

A CONTINUOUS SEMIGROUP OF NOTIONS OF INDEPENDENCE BETWEEN THE CLASSICAL AND THE FREE ONE

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In this paper, we investigate a continuous family of notions of independence which interpolates between the classical and free ones for noncommutative random variables. These notions are related to the liberation process introduced by Voiculescu. To each notion of independence correspond new convolutions of probability measures, for which we establish formulae and of which we compute simple examples. We prove that there exists no reasonable analogue of classical and free cumulants associated to these notions of independence.

Introduction. Let μ and ν be two Borel probability measures on the real line \mathbb{R} . The classical convolution of μ and ν is the probability measure on \mathbb{R} , denoted by $\mu * \nu$, which is the distribution of the sum of two classical independent random variables with respective distributions μ and ν . Let us describe $\mu * \nu$ in an alternative way. To each $n \times n$ matrix with eigenvalues $\lambda_1, \dots, \lambda_n$, we associate its spectral measure, which is the probability measure $\frac{1}{n} \sum_{i=1}^n \delta_{\lambda_i}$. Let $(A_n)_{n \geq 1}$ and $(B_n)_{n \geq 1}$ be two sequences of diagonal real matrices, with A_n and B_n of size n for all $n \geq 1$, such that the spectral measure of A_n (resp., of B_n) converges, as n tends to infinity, to μ (resp., to ν). For each $n \geq 1$, let S_n be a random matrix chosen uniformly among the $n!$ permutation matrices of size n . Then the spectral measure of $A_n + S_n B_n S_n^{-1}$ converges, as n tends to infinity, to $\mu * \nu$.

If we replace, for each $n \geq 1$, the matrix S_n by a random matrix U_n chosen in the unitary group $U(n)$ according to the Haar measure, then the spectral measure of $A_n + U_n B_n U_n^{-1}$ converges, as n tends to infinity, to the free convolution of μ and ν , a probability measure on \mathbb{R} denoted by $\mu \boxplus \nu$.

This way of describing classical and free convolutions suggests a natural way to interpolate between them. Indeed, consider, for all $n \geq 1$, a properly scaled Brownian motion $(U_{n,t})_{t \geq 0}$ issued from the identity matrix on the unitary group $U(n)$. Given $t \in [0, +\infty)$, one may consider the spectral measure of $A_n + U_{n,t} S_n B_n S_n^{-1} U_{n,t}^{-1}$, and ask for the limit of this distribution as n tends to infinity. For $t = 0$, the matrix $U_{n,0}$ is the identity matrix and we find the classical convolution of μ and ν . For $t = +\infty$, that is, when $U_{n,t}$ is replaced by its limit in distribution as t tends to infinity, which is a uniformly distributed unitary matrix, we recover the free convolution of μ and ν . For any other $t \in (0, +\infty)$, it turns out

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that one finds a probability measure which depends only on μ, ν and t and which we denote by $\mu *_t \nu$.

Consider, for example, the case where $\mu = \nu = \frac{1}{2}(\delta_1 + \delta_{-1})$. Then $\mu * \nu = \frac{1}{4}\delta_{-2} + \frac{1}{2}\delta_0 + \frac{1}{4}\delta_2$ and it is well known that $\mu \boxplus \nu = \mathbb{1}_{[-2,2]}(x) \frac{dx}{\pi\sqrt{4-x^2}}$, a dilation of the arcsine law ([10], Example 12.8). One may wonder which probability measures interpolate between $\mu * \nu$ and $\mu \boxplus \nu$. We will prove that

$$\forall t \geq 0 \quad \frac{\delta_1 + \delta_{-1}}{2} *_t \frac{\delta_1 + \delta_{-1}}{2} = \mathbb{1}_{[-2,2]}(x) \frac{\rho_{4t}(e^{4i \arccos(x/2)})}{\pi\sqrt{4-x^2}} dx.$$

Here, for all $t > 0$ and $\theta \in \mathbb{R}$, $\rho_t(e^{i\theta})$ is the density at $e^{i\theta}$, with respect to the uniform probability measure on the unit circle, of the distribution of the free unitary Brownian motion at time t . This distribution is also the limit, as n tends to infinity, of the spectral measure of $U_{n,t}$. There is no simple formula for this distribution, which apparently has to be taken as a fundamental function in any problem involving the large n asymptotics of the Brownian motion on the unitary group $U(n)$. However, the moments of this distribution are known since Biane first computed them [4]. It follows for instance from the previous expression that $\mu *_t \nu$ has a density with respect to the Lebesgue measure for all $t > 0$ and that its support, which one can compute for all $t \geq 0$, is the whole interval $[-2, 2]$ if and only if $t \geq 1$.

We have unfortunately not been able yet to give any other nontrivial examples where $\mu *_t \nu$ can be computed. However, the tools offered by the computation rules given in Section 3 allow one to develop a related notion of infinite divisibility, and this will be done in the forthcoming paper [2] where a random matrix model is presented to interpolate between Gaussian and Wigner distributions.

The family of operations $*_t$ is really just a by-product of a more fundamental construction, which is that of a continuous family of independence (or dependence) structures between noncommutative random variables which interpolates between classical independence and freeness. Indeed, we will define, for all $t \in [0, +\infty]$, a notion of independence between two subalgebras of a noncommutative probability space, which we call t -freeness, which, for $t = 0$ (resp., $t = +\infty$), coincides with classical independence (resp., freeness) and which, for all $t \geq 0$, is related to the so-called *liberation process* introduced by Voiculescu [12]. Once this structure is defined, it is straightforward to define additive or multiplicative convolution of t -free self-adjoint or unitary elements, thus giving rise to several operations on probability measures: additive or multiplicative convolution of probability measures with compact support on \mathbb{R} , denoted by $*_t$ and \odot_t ; multiplicative convolution of probability measures on the unit circle, also denoted by \odot_t .

The idea of seeking a continuous way of passing from classical to free independence is presumably as old as the theory of free probability itself, but the research of such a continuum has been broken off by a paper of Roland Speicher in 1997

[11], where he has shown that no other notion of independence than the classical and the free ones can be the base of a *reasonable* probability theory. Indeed t -freeness does not satisfy all the axioms enforced by Speicher because it is not an *associative* notion of independence. This axiom of associativity states, roughly, that if X, Y, Z are three random variables such that X is independent of Y and Z is independent of $\{X, Y\}$, then X must be independent of $\{Y, Z\}$. Instead of this, what is true with t -freeness is that for all $s, t \geq 0$, if X, Y, Z are three random variables such that X is t -free with Y and Y is s -free with Z , then under certain additional hypotheses, X will be $(s + t)$ -free with Z . This is of course related to the semi-group property of the Brownian motion.

There are several ways to characterize and deal with independence and freeness. The first one, which we have already mentioned, is to relate them with matrix models. The second one is to describe them by means of computation rules: the expectation factorizes with respect to independent subfamilies of random variables, whereas the expectation of a product of free elements can be computed using the fact that if x_1, \dots, x_n are centered and successively free, then their product is centered. The third way to describe independence and freeness is to identify integral transforms which linearize them (namely the logarithm of Fourier transform or the R -transform). This amounts to describing classical and free cumulants. The last way, a bit more abstract, is to consider tensor or free products: a family of random variables is independent (resp., free) if and only if it can be realized on a tensor product (resp., free product) of probability spaces.

In the present paper, we look for the analogues of all these approaches for the notion of t -freeness. We begin, in Section 2, by giving the definition of a t -free product and presenting the corresponding random matrix model. Then, in Section 3, we state the computation rules, which are best understood as a family of differential equations. Finally, in Section 4, we prove that no notion of cumulants of order greater than 6 can be associated to the notion of t -freeness. More precisely, we show that there does not exist a universally defined 7-linear form on any noncommutative probability space with the property that this form vanishes whenever it is evaluated on arguments which can be split into two nonempty subfamilies which are t -free, unless $t = 0$ or $t = +\infty$. This can be summarized in Figure 1.

	Matrix model	Computation rules	Cumulants	Algebraic structures
Indep.	$A + SBS^{-1}$	Factorization	Class. cumulants	Tensor product
t -freeness	$A + U_t SBS^{-1} U_t^{-1}$	Differential system	<i>Do not exist</i>	t -free product
Freeness	$A + UBU^{-1}$	$\varphi(x_1 \cdots x_n) = 0$	Free cumulants	Free product

FIG. 1. *The main computation rule for freeness is that $\varphi(x_1 \cdots x_n) = 0$ as soon as x_1, \dots, x_n are successively free and centered. For t -freeness, the computation rules are best expressed as a differential system relating the distributions of the t -free products of two families of random variables for different values of t .*

1. Preliminaries. In this section, we review the notions of noncommutative probability which are relevant to the definition of t -freeness.

1.1. *Probability space, distribution.* Noncommutative probability is based on the following generalization of the notion of probability space.

DEFINITION 1.1 (Noncommutative probability space). A *noncommutative probability space* is a pair (\mathcal{A}, φ) , where:

- \mathcal{A} is an algebra over \mathbb{C} with a unit element denoted by 1, endowed with an operation of adjunction $x \mapsto x^*$ which is \mathbb{C} -antilinear, involutive and satisfies $(xy)^* = y^*x^*$ for all $x, y \in \mathcal{A}$,
- $\varphi: \mathcal{A} \rightarrow \mathbb{C}$ is a linear form on \mathcal{A} , satisfying $\varphi(1) = 1$, $\varphi(xy) = \varphi(yx)$, $\varphi(x^*) = \overline{\varphi(x)}$ and $\varphi(xx^*) \geq 0$ for all $x, y \in \mathcal{A}$.

The linear form φ is often called the *expectation* of the noncommutative probability space.

Two fundamental examples are the algebra $L^{\infty}(\Omega, \Sigma, \mathbb{P})$ of complex-valued random variables with moments of all orders on a classical probability space, endowed with the complex conjugation and the expectation [we will say that this noncommutative probability space is *inherited* from $(\Omega, \mathcal{A}, \mathbb{P})$]; and the algebra $\mathcal{M}_n(\mathbb{C})$ endowed with the matricial adjunction and the normalized trace.

DEFINITION 1.2 (Noncommutative distribution). Let (\mathcal{A}, φ) be a noncommutative probability space. The *noncommutative distribution* of a family (a_1, \dots, a_n) of elements of \mathcal{A} with respect to φ is the linear map defined on the space of polynomials in the noncommutative variables $X_1, X_1^*, \dots, X_n, X_n^*$ which maps any such polynomial P to $\varphi(P(a_1, a_1^*, \dots, a_n, a_n^*))$.

