & \\ y_j - x_i & 1 & b_j \end{array} \right]}{\prod_{i=1}^n \prod_{j=1}^n (y_j - x_i)} \begin{vmatrix} 1 & x_1 & \cdots & x_1^{n-1} & a_1 & a_1 x_1 & \cdots & a_1 x_1^{n-1} \\ 1 & x_2 & \cdots & x_2^{n-1} & a_2 & a_2 x_2 & \cdots & a_2 x_2^{n-1} \\ & & & \cdots & & & & \cdots \\ 1 & x_n & \cdots & x_n^{n-1} & a_n & a_n x_n & \cdots & a_n x_n^{n-1} \\ 1 & y_1 & \cdots & y_1^{n-1} & b_1 & b_1 y_1 & \cdots & b_1 y_1^{n-1} \\ 1 & y_2 & \cdots & y_2^{n-1} & b_2 & b_2 y_2 & \cdots & b_2 y_2^{n-1} \\ & & & \cdots & & & & \cdots \\ 1 & y_n & \cdots & y_n^{n-1} & b_n & b_n y_n & \cdots & b_n y_n^{n-1} \end{vmatrix}. \quad (4.10)$$

Then we have the expression

$$M_{\text{GUE}}(2n, \alpha; N, \sigma^2) = \frac{1}{h_{2n}(\alpha)} \det_{1 \leq i, j \leq 2n} \left[\widehat{H}_{N+i-1}(\alpha_j; \sigma^2) \right], \quad (4.11)$$

where

$$\widehat{H}_i(\alpha; \sigma^2) \equiv \left(\frac{\sigma^2}{2} \right)^{i/2} H_i \left(\frac{\alpha}{\sqrt{2\sigma^2}} \right).$$

This determinantal expression (4.11) is also obtained from the general formula given by Brézin and Hikami as Eq.(14) in [6]. See [13] for recent development of this topic.

4.3 Fomin's determinant for loop-erased random walks and its scaling limit

We consider a network $\Gamma = (V, E, W)$, where $V = \{v_i\}$ and $E = \{e_i\}$ are sets of vertices and of edges of an undirected planar lattice, respectively, and $W = \{w(e)\}_{e \in E}$ is a set of the weight functions of edges. For $a, b \in V$, let π be a walk given by

$$\pi : a = v_0 \xrightarrow{e_1} v_1 \xrightarrow{e_2} v_2 \xrightarrow{e_3} \dots \xrightarrow{e_m} v_m = b$$

where the length of walk is $|\pi| = m \in \mathbb{N}$ and, for each $0 \leq i \leq m-1$, v_i and v_{i+1} are nearest-neighboring vertices in V and $e_i \in E$ is the edge connecting these two vertices. The weight of π is given by $w(\pi) = \prod_{i=1}^m w(e_i)$. For any two vertices of $a, b \in V$, the Green's function of walks $\{\pi : a \rightarrow b\}$ is defined by

$$W(a, b) = \sum_m \sum_{\pi: a \rightarrow b, |\pi|=m} w(\pi).$$

The matrix $W = (W(a, b))_{a, b \in V}$ is called the *walk matrix* of the network Γ .

The loop-erased part of π , denoted by $\text{LE}(\pi)$, is defined recursively as follows. If π does not have self-intersections, that is, all vertices $v_i, 0 \leq i \leq m$ are distinct, then $\text{LE}(\pi) = \pi$. Otherwise, set $\text{LE}(\pi) = \text{LE}(\pi')$, where π' is obtained by removing the first loop it makes. The loop-erasing operator LE maps arbitrary walks to *self-avoiding walks* (SAWs). Note that the map is many-to-one. For each SAW, ζ , the weight $\tilde{w}(\zeta)$ is given by

$$\tilde{w}(\zeta) = \sum_{\pi: \text{LE}(\pi)=\zeta} w(\pi). \quad (4.12)$$

We consider the statistical ensemble of SAWs with the weight (4.12) and call it *loop-erased random walks* (LERWs) [61].