The link between the classical notion of distribution and the noncommutative one is the following. Consider a self-adjoint element a in a noncommutative probability space (\mathcal{A}, φ) , that is, an element such that $a = a^*$. Since $\varphi(xx^*) \geq 0$ for all $x \in \mathcal{A}$, the distribution of a is a linear form on $\mathbb{C}[X]$ which is nonnegative on the polynomials which are nonnegative on the real line. Hence, it can be represented as the integration with respect to a probability measure on the real line. This probability measure is unique if and only if it is determined by its moments, which is in particular the case when it has compact support, or equivalently when there exists a constant M such that for all $n \geq 0$, one has $\varphi(a^{2n}) \leq M^{2n}$.

Similarly, the distribution of a unitary element u , that is, an element such that $uu^* = u^*u = 1$, is the integration with respect to a probability measure on the unit circle of \mathbb{C} . Since the circle is compact, there is no issue of uniqueness in this case.

1.2. *Independence, freeness and random matrices.*

1.2.1. *Definitions and basic properties.* We shall recall the definitions of the two notions of independence in a noncommutative probability space between which our main purpose is to interpolate. The first one is a straightforward translation of the classical notion of independence in the noncommutative setting, which coincides with the original notion in the case of a noncommutative probability space inherited from a classical one. The second one is the notion of freeness, as defined by Voiculescu [13], which is called freeness.

In this paper, by a subalgebra of the algebra of a noncommutative probability space, we shall always mean a subalgebra which contains 1 and which is stable under the operation $x \mapsto x^*$.

DEFINITION 1.3 (Independence and freeness). Let (\mathcal{M}, φ) be a noncommutative probability space. The kernel of φ will be called the set of *centered elements*. Consider a family $(\mathcal{A}_i)_{i \in I}$ of subalgebras of \mathcal{M} .

- The family $(\mathcal{A}_i)_{i \in I}$ is said to be *independent* if
 - (i) for all $i \neq j \in I$, \mathcal{A}_i and \mathcal{A}_j commute,
 - (ii) for all $n \geq 1$, $i_1, \dots, i_n \in I$ pairwise distinct, for all family $(a_1, \dots, a_n) \in \mathcal{A}_{i_1} \times \dots \times \mathcal{A}_{i_n}$ of centered elements, the product $a_1 \cdots a_n$ is also centered.
- The family $(\mathcal{A}_i)_{i \in I}$ is said to be *free* if for all $n \geq 1$, $i_1, \dots, i_n \in I$ such that $i_1 \neq i_2, i_2 \neq i_3, \dots, i_{n-1} \neq i_n$, for all family $(a_1, \dots, a_n) \in \mathcal{A}_{i_1} \times \dots \times \mathcal{A}_{i_n}$ of centered elements, the product $a_1 \cdots a_n$ is also centered.

On a classical probability space $(\Omega, \Sigma, \mathbb{P})$, a family $(\Sigma_i)_{i \in I}$ of sub- σ -fields of Σ is independent with respect to \mathbb{P} if and only if the subalgebras $(L^{\infty-}(\Omega, \Sigma_i, \mathbb{P}))_{i \in I}$ of $(L^{\infty-}(\Omega, \Sigma, \mathbb{P}), \mathbb{E})$ are independent in the sense of the definition above.

In the classical setting again, a family of random variables is independent if and only if its joint distribution is the tensor product of the individual ones. In the following definition and proposition, we translate this statement into our vocabulary, and give its analogue for freeness. These definitions prepare those which we will give later for t -freeness.

DEFINITION 1.4 (Tensor and free product). Let $(\mathcal{A}_1, \varphi_1)$ and $(\mathcal{A}_2, \varphi_2)$ be two noncommutative probability spaces.

- Their *tensor product*, denoted by $(\mathcal{A}_1, \varphi_1) \otimes (\mathcal{A}_2, \varphi_2)$, is the noncommutative probability space with algebra the tensor product of unital algebras $\mathcal{A}_1 \otimes \mathcal{A}_2$, on which the adjoint operation and the expectation are defined by

$$\forall (x_1, x_2) \in \mathcal{A}_1 \times \mathcal{A}_2 \quad (x_1 \otimes x_2)^* = x_1^* \otimes x_2^*, \quad \varphi(x_1 \otimes x_2) = \varphi_1(x_1)\varphi_2(x_2).$$

- Their *free product*, denoted by $(\mathcal{A}_1, \varphi_1) * (\mathcal{A}_2, \varphi_2)$, is the noncommutative probability space with algebra the free product of unital algebras $\mathcal{A}_1 * \mathcal{A}_2$, with adjoint operation and expectation defined uniquely by the fact that for all $n \geq 1$, for all $i_1 \neq \dots \neq i_n \in \{1, 2\}$, for all $(x_1, \dots, x_n) \in \mathcal{A}_{i_1} \times \dots \times \mathcal{A}_{i_n}$,

$$(x_1 \cdots x_n)^* = x_n^* \cdots x_1^*$$

and $x_1 \cdots x_n$ is centered whenever all x_i 's are.

This definition can easily be extended to products of finite or infinite families of noncommutative probability spaces, but we have restricted ourselves to what is needed in this article. We can now explain the link between these products and the notions of independence and freeness.

PROPOSITION 1.5 (Characterization of independence and freeness). *Let (\mathcal{M}, φ) be a noncommutative probability space. Let $\mathcal{A}_1, \mathcal{A}_2$ be subalgebras of \mathcal{A} . Then the family $(\mathcal{A}_1, \mathcal{A}_2)$ is*

- *independent if and only if \mathcal{A}_1 commutes with \mathcal{A}_2 and the unique algebra morphism defined from $\mathcal{A}_1 \otimes \mathcal{A}_2$ to \mathcal{M} which, for all $(a_1, a_2) \in \mathcal{A}_1 \times \mathcal{A}_2$, maps $a_1 \otimes 1$ to a_1 and $1 \otimes a_2$ to a_2 , preserves the expectation from $(\mathcal{A}_1, \varphi|_{\mathcal{A}_1}) \otimes (\mathcal{A}_2, \varphi|_{\mathcal{A}_2})$ to (\mathcal{M}, φ) ,*
- *free if and only if the unique algebra morphism defined from the free product of unital algebras $\mathcal{A}_1 * \mathcal{A}_2$ to \mathcal{M} which, restricted to $\mathcal{A}_1 \cup \mathcal{A}_2$ is the canonical injection, preserves the expectation from $(\mathcal{A}_1, \varphi|_{\mathcal{A}_1}) * (\mathcal{A}_2, \varphi|_{\mathcal{A}_2})$ to (\mathcal{M}, φ) .*

Let us finally recall the definition of the free analogue of the classical convolution, which is meaningful thanks to the last proposition.

DEFINITION 1.6 (Additive free convolution). Let μ and ν be two probability measures on \mathbb{R} . The distribution of the sum of two free self-adjoint elements with respective distributions μ and ν depends only on μ and ν and will be called the *free additive convolution* of μ and ν , and be denoted by $\mu \boxplus \nu$.

1.2.2. *Asymptotic behavior of random matrices.* In this section, we recall matrix models for the classical and free convolution. The main notion of convergence which is involved is the following.

DEFINITION 1.7 (Convergence in noncommutative distribution). Let p be a positive integer and let, for each $n \geq 1$, $(M(1, n), \dots, M(p, n))$ be a family of $n \times n$ random matrices. This family is said to *converge in noncommutative distribution* if its noncommutative distribution converges in probability to a nonrandom one, that is, if the normalized trace of any word in the $M(i, n)$'s and the $M(i, n)$'s* converges in probability to a constant.

THEOREM 1.8 (Asymptotic independence and asymptotic freeness). *Let us fix $p, q \geq 1$. For each $n \geq 1$, let $\mathcal{F}_n = (A(1, n), \dots, A(p, n), B(1, n), \dots, B(q, n))$ be a family of $n \times n$ random matrices and assume that the sequence $(\mathcal{F}_n)_{n \geq 1}$ converges in noncommutative distribution. Assume also that for all $r \geq 1$, the entries of these random matrices are uniformly bounded in L^r .*

- *Assume that these matrices are diagonal and consider, for each n , the matrix S_n of a uniformly distributed random permutation of $\{1, \dots, n\}$ independent of the family \mathcal{F}_n . Then the family*

$$(1) \quad (A(1, n), \dots, A(p, n), S_n B(1, n) S_n^{-1}, \dots, S_n B(q, n) S_n^{-1})$$

converges in distribution to the distribution of a commutative family $(a_1, \dots, a_p, b_1, \dots, b_q)$ of elements of a noncommutative probability space such that the algebras generated by $\{a_1, \dots, a_p\}$ and $\{b_1, \dots, b_q\}$ are independent.

- *Consider, for each n , the matrix U_n of a uniformly distributed random unitary n by n matrix independent of the family \mathcal{F}_n . Then the family*

$$(2) \quad (A(1, n), \dots, A(p, n), U_n B(1, n) U_n^{-1}, \dots, U_n B(q, n) U_n^{-1})$$

converges in distribution to the distribution of a family $(a_1, \dots, a_p, b_1, \dots, b_q)$ of elements of a noncommutative probability space such that the algebras generated by $\{a_1, \dots, a_p\}$ and $\{b_1, \dots, b_q\}$ are free.

REMARK 1.9. The hypothesis of uniform boundedness of the entries of the matrices in each L^r could be sharply weakened for the first part of the theorem if, instead of asking for the convergence of the noncommutative distribution of the family (1), one would ask for the weak convergence of the empirical joint spectral measure. This would amount to choosing, as set of test functions, the set of bounded continuous functions of $p + q$ variables instead of the set of polynomials in $p + q$ variables (see [3], where this is precisely proved).

The first part of this theorem is much simpler than the second but seems to be also less well known. It is in any case harder to locate a proof in the literature, so that we offer one. We shall need the following lemma. We denote by $\|\cdot\|_2$ the usual Hermitian norm on \mathbb{C}^n .

LEMMA 1.10. *Let, for each $n \geq 1$, $x(n) = (x_{n,1}, \dots, x_{n,n})$ and $y(n) = (y_{n,1}, \dots, y_{n,n})$ be two complex random vectors defined on the same probability space such that the random variables*

$$\overline{x(n)} = \frac{x_{n,1} + \dots + x_{n,n}}{n}, \quad \overline{y(n)} = \frac{y_{n,1} + \dots + y_{n,n}}{n}$$

converge in probability to constant limits x, y as n tends to infinity. Suppose moreover that the sequences $\frac{1}{n} \|x(n)\|_2^2$ and $\frac{1}{n} \|y(n)\|_2^2$ are bounded in L^2 . Consider, for

all n , a uniformly distributed random permutation σ_n of $\{1, \dots, n\}$, independent of $(x(n), y(n))$, and define $y_{\sigma_n}(n) := (y_{n,\sigma_n(1)}, \dots, y_{n,\sigma_n(n)})$. Then the scalar product

$$\frac{1}{n} \langle x(n), y_{\sigma_n}(n) \rangle = \frac{x_{n,1}y_{n,\sigma_n(1)} + \dots + x_{n,n}y_{n,\sigma_n(n)}}{n}$$

converges in probability to xy as n tends to infinity.

PROOF. First of all, note that one can suppose that for all n , $\overline{x(n)} = \overline{y(n)} = 0$ almost surely. Indeed, if the result is proved under this additional hypothesis, then since for all n , one has

$$\frac{1}{n} \langle x(n), y_{\sigma_n}(n) \rangle = \frac{1}{n} \langle x(n) - \overline{x(n)} \cdot 1_n, y_{\sigma_n}(n) - \overline{y(n)} \cdot 1_n \rangle + \overline{x(n)} \cdot \overline{y(n)}$$

[with $1_n = (1, \dots, 1)$],

the result holds for general $x(n), y(n)$. So we henceforth assume that for all n , $\overline{x(n)} = \overline{y(n)} = 0$. The equality $\overline{y(n)} = 0$ implies, for all n and all $i, j = 1, \dots, n$, that

$$\mathbb{E}[y_{n,\sigma_n(i)}y_{n,\sigma_n(j)} \mid x(n), y(n)] = \begin{cases} \frac{1}{n} \|y(n)\|_2^2, & \text{if } i = j, \\ -\frac{1}{n(n-1)} \|y(n)\|_2^2, & \text{if } i \neq j. \end{cases}$$

Then, using the fact that $\overline{x(n)} = 0$, we have

$$\begin{aligned} \mathbb{E}\left[\frac{1}{n^2} \langle x(n), y_{\sigma_n}(n) \rangle^2\right] &= \mathbb{E}\left[\frac{1}{n^3(n-1)} \|x(n)\|_2^2 \|y(n)\|_2^2 + \frac{1}{n^3} \|x(n)\|_2^2 \|y(n)\|_2^2\right] \\ &= O\left(\frac{1}{n}\right), \end{aligned}$$

which completes the proof. \square

PROOF OF THEOREM 1.8. The second point is a well-known result of Voiculescu (see [13]). To prove the first one, we shall prove that the normalized trace any word in the random matrices $A(1, n), \dots, A(p, n), S_n B(1, n)S_n^{-1}, \dots, S_n B(q, n)S_n^{-1}$ converges to a constant which is the product in two terms: the limiting normalized trace of the $A(i, n)$'s and the $A(i, n)^*$'s which appear in the word on one hand and the limiting normalized trace of the $B(j, n)$'s and the $B(j, n)^*$'s which appear in the word on the other hand. Since the $A(i, n)$'s, the $A(i, n)^*$'s, the $S_n B(j, n)S_n^{-1}$'s and the $S_n B(j, n)^*S_n^{-1}$'s commute, are uniformly bounded and their noncommutative distribution converges, this amounts to proving that if $M(n), N(n)$ are two diagonal random matrices with entries uniformly bounded in L^r for all $r \geq 1$, whose normalized traces converge in probability to constants m, n , then for S_n the matrix of a uniform random permutation of $\{1, \dots, n\}$ independent of $(M(n), N(n))$, the normalized trace of $M(n)S_n N(n)S_n^{-1}$ converges to mn . This follows directly from the previous lemma and the proof is complete. \square

COROLLARY 1.11 (Matricial model for classical and free convolutions). *Let μ, ν be two probability measures on the real line. Let, for each $n \geq 1$, M_n, N_n be n by n diagonal random matrices with empirical spectral measures converging weakly in probability to μ and ν , respectively. For each $n \geq 1$, let S_n (resp., U_n) be a uniformly distributed n by n permutation (resp., unitary) random matrix independent of (M_n, N_n) . Then*

- *the empirical spectral measure of $M_n + S_n N_n S_n^{-1}$ converges weakly in probability to the classical convolution $\mu * \nu$ of μ and ν ,*
- *the empirical spectral measure of $M_n + U_n N_n U_n^{-1}$ converges weakly in probability to the free convolution $\mu \boxplus \nu$ of μ and ν .*

PROOF. In the case where μ, ν have compact supports and the entries of the diagonal matrices M_n, N_n are uniformly bounded, it is a direct consequence of the previous theorem. The general case can easily be deduced using functional calculus, like in the proof of Theorem 3.13 of [1]. \square

1.3. *Unitary Brownian motion, free unitary Brownian motion.* In this paragraph, we give a brief survey of the definition and the main convergence result for the Brownian motion on the unitary group.

Let $n \geq 1$ be an integer. Let \mathcal{H}_n denote the n^2 -dimensional real linear subspace of $\mathcal{M}_n(\mathbb{C})$ which consists of Hermitian matrices. On $\mathcal{M}_n(\mathbb{C})$, we denote by Tr the usual trace and by $\text{tr} = \frac{1}{n} \text{Tr}$ the normalized trace. Let us endow \mathcal{H}_n with the scalar product $\langle \cdot, \cdot \rangle$ defined by

$$\forall A, B \in \mathcal{H}_n \quad \langle A, B \rangle = n \text{Tr}(A^* B) = n \text{Tr}(AB).$$

There is a linear Brownian motion canonically attached to the Euclidean space $(\mathcal{H}_n, \langle \cdot, \cdot \rangle)$. It is the unique Gaussian process H indexed by \mathbb{R}_+ with values in \mathcal{H}_n such that for all $s, t \in \mathbb{R}_+$ and all $A, B \in \mathcal{H}_n$, one has

$$\mathbb{E}[\langle H_s, A \rangle \langle H_t, B \rangle] = \min(s, t) \langle A, B \rangle.$$

Let us consider the following stochastic differential equation:

$$U_0 = I_n, \quad dU_t = i(dH_t)U_t - \frac{1}{2}U_t dt,$$

where $(U_t)_{t \geq 0}$ is a stochastic process with values in $\mathcal{M}_n(\mathbb{C})$. This linear equation admits a strong solution. The process $(U_t^*)_{t \geq 0}$, where U_t^* denotes the adjoint of U_t , satisfies the stochastic differential equation

$$U_0^* = I_n, \quad dU_t^* = -iU_t^* dH_t - \frac{1}{2}U_t^* dt.$$

An application of Itô's formula to the process $U_t U_t^*$ shows that, for all $t \geq 0$, $U_t U_t^* = I_n$. This proves that the process $(U_t)_{t \geq 0}$ takes its values in the unitary group $U(n)$.

DEFINITION 1.12. The process $(U_t)_{t \geq 0}$ is called the *unitary Brownian motion of dimension n* .

As n tends to infinity, the unitary Brownian motion has a limit in distribution which we now describe. For all $t \geq 0$, the numbers

$$e^{-kt/2} \sum_{j=0}^{k-1} \frac{(-t)^j}{j!} \binom{k}{j+1} k^{j-1}, \quad k \geq 0,$$

are the moments of a unique probability measure on the set $\mathbb{U} = \{z \in \mathbb{C} : |z| = 1\}$ invariant by the complex conjugation. We denote this probability measure by ν_t . The following definition was given by Biane in [4].

DEFINITION 1.13. Let (\mathcal{A}, τ) be a noncommutative probability space. We say that a collection $(u_t)_{t \geq 0}$ of unitary elements of \mathcal{A} is a *free unitary Brownian motion* if the following conditions hold.

- For all $s, t \geq 0$ such that $s \leq t$, the distribution of $u_t u_s^*$ is the probability measure ν_{t-s} .
- For all positive integer m , for all $0 \leq t_1 \leq t_2 \leq \dots \leq t_m$, the elements $u_{t_1} u_0^*, u_{t_2} u_{t_1}^*, \dots, u_{t_m} u_{t_{m-1}}^*$ are free.

In the same paper, Biane has proved the following convergence result.

THEOREM 1.14. For each $n \geq 1$, let $(U_{n,t})_{t \geq 0}$ be a Brownian motion on the unitary group $U(n)$. As n tends to infinity, the collection of random matrices $(U_{n,t})_{t \geq 0}$ converges in noncommutative distribution to a free unitary Brownian motion.

2. A continuum of notions of independence. In this section, we shall define a family indexed by a real number $t \in [0, +\infty]$ of relations between two subalgebras of a noncommutative probability space which passes from the classical independence (which is the case $t = 0$) to freeness (which is the “limit” when t tends to infinity). We start with the definition of the t -free product of two noncommutative probability spaces. In a few words, it is the space obtained by conjugating one of them, in their tensor product, by a free unitary Brownian motion at time t , free with the tensor product.

Fix $t \in [0, +\infty]$ and let $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$ be two noncommutative probability spaces. Let $(\mathcal{U}^{(t)}, \varphi_{\mathcal{U}^{(t)}})$ be the noncommutative probability space generated by a single unitary element u_t whose distribution is that of a free unitary Brownian motion at time t (with the convention that a free unitary Brownian motion at time $+\infty$ is a Haar unitary element, that is, a unitary element whose distribution is the uniform law on the unit circle of \mathbb{C}).

DEFINITION 2.1 (*t*-free product). The *t*-free product of $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$, defined up to an isomorphism of noncommutative probability spaces, is the noncommutative probability space $(\mathcal{C}, \varphi|_{\mathcal{C}})$, where \mathcal{C} is the subalgebra generated by \mathcal{A} and $u_t \mathcal{B} u_t^*$ in

$$(\mathcal{X}, \varphi) := [(\mathcal{A}, \varphi_{\mathcal{A}}) \otimes (\mathcal{B}, \varphi_{\mathcal{B}})] * (\mathcal{U}^{(t)}, \varphi_{\mathcal{U}^{(t)}}).$$

A few simple observations are in order.