Assume that $A = \{a_1, a_2, \dots, a_N\} \subset V$ and $B = \{b_1, b_2, \dots, b_N\} \subset V$ are chosen so that any walk from a_i to b_j intersects any walk from $a_{i'}, i' > i$, to $b_{j'}, j' < j$. The weight of N -tuples of independent walks $a_1 \xrightarrow{\pi_1} b_1, \dots, a_N \xrightarrow{\pi_N} b_N$ is given by the product of N weights $\prod_{i=1}^N w(\pi_i)$. Then we consider N -tuples of walks $(\pi_1, \pi_2, \dots, \pi_N)$ conditioned so that, for any $1 \leq i < j \leq N$, the walk π_j has no common vertices with the loop-erased part of π_i ;

$$\text{LE}(\pi_i) \cap \pi_j = \emptyset, \quad 1 \leq i < j \leq N. \quad (4.13)$$

See Fig.12. By definition, $\text{LE}(\pi_j)$ is a part of π_j , and thus nonintersection of any pair of loop-erased parts is concluded from (4.13);

$$\text{LE}(\pi_i) \cap \text{LE}(\pi_j) = \emptyset, \quad 1 \leq i < j \leq N.$$

Fomin proved that total weight of N -tuples of walks satisfying such a version of nonintersection condition is given by the minor of walk matrix, $\det(W_{A,B}) \equiv \det_{a \in A, b \in B} (W(a, b))$

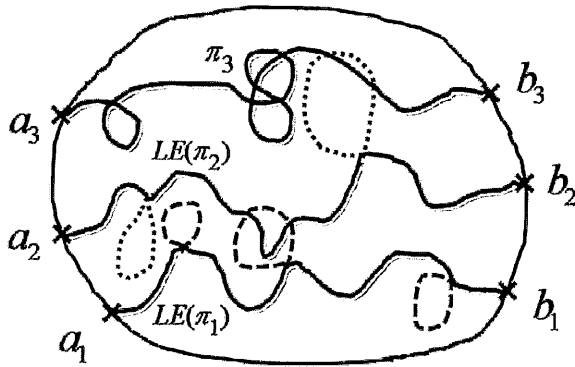


Figure 12: The situation $LE(\pi_j) \cap \pi_3 = \emptyset, j = 1, 2$ is illustrated in a planar domain D , where $A = (a_1, a_2, a_3)$ and $B = (b_1, b_2, b_3)$ are all boundary points of ∂D . In this figure, $LE(\pi_1)$ and $LE(\pi_2)$ denoted by solid curves are the loop-erased parts of the walks $\pi_1 : a_1 \rightarrow b_1$ and $\pi_2 : a_2 \rightarrow b_2$, respectively. The third walk $\pi_3 : a_3 \rightarrow b_3$ can be self-intersecting, but it does not intersect with $LE(\pi_1)$ nor $LE(\pi_2)$.

[21]. This minor is called *Fomin's determinant* and Fomin's formula is expressed by the equality [21, 61]

$$\det(W_{A,B}) = \sum_{LE(\pi_i) \cap \pi_j = \emptyset, i < j} \prod_{k=1}^N w(\pi_k). \quad (4.14)$$

Kozdron and Lawler [55] consider continuum limit (the diffusion scaling limit) of Fomin's determinantal system of loop-erased random walks in the complex plane \mathbb{C} , where the initial and the final points $A = \{a_i\}$ and $B = \{b_i\}$ of paths can be put on the boundaries of the domains ∂D . By the diffusion scaling limit each random walk will converge to a path of complex BM. We should note that, however, the characteristics of BM look more similar to those of a surface than those of a curve. It implies that the Brownian path has loops on every scale and then the loop-erasing procedure mentioned above does not make sense for BM in the plane, since we can not decide which loop is the first one. Kozdron and Lawler proved explicitly, however, that the continuum limit of Fomin's determinant of the Green's functions of random walks converges to that of the Green's functions of Brownian motions [55]. This will enable us to discuss *nonintersecting systems of loop-erased Brownian paths* in the sense of Fomin (4.13). Moreover, Kozdron [54] showed that 2×2 Fomin's determinant representing the event $LE(\beta_1) \cap \beta_2 = \emptyset$ for two complex Brownian paths (β_1, β_2) is proportional to the probability that $\gamma_{SLE^{(3)}} \cap \beta = \emptyset$, where $\gamma_{SLE^{(3)}}$ and β denote the $SLE^{(3)}$ path and a complex Brownian path, respectively. On the other hand, Lawler and Werner gave a correct way to add 'Brownian loops' to an $SLE^{(3)}$ path to obtain a complex Brownian path [63]. These results imply that the scaling limit of loop-erased part of complex Brownian path is described by the $SLE^{(3)}$ path, as announced in Section 1.4.

Setting a sequence of chambers in a planar domain, Sato and the present author observe the first passage points at which N -tuples of complex Brownian paths $(\beta_1, \dots, \beta_N)$

first enter each chamber, under the condition that the loop-erased parts $(LE(\beta_1), \dots, LE(\beta_N))$ make a nonintersecting system in the domain in the sense of Fomin (4.13) [75]. It is proved that the system of first passage points is a determinantal point process in the planar domain, in which the correlation kernel is of Eynard-Mehta type [75]. Interpretation of this result in terms of ‘mutually avoiding SLE paths’ [56, 15] will be an interesting future problem.

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