REMARK 2.2. Both $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$ can be identified with subalgebras of the algebra of their *t*-free product [namely with $(\mathcal{A}, \varphi|_{\mathcal{A}})$ and $(u_t \mathcal{B} u_t^*, \varphi|_{u_t \mathcal{B} u_t^*})$]. More specifically, if one defines

$$\begin{aligned} \mathcal{A}_{st} &:= \{a \in \mathcal{A}; \varphi_{\mathcal{A}}(a) = 0, \varphi_{\mathcal{A}}(aa^*) = 1\}, \\ \mathcal{B}_{st} &:= \{b \in \mathcal{B}; \varphi_{\mathcal{B}}(b) = 0, \varphi_{\mathcal{B}}(bb^*) = 1\}, \end{aligned}$$

then any element in the algebra of the *t*-free product $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$ can be uniquely written as a constant term plus a linear combination of words in the elements of $\mathcal{A}_{st} \cup u_t \mathcal{B}_{st} u_t^*$ where no two consecutive letters both belong to \mathcal{A}_{st} or to $u_t \mathcal{B}_{st} u_t^*$.

REMARK 2.3. As a consequence, since u_t is unitary and (u_t, u_t^*) has the same noncommutative distribution as (u_t^*, u_t) , the *t*-free product of $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$ is clearly isomorphic, as a noncommutative probability space, to the *t*-free product of $(\mathcal{B}, \varphi_{\mathcal{B}})$ and $(\mathcal{A}, \varphi_{\mathcal{A}})$.

REMARK 2.4. Another consequence of Remark 2.2 is that as a unital algebra, the algebra of the *t*-free product of $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$ is isomorphic to the free product of the unital algebras $\mathcal{A}/\tilde{\mathcal{A}}$ and $\mathcal{B}/\tilde{\mathcal{B}}$, where $\tilde{\mathcal{A}}$ (resp., $\tilde{\mathcal{B}}$) is the bilateral ideal of the elements x of \mathcal{A} (resp., of \mathcal{B}) such that $\varphi_{\mathcal{A}}(xx^*) = 0$ [resp., $\varphi_{\mathcal{B}}(xx^*) = 0$]. Thus, if \mathcal{A} and \mathcal{B} are subalgebras of the algebra of a noncommutative probability space (\mathcal{M}, φ) , there is a canonical algebra morphism from the algebra of the *t*-free product of $(\mathcal{A}, \varphi|_{\mathcal{A}})$ and $(\mathcal{B}, \varphi|_{\mathcal{B}})$ to \mathcal{M} whose restriction to $\mathcal{A} \cup \mathcal{B}$ preserves the expectation.

Now, we can give the definition of *t*-freeness. A real $t \in [0, +\infty]$ is still fixed.

DEFINITION 2.5 (*t*-freeness). Let (\mathcal{M}, τ) be a noncommutative probability space.

- Two subalgebras \mathcal{A}, \mathcal{B} of \mathcal{M} are said to be *t*-free if the canonical algebra morphism from the algebra of the *t*-free product of $(\mathcal{A}, \varphi_{\mathcal{A}})$ and $(\mathcal{B}, \varphi_{\mathcal{B}})$ to \mathcal{M} mentioned in Remark 2.4 preserves the expectation.

- Two subsets X, Y of \mathcal{M} are said to be t -free if the subalgebras they generate are t -free.

REMARK 2.6. The notion of t -freeness is a generalization of the notions of independence and freeness. Moreover, it appears already implicitly in some works of Voiculescu about the analogues of entropy and Fisher’s information measure in free probability theory:

- For $t = 0$, t -freeness is simply the independence, whereas it follows from [8] that in the case where $t = +\infty$, it is the freeness.
- If \mathcal{A}, \mathcal{B} are independent, then for all $t \geq 0$, their images at time t by the so-called *liberation process starting at* $(\mathcal{A}, \mathcal{B})$, defined by Voiculescu in [12], Section 2.1, are t -free.

The following proposition is obvious from the definition of t -freeness.

PROPOSITION 2.7. *Let (\mathcal{M}, τ) be a noncommutative probability space. Let $\{a_1, \dots, a_n\}$ and $\{b_1, \dots, b_m\}$ be two t -free subsets of \mathcal{M} . Then the joint noncommutative distribution of the family $(a_1, \dots, a_n, b_1, \dots, b_m)$ depends only on t and on the distributions of the families (a_1, \dots, a_n) and (b_1, \dots, b_m) .*

PROPOSITION–DEFINITION 2.8 (Additive and multiplicative t -free convolutions). *Let us fix $t \in [0, +\infty)$. Let μ, ν be compactly supported probability measures on the real line (resp., on $[0, +\infty)$, on the unit circle). Let a, b are t -free self-adjoint elements (resp., positive elements, unitary elements) with distributions μ, ν . Then the distribution of $a + b$ (resp., of $\sqrt{b}a\sqrt{b}$, of ab) is a compactly supported probability measure on the real line which depends only on t, μ and ν , and which will be denoted by $\mu *_t \nu$ (resp., $\mu \odot_t \nu$).*

PROOF. Let us treat the case of the sum of two self-adjoint elements. The other cases can be treated analogously. From Proposition 2.7, it follows that the moments of $a + b$ depend only on μ and ν . To see that these are the moments of a compactly supported probability measure on the real line, introduce $M > 0$ such that the supports of μ and ν are both contained in $[-M, M]$. Then for all $n \geq 1$, by Hölder inequalities in a noncommutative probability space [9], $\varphi((a + b)^{2n}) \leq 2^{2n} M^{2n}$. By the remark made after Definition 1.2, the result follows. \square

PROPOSITION 2.9 (Matricial model for the t -freeness). *For each $n \geq 1$, let M_n and N_n be diagonal random matrices whose noncommutative distributions have limits. Let also, for each n , S_n be the matrix of a uniform random permutation of $\{1, \dots, n\}$ and $U_{n,t}$ be a random $n \times n$ unitary matrix distributed according to the law of a Brownian motion on the unitary group at time t . Suppose that for each*

n , the sets of random variables $\{M_n, N_n\}, \{S_n\}, \{U_{n,t}\}$ are independent. Then as n tends to infinity, the noncommutative distribution of

$$(M_n, U_{n,t} S_n N_n S_n^* U_{n,t}^*)$$

converges in probability to that of a pair (a, b) of self-adjoint elements of a noncommutative probability space which are t -free.

PROOF. By Theorem 1.8, the noncommutative distribution of $(M_n, S_n N_n S_n^*)$ converges to the one of a pair (x, y) of independent elements. Moreover, since for all n , the law of U_n is invariant by conjugation, by Theorems 1.8 and 1.14, the family of sets

$$(\{M_n, S_n N_n S_n^*\}, \{U_{n,t}\})$$

is asymptotically free and the limit distribution of $U_{n,t}$ is that of a free unitary Brownian motion at time t . By definition of t -freeness, this concludes the proof. \square

In the next result, the convergences in probability of random measures toward nonrandom limits are understood with respect to the weak topology on the space of probability measures on the real line.

COROLLARY 2.10. For each n , let M_n, N_n be random $n \times n$ diagonal matrices, one of them having a distribution which is invariant under the action of the symmetric group by conjugation. Suppose that the spectral law of M_n (resp., N_n) converges in probability to some compactly supported probability measure μ (resp., ν) on the real line. Then the spectral law of $M_n + U_{n,t} N_n U_{n,t}^*$ converges in probability to the measure $\mu *_t \nu$.

3. Computation rules for t -freeness.

3.1. Multivariate free Itô calculus.

3.1.1. *Technical preliminaries.* In this section, we shall extend some results of [7] to the multivariate case. Let us first recall basics of free stochastic calculus. For more involved definitions, the reader should refer to Sections 1 and 2 of [7]. Let (\mathcal{M}, τ) be a faithful¹ noncommutative probability space endowed with a filtration $(\mathcal{M}_t)_{t \geq 0}$ and an $(\mathcal{M}_t)_{t \geq 0}$ -free additive Brownian motion $(X_t)_{t \geq 0}$. Let \mathcal{M}^{op} be the opposite algebra of \mathcal{M} (it is the same vector space, but it is endowed with the product $a \times_{op} b = ba$). We shall denote by \sharp the left actions of the algebra $\mathcal{M} \otimes$

¹A noncommutative probability space (\mathcal{M}, τ) is said to be *faithful* if for all x in $\mathcal{M} \setminus \{0\}$, $\tau(xx^*) > 0$. Any noncommutative probability space can be quotiented by a bilateral ideal into a faithful space.

\mathcal{M}^{op} on \mathcal{M} and $\mathcal{M} \otimes \mathcal{M}$ defined by $(a \otimes b)\sharp u = au b$ and $(a \otimes b)\sharp(u \otimes v) = au \otimes vb$. The algebras \mathcal{M} and $\mathcal{M} \otimes \mathcal{M}^{op}$ are endowed with the inner products defined by $\langle a, b \rangle = \tau(ab^*)$ and $\langle a \otimes b, c \otimes d \rangle = \tau(ac^*)\tau(bd^*)$. The Riemann integral of functions defined on a closed interval with left and right limits at any point with values in the Hilbert space² $L^2(\mathcal{M}, \tau)$ is a well-known notion. Now, we shall recall the definition of the stochastic integral. A simple adapted biprocess is a piecewise constant map U from $[0, +\infty)$ to $\mathcal{M} \otimes \mathcal{M}^{op}$ vanishing for t large enough such that $U_t \in \mathcal{M}_t \otimes \mathcal{M}_t$ for all t . The set of simple biprocesses is endowed with the inner product

$$\langle U, V \rangle = \int_0^\infty \langle U_t, V_t \rangle dt.$$

We shall denote by \mathcal{B}_2^a the closure of the set of simple adapted biprocesses with respect to this inner product. Let U be a simple adapted biprocess. Then there exists times $0 = t_0 \leq t_1 \leq \dots \leq t_m$ such that L (resp., U) is constant on each $[t_i, t_{i+1})$ and vanishes on $[t_m, +\infty)$. Then we define

$$\int_0^\infty U_t dX_t = \sum_{i=0}^{m-1} U_{t_i} \sharp (X_{t_{i+1}} - X_{t_i}).$$

It can be proved (Corollary 3.1.2 of [7]) that the map $U \mapsto \int_0^\infty U_t dX_t$ can be extended isometrically from \mathcal{B}_2^a to $L^2(\mathcal{M}, \tau)$.

3.1.2. *Free Itô processes.* We shall call a *free Itô process* any process

$$(3) \quad A_t = A_0 + \int_0^t L_s ds + \int_0^t U_s dX_s,$$

where $A_0 \in \mathcal{M}_0$, L is an adapted process with left and right limit at any point and $U \in \mathcal{B}_2^a$. In this case, we shall denote

$$(4) \quad dA_t = L_t dt + U_t \sharp dX_t.$$

The part $U_t \sharp dX_t$ of this expression is called the *martingale part* of A . Note that the process A is determined by A_0 and dA_t .

We shall use the following lemma, which follows from Proposition 2.2.2 of [7] and from the linearity of τ .

LEMMA 3.1. *Let A_t be as in (3). Then $\tau(A_t) = \tau(A_0) + \int_0^t \tau(L_s) ds$.*

²The Hilbert space $L^2(\mathcal{M}, \tau)$ is the completion of \mathcal{M} with respect to the inner product $\langle x, y \rangle = \tau(xy^*)$.

3.1.3. *Multivariate free Itô calculus.* Consider n elements $a_1, \dots, a_n \in \mathcal{M}$ for some $n \geq 2$. Consider also two elements $u = \sum_k x_k \otimes y_k, v = \sum_l z_l \otimes t_l$ of $\mathcal{M} \otimes \mathcal{M}^{op}$. For all $1 \leq i < j \leq n$, we define an element of \mathcal{M} by setting

$$\begin{aligned} &\langle\langle a_1, \dots, a_{i-1}, u, a_{i+1}, \dots, a_{j-1}, v, a_{j+1}, \dots, a_n \rangle\rangle_{i,j} \\ &= \sum_{k,l} a_1 \cdots a_{i-1} x_k \tau(y_k a_{i+1} \cdots a_{j-1} z_l) t_l a_{j+1} \cdots a_n. \end{aligned}$$

The following theorem follows from Theorem 4.1.12 and the remark following in [7].

THEOREM 3.2. *Let $A_t = A_0 + \int_0^t L_s ds + \int_0^t U_s dX_s$ and $B_t = B_0 + \int_0^t K_s ds + \int_0^t V_s dX_s$ be two Itô processes with respect to the same free Brownian motion (X_t) . Then AB is a free Itô process and with the notations of (4),*

$$d(AB)_t = A_t dB_t + (dA_t)B_t + \langle\langle U_t, V_t \rangle\rangle_{1,2} dt.$$

In order to prove computation rules for t -freeness, we shall need the following theorem.

THEOREM 3.3. *Let A_1, \dots, A_n be free Itô processes with respect to the same Brownian motion. For all k , denote $A_{k,t} = A_{k,0} + \int_0^t L_{k,s} ds + \int_0^t U_{k,s} dX_s$. Then $A_1 \cdots A_n$ is a free Itô process and*

$$\begin{aligned} d(A_1 \cdots A_n)_t &= \sum_{k=1}^n A_{1,t} \cdots A_{k-1,t} (dA_{k,t}) A_{k+1,t} \cdots A_{n,t} \\ &+ \sum_{1 \leq k < l \leq n} \langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, \\ &A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{n,t} \rangle\rangle_{k,l} dt. \end{aligned}$$

PROOF. Let us prove this theorem by induction on n . For $n = 1$, it is obvious. Let us suppose the result to hold at rank n . Then the martingale part of $A_1 \cdots A_n$ is

$$\sum_{k=1}^n A_{1,t} \cdots A_{k-1,t} (U_{k,t} \sharp dX_{k,t}) A_{k+1,t} \cdots A_{n,t}.$$

Thus, by Theorem 3.2, $A_1 \cdots A_{n+1}$ is a free Itô process and

$$\begin{aligned} &d(A_1 \cdots A_{n+1})_t \\ &= (A_1 \cdots A_n)_t dA_{n+1,t} + (d(A_1 \cdots A_n)_t) A_{n+1,t} \\ &+ \sum_{k=1}^n \langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, A_{k+1,t}, \dots, A_{n,t}, U_{n+1,t} \rangle\rangle_{k,n+1} dt \end{aligned}$$

$$\begin{aligned}
 &= \sum_{k=1}^{n+1} A_{1,t} \cdots A_{k-1,t} (dA_{k,t}) A_{k+1,t} \cdots A_{n+1,t} \\
 &\quad + \sum_{1 \leq k < l \leq n} \langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, \\
 &\qquad\qquad\qquad A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{n,t}, A_{n+1,t} \rangle\rangle_{k,l} dt \\
 &\quad + \sum_{k=1}^n \langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, A_{k+1,t}, \dots, A_{n,t}, U_{n+1,t} \rangle\rangle_{k,n+1} dt,
 \end{aligned}$$

which concludes the proof. \square

3.2. Computation rules for t -freeness.

3.2.1. *Main result.* In order to do computations with elements which are t -free, we have to find out a formula for the expectation of a product of elements of the type

$$(5) \quad x_1 u_t y_1 u_t^* x_2 u_t y_2 u_t^* \cdots x_n u_t y_n u_t^*$$

for $\{x_1, \dots, x_n\}$ independent with $\{y_1, \dots, y_n\}$ and $\{x_1, y_1, \dots, x_n, y_n\}$ free with u_t , free unitary Brownian motion. Actually, for the result which follows, the independence of the x_i 's and the y_i 's will not be useful, thus we consider a noncommutative probability space (\mathcal{M}, τ) , an integer $n \geq 1$, $a_1, \dots, a_{2n} \in \mathcal{M}$ and a free unitary Brownian motion (u_t) which is free with $\{a_1, \dots, a_{2n}\}$. In order to have some more concise formulae, it will be useful to multiply the product of (5) by e^{nt} . So we define

$$f_{2n}(a_1, \dots, a_{2n}, t) = e^{nt} \tau(a_1 u_t a_2 u_t^* \cdots a_{2n-1} u_t a_{2n} u_t^*).$$

We shall use the convention $f_0(a, t) = \tau(a)$ for all $a \in \mathcal{M}$.

Since $f_{2n}(a_1, \dots, a_{2n}, 0) = \tau(a_1 \cdots a_{2n})$, the following theorem allows us to deduce all functions $f_{2n}(a_1, \dots, a_{2n}, t)$ [thus, the expectation of any product of the type of (5)] from the joint distribution of the a_i 's.

THEOREM 3.4. *For all $n \geq 1$ and all $a_1, \dots, a_{2n} \in \mathcal{M}$ free with the process (u_t) , the following differential relations are satisfied:*

$$\begin{aligned}
 &\frac{\partial}{\partial t} f_{2n}(a_1, \dots, a_{2n}, t) \\
 &= - \sum_{\substack{1 \leq k < l \leq 2n \\ k=l \bmod 2}} f_{2n-(l-k)}(a_1, \dots, a_k, a_{l+1}, \dots, a_{2n}, t) f_{l-k}(a_{k+1}, \dots, a_l, t) \\
 &\quad + e^t \sum_{\substack{1 \leq k < l \leq 2n \\ k \neq l \bmod 2}} f_{2n-(l-k)-1}(a_1, \dots, a_{k-1}, a_k a_{l+1}, a_{l+2}, \dots, a_{2n}, t) \\
 &\qquad\qquad\qquad \times f_{l-k-1}(a_l a_{k+1}, a_{k+2}, \dots, a_{l-1}, t).
 \end{aligned}$$

PROOF. Let us introduce the process (v_t) defined by $v_t = e^{t/2}u_t$ for all t . As explained in the beginning of Section 2.3 of [4], this process can be realized as an Itô process, with the formula

$$v_t = 1 + i \int_0^t dX_s v_s.$$

Thus, one can realize the family of noncommutative random variables $a_1, \dots, a_{2n}, (v_t)_{t \geq 0}$ in a faithful noncommutative probability space (\mathcal{M}, τ) endowed with a filtration $(\mathcal{M}_t)_{t \geq 0}$ and an additive free Brownian motion $(X_t)_{t \geq 0}$ such that $a_1, \dots, a_{2n} \in \mathcal{M}_0$ and for all $t, v_t = 1 + i \int_0^t dX_s v_s$ and $v_t^* = 1 - i \int_0^t v_s^* dX_s$. By definition of $f_{2n}(a_1, \dots, a_{2n}, t)$, one has

$$f_{2n}(a_1, \dots, a_{2n}, t) = \tau(a_1 v_t a_2 v_t^* \cdots a_{2n-1} v_t a_{2n} v_t^*).$$

Note that since all a_i 's belong to \mathcal{M}_0 , the processes $A_1 := (a_1 v_t), A_2 := (a_2 v_t^*), \dots, A_{2n-1} := (a_{2n-1} v_t), A_{2n} := (a_{2n} v_t^*)$ are all free Itô processes: if one defines $U_{k,t} = a_k \otimes i v_t$ for k odd and $U_{k,t} = -i a_k v_t^* \otimes 1$ for k even, then for all $k, dA_{k,t} = U_{k,t} \sharp dX_t$. Thus, by Theorem 3.3, $A_1 \cdots A_{2n}$ is an Itô process such that for all $t,$

$$\begin{aligned} (A_1 \cdots A_{2n})_t &= (A_1 \cdots A_{2n})_0 + \int_0^t \sum_{k=1}^{2n} A_{1,s} \cdots A_{k-1,s} (U_{k,s} \sharp dX_s) A_{k+1,s} \cdots A_{2n,s} \\ &\quad + \int_0^t \sum_{1 \leq k < l \leq 2n} \langle\langle A_{1,s}, \dots, A_{k-1,s}, U_{k,s}, A_{k+1,s}, \dots, A_{l-1,s}, \\ &\hspace{15em} U_{l,s}, A_{l+1,s}, \dots, A_{2n,s} \rangle\rangle_{k,l} ds. \end{aligned}$$

Hence, by Lemma 3.1, for all $t,$

$$\begin{aligned} &\frac{\partial}{\partial t} f_{2n}(a_1, \dots, a_{2n}, t) \\ (6) \quad &= \sum_{1 \leq k < l \leq 2n} \tau(\langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, \\ &\hspace{15em} A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{2n,t} \rangle\rangle_{k,l}). \end{aligned}$$

Now, fix $1 \leq k < l \leq 2n$ and discuss according to the parity of k and l .

- If $k = l \pmod 2$. Suppose, for example, that k, l are both odd (the other case can be treated in the same way). Then $U_{k,t} = a_k \otimes i v_t$ and $U_{l,t} = a_l \otimes i v_t$, which implies that

$$\begin{aligned} &\tau(\langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{2n,t} \rangle\rangle_{k,l}) \\ &= i \tau(a_1 v_t a_2 v_t^* \cdots a_{k-1} v_t^* a_k v_t a_{l+1} v_t \cdots a_{2n} v_t^*) i \tau(v_t a_{k+1} v_t^* \cdots a_{l-1} v_t^* a_l). \end{aligned}$$

Note that since τ is tracial and the joint distribution of $a_1, \dots, a_{2n}, (v_t)_{t \geq 0}$ is the same as the one of $a_1, \dots, a_{2n}, (v_t^*)_{t \geq 0}$, we have $\tau(v_t a_{k+1} v_t^* \cdots a_{l-1} v_t^* a_l) =$

$\tau(a_{k+1}v_t \cdots a_{l-1}v_t a_l v_t^*)$. Hence,

$$(7) \quad \begin{aligned} &\tau(\langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{2n,t} \rangle\rangle_{k,l}) \\ &= -f_{2n-(l-k)}(a_1, \dots, a_k, a_{l+1}, \dots, a_{2n}, t) f_{l-k}(a_{k+1}, \dots, a_l, t). \end{aligned}$$

• If $k \neq l \pmod 2$. Suppose, for example, k to be odd and l to be even (the other case can be treated in the same way). Then $U_{k,t} = a_k \otimes i v_t$ and $U_{l,t} = -a_l i v_t^* \otimes 1$, which implies that

$$\begin{aligned} &\tau(\langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{2n,t} \rangle\rangle_{k,l}) \\ &= \tau(a_1 v_t a_2 v_t^* \cdots a_{k-1} v_t^* a_k a_{l+1} v_t \cdots a_{2n} v_t^*) (-i^2) \tau(v_t a_{k+1} v_t^* \cdots a_{l-1} v_t a_l v_t^*). \end{aligned}$$

Note that since $v_t^* v_t = e^t$, τ is tracial and the joint distribution of $a_1, \dots, a_{2n}, (v_t)_{t \geq 0}$ is the same as the one of $a_1, \dots, a_{2n}, (v_t^*)_{t \geq 0}$, we have $\tau(v_t a_{k+1} v_t^* \cdots \times a_{l-1} v_t a_l v_t^*) = e^t \tau(a_l a_{k+1} v_t \cdots a_{l-1} v_t^*)$. Hence,

$$(8) \quad \begin{aligned} &\tau(\langle\langle A_{1,t}, \dots, A_{k-1,t}, U_{k,t}, A_{k+1,t}, \dots, A_{l-1,t}, U_{l,t}, A_{l+1,t}, \dots, A_{2n,t} \rangle\rangle_{k,l}) \\ &= e^t f_{2n-(l-k)-1}(a_1, \dots, a_{k-1}, a_k a_{l+1}, a_{l+2}, \dots, a_{2n}, t) \\ &\quad \times f_{l-k-1}(a_l a_{k+1}, a_{k+2}, \dots, a_{l-1}, t). \end{aligned}$$

Equations (6)–(8) together conclude the proof. \square

The following proposition, which we shall use later, is a consequence of the previous theorem.

PROPOSITION 3.5. *In a noncommutative probability space (\mathcal{M}, τ) , consider two independent normal elements a, b with symmetric compactly supported laws. Let (u_t) be a free unitary Brownian motion which is free with $\{a, b\}$. Then the function*

$$G(t, z) = \sum_{n \geq 1} \tau((a u_t b u_t^*)^{2n}) e^{2nt} z^n$$

is the only solution, in a neighborhood of $(0, 0)$ in $[0, +\infty) \times \mathbb{C}$, to the nonlinear, first-order partial differential equation

$$(9) \quad \frac{\partial G}{\partial t} + 4zG \frac{\partial G}{\partial z} = 0,$$

$$(10) \quad G(0, z) = \sum_{n \geq 1} \tau(a^{2n}) \tau(b^{2n}) z^n.$$

PROOF. Let us define, for all $n \geq 1$, $g_n(t) = \tau((a u_t b u_t^*)^n) e^{nt}$. For $n = 0$, we set $g_0(t) = 0$. Let us fix $n \geq 1$. In order to apply the previous theorem, let us define,

for $i = 1, \dots, 2n$, $a_i = a$ if i is odd and $a_i = b$ if i is even. By the previous theorem, for all $n \geq 1$, we have

$$\begin{aligned}
 \frac{\partial}{\partial t} g_n(t) &= - \sum_{\substack{1 \leq k < l \leq 2n \\ k=l \bmod 2}} g_{n-(l-k)/2}(t) g_{(l-k)/2}(t) \\
 (11) \quad &+ e^t \sum_{\substack{1 \leq k < l \leq 2n \\ k \neq l \bmod 2}} f_{2n-(l-k)-1}(a_1, \dots, a_{k-1}, a_k a_{l+1}, a_{l+2}, \dots, a_{2n}, t) \\
 &\quad \times f_{l-k-1}(a_l a_{k+1}, a_{k+2}, \dots, a_{l-1}, t).
 \end{aligned}$$

Now, note that since for any $\varepsilon, \varepsilon' = \pm 1$, the joint distribution of (a, b, u_t) is the same as the one of $(\varepsilon a, \varepsilon' b, u_t)$, $g_p(t) = 0$ when p is odd. Thus, in the first sum of (11) only pairs (k, l) such that $k = l \bmod 4$ have a nonnull contribution. For the same reason, all terms in the second sum are null. Indeed, for any $1 \leq k < l \leq 2n$ such that $k \neq l \bmod 2$, the set $\{k + 1, k + 2, \dots, l\}$, whose cardinality is odd, has either an odd number of odd elements or an odd number of even elements. To sum up, for all $n \geq 1$, we have

$$\begin{aligned}
 \frac{\partial}{\partial t} g_{2n}(t) &= - \sum_{\substack{1 \leq k < l \leq 4n \\ k=l \bmod 4}} g_{2n-(l-k)/2}(t) g_{(l-k)/2}(t) = -4 \sum_{i=1}^{n-1} (n-i) g_{2(n-i)}(t) g_{2i}(t) \\
 &= -2n \sum_{i=1}^{n-1} g_{2(n-i)}(t) g_{2i}(t).
 \end{aligned}$$

Thus, since $g_0(t) = 0$ and $G(t, z) = \sum_{n \geq 1} g_{2n}(t) z^n = \sum_{n \geq 0} g_{2n}(t) z^n$, the last computation implies

$$\frac{\partial G}{\partial t} = -2z \frac{\partial G^2}{\partial z},$$

which proves (9). The formula (10) is obvious.

To prove the uniqueness, let $H(t, z) = \sum_{n \geq 0} h_n(t) z^n$ be another solution of (9) and (10). By (10), for all $n \geq 0$, we have $h_n(0) = g_{2n}(0)$ and by (9), for all $n \geq 0$, we have $\frac{\partial}{\partial t} h_n(t) = -2n \sum_{m=0}^n h_{n-m}(t) h_m(t)$, which implies that $h_0 = 0$ and that by induction on n , $h_n = g_{2n}$. \square

3.2.2. Examples. Let us give examples of applications of the computation rules that we have just established. The third example below is a rather big formula, but we shall need it when we study the problem of existence of t -free cumulants. So, let \mathcal{A} and \mathcal{B} be two independent subalgebras of a noncommutative probability space (\mathcal{M}, τ) and (u_t) be a free unitary Brownian motion free from $\mathcal{A} \cup \mathcal{B}$.

(1) For $a \in \mathcal{A}, b \in \mathcal{B}$, for all $t \geq 0$, we have

$$(12) \quad \tau(a u_t b u_t^*) = \tau(a) \tau(b).$$

[In fact, it even follows from Theorem 3.4 that without the assumption that a and b are independent, for all t , we have $\tau(au_tbu_t^*) = e^{-t}\tau(ab) + (1 - e^{-t})\tau(a)\tau(b)$.]

(2) For $a, a' \in \mathcal{A}, b, b' \in \mathcal{B}$, for all $t \geq 0$, we have

$$\begin{aligned} & \tau(au_tbu_t^*a'u_tb'u_t^*) \\ (13) \quad & = (\tau(a)\tau(a')\tau(bb') + \tau(aa')\tau(b)\tau(b') - \tau(a)\tau(a')\tau(b)\tau(b'))(1 - e^{-2t}) \\ & \quad + \tau(aa')\tau(bb')e^{-2t}. \end{aligned}$$

(3) For $a, a', a'' \in \mathcal{A}$ and $b, b', b'' \in \mathcal{B}$, we have

$$\begin{aligned} & \tau(au_tbu_t^*a'u_tb'u_t^*a''u_tb''u_t^*) \\ & = \tau(a)\tau(a')\tau(a'')\tau(b)\tau(b')\tau(b'')(2 - 6e^{-2t} + 4e^{-3t}) \\ & \quad - (1 - 3e^{-2t} + 2e^{-3t})\tau(a)\tau(a')\tau(a'')\tau(b)\tau(b'b'') \\ & \quad - (1 - 3e^{-2t} + 2e^{-3t})\tau(a)\tau(a')\tau(a'')\tau(bb')\tau(b'') \\ & \quad - (1 - 3e^{-2t} + 2e^{-3t})\tau(a)\tau(a')\tau(a'')\tau(bb'')\tau(b') \\ & \quad - (1 - 3e^{-2t} + 2e^{-3t})\tau(aa')\tau(a'')\tau(b)\tau(b')\tau(b'') \\ & \quad - (1 - 3e^{-2t} + 2e^{-3t})\tau(aa'')\tau(a')\tau(b)\tau(b')\tau(b'') \\ & \quad - (1 - 3e^{-2t} + 2e^{-3t})\tau(a)\tau(a'a'')\tau(b)\tau(b')\tau(b'') \\ & \quad + (1 - 3e^{-2t} + 2e^{-3t})[\tau(a)\tau(a')\tau(a'')\tau(bb'b'')] \\ (14) \quad & \quad + \tau(aa'a'')\tau(b)\tau(b')\tau(b'')] \\ & \quad - (e^{-2t} - e^{-3t})[\tau(aa')\tau(a'')\tau(bb')\tau(b'') + \tau(aa')\tau(a'')\tau(bb'')\tau(b') \\ & \quad \quad + \tau(aa'')\tau(a')\tau(bb'')\tau(b') + \tau(aa'')\tau(a')\tau(b)\tau(b'b'') \\ & \quad \quad + \tau(a)\tau(a'a'')\tau(b'')\tau(bb') + \tau(a)\tau(a'a'')\tau(b)\tau(b'b'')] \\ & \quad + (1 - 2e^{-2t} + e^{-3t})[\tau(aa')\tau(a'')\tau(b'b'')\tau(b) \\ & \quad \quad + \tau(aa'')\tau(a')\tau(bb')\tau(b'') \\ & \quad \quad + \tau(a)\tau(a'a'')\tau(bb'')\tau(b')] \\ & \quad + (e^{-2t} - e^{-3t})\tau(bb'b'')[\tau(a)\tau(a'a'') + \tau(aa')\tau(a'') + \tau(aa'')\tau(a')] \\ & \quad + (e^{-2t} - e^{-3t})\tau(aa'a'')[\tau(b)\tau(b'b'') + \tau(bb')\tau(b'') + \tau(bb'')\tau(b')] \\ & \quad + e^{-3t}\tau(aa'a'')\tau(bb'b''). \end{aligned}$$

It can be verified that the last formula actually corresponds to the formula of $\mathbb{E}(aba'b'a''b'')$ with $\{a, a', a''\}$ and $\{b, b', b''\}$ independent when $t = 0$, and to the formula of $\tau(aba'b'a''b'')$ with $\{a, a', a''\}$ and $\{b, b', b''\}$ free when t tends to infinity.

3.3. *Multiplicative and additive t -free convolutions of two symmetric Bernoulli laws.*

3.3.1. *Multiplicative case.* Here, we shall compute the multiplicative t -free convolution of $\frac{\delta_{-1+\delta_1}}{2}$ (considered as a law on the unit circle) with itself.

THEOREM 3.6. *For all $t \geq 0$, $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ is the only law on the unit circle which is invariant under the symmetries with respect to the real and imaginary axes and whose push-forward by the map $z \mapsto z^2$ has the law of u_{4t} , a free unitary Brownian motion taken at time $4t$.*

REMARK 3.7. The moments of u_{4t} have been computed by Biane at Lemma 1 of [4]: for all $n \geq 1$,

$$(15) \quad \tau(u_{4t}^n) = \frac{e^{-2nt}}{n} \sum_{k=0}^{n-1} \frac{(-4nt)^k}{k!} \binom{n}{k+1}.$$

PROOF OF THEOREM 3.6. In a noncommutative probability space (\mathcal{M}, τ) , consider two independent normal elements a, b with law $\frac{\delta_{-1+\delta_1}}{2}$. Let (u_t) be a free unitary Brownian motion which is free with $\{a, b\}$. Then $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ is the distribution of the unitary element $au_tbu_t^*$. Since the joint distribution of (a, b, u_t) is the same as the one of $(-a, b, u_t)$, $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ is invariant under the transformation $z \mapsto -z$. Moreover, $(au_tbu_t^*)^* = u_tbu_t^*a$ has the same distribution as $au_tbu_t^*$ (because τ is tracial and u_t has the same law as u_t^*), hence $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ is invariant under the transformation $z \mapsto \bar{z}$. This proves that $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ is invariant under the symmetries with respect to the real and imaginary axes.

Since any distribution on the unit circle is determined by its moments, to prove that the push-forward of $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ by the map $z \mapsto z^2$ is the law of u_{4t} , it suffices to prove that for all $n \geq 1$,

$$\tau((au_tbu_t^*)^{2n}) = \tau(u_{4t}^n),$$

that is, to prove that the functions

$$F_1(t, z) = \sum_{n \geq 1} \tau((au_tbu_t^*)^{2n}) e^{2nt} z^n \quad \text{and} \quad F_2(t, z) = \sum_{n \geq 1} \tau(u_{4t}^n) e^{2nt} z^n$$

are equal. It follows from Proposition 3.5 that F_1 is the only solution, in a neighborhood of $(0, 0)$ in $[0, +\infty) \times \mathbb{C}$, to equation (9) satisfying $F_1(0, z) = \frac{z}{1-z}$. But it follows from Lemma 1 of [4] that F_2 is also a solution of (9) with the same initial conditions. By uniqueness, it closes the proof. \square

For all $t \in [0, 1]$, let us define $\beta(t) = 2\sqrt{t(1-t)} + \arccos(1-2t)$. Then $\beta(t)$ is an increasing function of t which goes from 0 to π when t goes from 0 to 1. Biane

has proved in [5], Proposition 10, that the distribution of u_{A_t} is absolutely continuous with respect to the Lebesgue measure on the unit circle, that its support is the full unit circle for $t \geq 1$, and the set $\{e^{i\theta} : |\theta| \leq \beta(t)\}$ for all $t \in [0, 1]$. Moreover, the density of this distribution with respect to the uniform probability measure on the unit circle, which we denote by ρ_{A_t} , is positive and analytic on the interior of its support for all $t \geq 0$, except at -1 for $t = 1$.

REMARK 3.8. Since there is no simple formula for the density of u_t , it may be worth explaining how we were able to produce the picture shown in Figure 2. The expression of the moments of the distribution of u_t given by (15) is numerically ineffective, because it is an alternated sum of very large numbers. It only allows one to compute the first few dozens of moments of the distribution. Nevertheless, this expression of the moments allows one to prove that, for all $t \geq 0$, the function κ_t defined on the interior of the complex unit disk by the formula

$$\kappa_t(z) = \int_0^{2\pi} \frac{e^{i\theta} + z}{e^{i\theta} - z} \rho_t(e^{i\theta}) \frac{d\theta}{2\pi} = 1 + 2 \sum_{k=1}^{+\infty} \tau(u_t^k) z^k,$$

satisfies the equation

$$(16) \quad \frac{\kappa_t(z) - 1}{\kappa_t(z) + 1} e^{(t/2)\kappa_t(z)} = z.$$

This fact can be established using the Lagrange inversion formula (see [4]), see also [5], Section 4.2.2. Now, on one hand, a computer seems to be able to solve

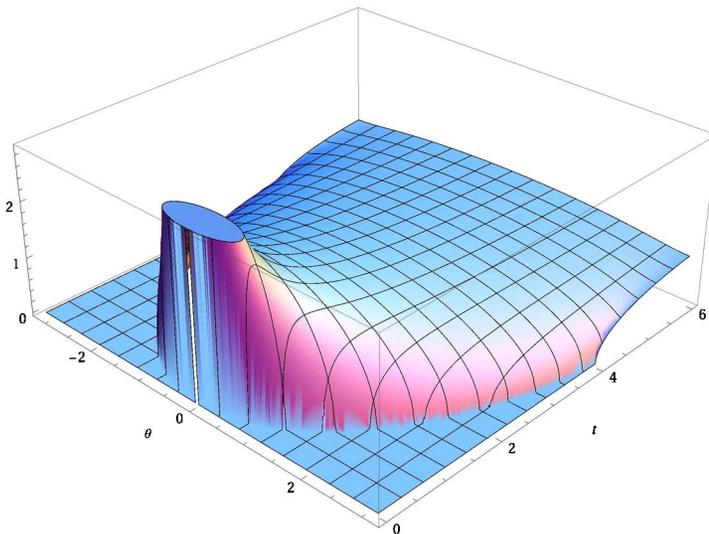


FIG. 2. The density of the distribution of u_t at the point $e^{i\theta}$ in function of θ and t . One sees the support progressively filling the circle when t increases from 0 to 4, and then the distribution rapidly converging toward the uniform distribution.

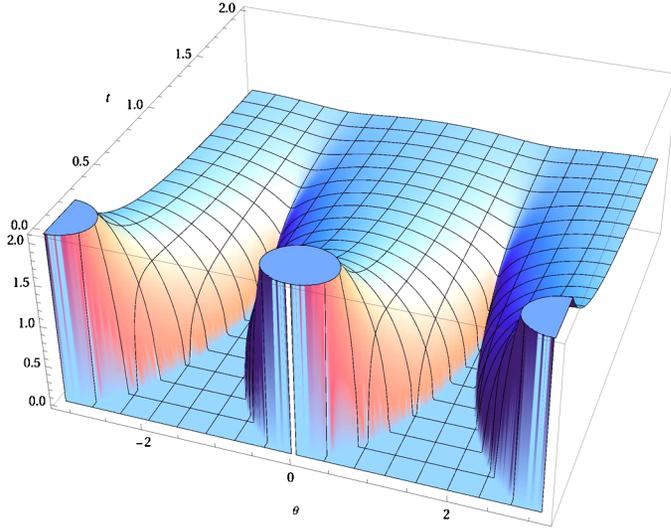


FIG. 3. The density of $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ at the point $e^{i\theta}$ in function of θ and t . The support progressively fills the circle when t increases from 0 to 1, and then the distribution converges rapidly toward the uniform distribution.

this equation more reliably than it computes the moments of the distribution. On the other hand, κ_t is the holomorphic function in the unit disk whose real part is the harmonic extension of the density of the distribution of u_t . Thus, we evaluated $\rho_t(e^{i\theta})$ by taking the real part of a numerical solution of (16) with $z = e^{i\theta}$.

From the facts exposed above Remark 3.8, one deduces easily the next result, illustrated by Figure 3.

COROLLARY 3.9. For all $t > 0$, the measure $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$ has a density with respect to the uniform probability measure on the unit circle, which we shall denote by σ_t and which is given by the formula $\sigma_t(z) = \rho_{4t}(z^2)$ for all z in the unit circle. In particular, the support of this measure is the full unit circle for $t \geq 1$ and the set $\{e^{i\theta} : |\theta| \leq \frac{1}{2}\beta(t) \text{ or } |\pi - \theta| \leq \frac{1}{2}\beta(t)\}$ for $t \in [0, 1]$. The density σ_t is positive and analytic on the interior of its support for all $t \geq 0$, except at $\pm i$ for $t = 1$.

3.3.2. Additive case. Here, we shall compute the additive t -free convolution of $\frac{\delta_{-1+\delta_1}}{2}$ (considered as a law on the real line) with itself (see also Figure 4).

THEOREM 3.10. For all $t \geq 0$, $\frac{\delta_{-1+\delta_1}}{2} *_t \frac{\delta_{-1+\delta_1}}{2}$ is the only symmetric law on the real line whose push-forward by the map $x \mapsto x^2$ has the law of $2 + v + v^*$, with v unitary element distributed according to $\frac{\delta_{-1+\delta_1}}{2} \odot_t \frac{\delta_{-1+\delta_1}}{2}$.

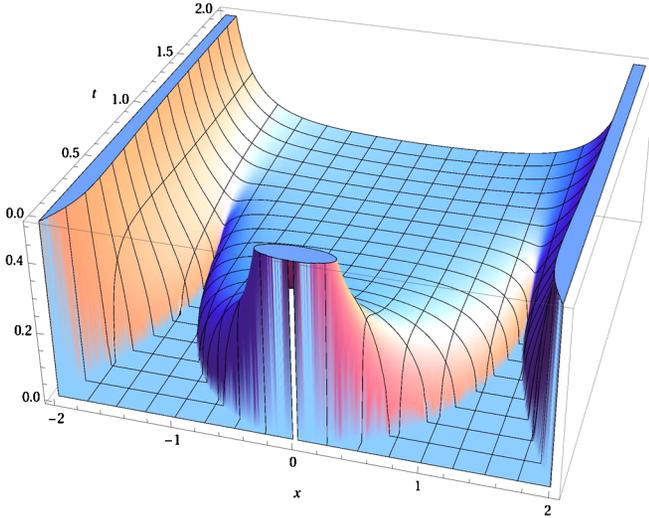


FIG. 4. The density of $\frac{\delta_{-1+\delta_1}}{2} *_t \frac{\delta_{-1+\delta_1}}{2}$ at the point x in function of x and t . The support fills progressively the interval $[-2, 2]$ when t increases from 0 to 1, and then the distribution converges rapidly toward the arcsine distribution.

PROOF. In a noncommutative probability space (\mathcal{M}, τ) , consider two independent normal elements a, b with law $\frac{\delta_{-1+\delta_1}}{2}$. Let (u_t) be a free unitary Brownian motion which is free with $\{a, b\}$. Then $\frac{\delta_{-1+\delta_1}}{2} *_t \frac{\delta_{-1+\delta_1}}{2}$ is the distribution of $a + u_t b u_t^*$. Since the joint distribution of (a, b, u_t) is the same as the one of $(-a, -b, u_t)$, $\frac{\delta_{-1+\delta_1}}{2} *_t \frac{\delta_{-1+\delta_1}}{2}$ is symmetric. Note that since a^2 and b^2 have δ_1 for distribution, one can suppose that $a^2 = b^2 = 1$. In this case,

$$(a + u_t b u_t^*)^2 = 2 + a u_t b u_t^* + u_t b u_t^* a = 2 + a u_t b u_t^* + (a u_t b u_t^*)^*$$

and the result is obvious by definition of \odot_t . \square

From the last result and Corollary 3.9, we deduce the following.

COROLLARY 3.11. For all $t > 0$, the measure $\frac{\delta_{-1+\delta_1}}{2} *_t \frac{\delta_{-1+\delta_1}}{2}$ has a density with respect to the Lebesgue measure on $[-2, 2]$, which we shall denote by η_t and which is given by the formula

$$\forall x \in [-2, 2] \quad \eta_t(x) = \rho_{4t} (e^{4i \arccos(x/2)}) \frac{1}{\pi \sqrt{4 - x^2}}.$$

The support of this measure is the interval $[-2, 2]$ for $t \geq 1$, and the set

$$\left[-2, -2 \cos \frac{\beta(t)}{4}\right] \cup \left[-2 \sin \frac{\beta(t)}{4}, 2 \sin \frac{\beta(t)}{4}\right] \cup \left[2 \cos \frac{\beta(t)}{4}, 2\right]$$

for $t \in [0, 1]$. The density η_t is positive and analytic on the interior of its support for all $t \geq 0$, except at $\pm\sqrt{2}$ for $t = 1$.

4. The lack of cumulants. In this section, we investigate the existence of an analogue of classical and free cumulants in the context of t -freeness. Informally, the problem is to find multilinear forms defined on any noncommutative probability space which vanish when evaluated on a family of elements which can be split into two nonempty subfamilies which are t -free.

More precisely, given a noncommutative probability space (\mathcal{M}, φ) , we would like to know if there exists a family $(k_n)_{n \geq 2}$ of multilinear forms on \mathcal{M} , with k_n an n -linear form for all $n \geq 2$, such that, for all $n \geq 2$, all $n_1, n_2 > 0$ such that $n_1 + n_2 = n$, all m_1, \dots, m_n in \mathcal{M} such that $\{m_1, \dots, m_{n_1}\}$ and $\{m_{n_1+1}, \dots, m_{n_1+n_2}\}$ are t -free, and finally for all $\sigma \in \mathfrak{S}_n$, one has $k_n(m_{\sigma(1)}, \dots, m_{\sigma(n)}) = 0$.

Our main result is negative: there does not exist in general such a family of multilinear forms, at least in a large class which we describe now.

DEFINITION 4.1. Let (\mathcal{M}, φ) be a noncommutative probability space. Let $n \geq 1$ be an integer. Let σ be an element of \mathfrak{S}_n . We define the n -linear form φ_σ on \mathcal{M} as follows:

$$\forall m_1, \dots, m_n \in \mathcal{M} \quad \varphi_\sigma(m_1, \dots, m_n) = \prod_{\substack{c \text{ cycle of } \sigma \\ c=(i_1 \dots i_r)}} \varphi(m_{i_1} \cdots m_{i_r}).$$

Using only the algebra structure of \mathcal{M} and the linear form φ , a linear combination of the forms $\{\varphi_\sigma : \sigma \in \mathfrak{S}_n\}$ seems to be the most general n -linear form that one can construct on \mathcal{M} . We seek cumulants within this wide class of n -linear forms. Our definition does not require that the vanishing of cumulants characterize t -freeness. We only insist that mixed cumulants of t -free variable vanish.

DEFINITION 4.2. Let $n \geq 2$ be an integer. Let $t \geq 0$ be a real number. A t -free cumulant of order n is a collection $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ of complex numbers such that $\sum_{\sigma \text{ n-cycle}} c(\sigma) \neq 0$ and the following property holds for every noncommutative probability space (\mathcal{M}, φ) : for any pair $(\mathcal{A}, \mathcal{B})$ of sub-algebras of \mathcal{M} which are t -free with respect to φ , for any family (m_1, \dots, m_n) of elements of $\mathcal{A} \cup \mathcal{B}$, which do not all belong to \mathcal{A} , and not all to \mathcal{B} , we have

$$(17) \quad \sum_{\sigma \in \mathfrak{S}_n} c(\sigma) \varphi_\sigma(m_1, \dots, m_n) = 0.$$

Let us emphasize that what we call *cumulant* is not a specific multilinear form, but rather a collection of coefficients which allows one to define a multilinear form on any noncommutative probability space.

If $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ is a t -free cumulant of order n and m_1, \dots, m_n are elements of a noncommutative probability space (\mathcal{M}, φ) , at least one of which is equal to 1, then

$$(18) \quad \sum_{\sigma \in \mathfrak{S}_n} c(\sigma) \varphi_\sigma(m_1, \dots, m_n) = 0.$$

Indeed, the subalgebra $\mathbb{C}.1$ of \mathcal{M} is t -free with any subalgebra of \mathcal{M} .

We extend the previous definition by including the free case $t = +\infty$. We will mainly consider collections $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ with the property that $c(\rho\sigma\rho^{-1}) = c(\sigma)$ for all $\sigma, \rho \in \mathfrak{S}_n$. We call such collections *conjugation-invariant*. They are in fact indexed by conjugacy classes of \mathfrak{S}_n , that is, integer partitions of n . Thus, we write use as well the notation $(c_\lambda)_{\lambda \vdash n}$ for a conjugation-invariant collection.

Our main results are the following.

THEOREM 4.3. *For all $t \in [0, +\infty]$ and all $n \in \{2, 3, 4, 5, 6\}$, there exists, up to scaling, a unique conjugation-invariant t -free cumulant of order n .*

THEOREM 4.4. *There exists a t -free cumulant of order 7 if and only if $t = 0$ or $t = +\infty$.*

Let us start by proving that we lose nothing by focusing on conjugation-invariant t -free cumulants.

LEMMA 4.5. *If for some t and some n there exists a t -free cumulant of order n , then there exists such a cumulant $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ such that moreover $c(\sigma) = c(\rho\sigma\rho^{-1})$ for all $\sigma, \rho \in \mathfrak{S}_n$.*

PROOF. The point is that the order of the arguments is arbitrary in (17). Hence, if (17) is satisfied, then for all $\rho \in \mathfrak{S}_n$,

$$\begin{aligned} 0 &= \sum_{\sigma \in \mathfrak{S}_n} c(\sigma)\varphi_\sigma(m_{\rho^{-1}(1)}, \dots, m_{\rho^{-1}(n)}) = \sum_{\sigma \in \mathfrak{S}_n} c(\sigma)\varphi_{\rho^{-1}\sigma\rho}(m_1, \dots, m_n) \\ &= \sum_{\sigma \in \mathfrak{S}_n} c(\rho\sigma\rho^{-1})\varphi_\sigma(m_1, \dots, m_n). \end{aligned}$$

Hence, if $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ is a t -free cumulant, then so is $(c(\rho\sigma\rho^{-1}))_{\sigma \in \mathfrak{S}_n}$. By averaging over ρ , we get a conjugation-invariant cumulant. \square

Observe that the assumption made in the definition of a cumulant that the sum of $c(\sigma)$ when σ spans the n -cycles is nonzero implies that $c_n \neq 0$ for any conjugation-invariant cumulant.

Let us introduce some notation. Given a permutation σ of $\{1, \dots, n\}$, we denote by $\{\{\sigma\}\}$ the partition of $\{1, \dots, n\}$ whose blocks are the sets underlying the cycles of σ . Let $\mathcal{P}(n)$ denote the set of partitions of the set $\{1, \dots, n\}$. Let (\mathcal{A}, φ) be a commutative noncommutative probability space. For each partition $\pi \in \mathcal{P}(n)$, we define an n -linear form φ_π on \mathcal{A} by setting $\varphi_\pi = \varphi_\sigma$, where σ is any permutation of $\{1, \dots, n\}$ such that $\{\{\sigma\}\} = \pi$. Since \mathcal{A} is commutative, this definition does not depend on the choice of σ . Finally, when σ is a permutation of $\{1, \dots, n\}$, we say that $i, j \in \{1, \dots, n\}$ are *consecutive* in a cycle of σ if $\sigma(i) = j$ or $\sigma(j) = i$. We will use repeatedly the following fact, which is a consequence of Proposition 2.7 and Proposition 1 of [6].

LEMMA 4.6. *Choose two integers $k, l > 0$, $t \geq 0$, and set $n = k + l$.*

(1) *There exists universal coefficients $(C(\sigma, \pi, \pi'))_{\sigma \in \mathfrak{S}_n, \pi \in \mathcal{P}(k), \pi' \in \mathcal{P}(l)}$, depending on t , such that the following property holds:*

Let \mathcal{A} and \mathcal{B} be two commutative sub-algebras of some noncommutative probability space (\mathcal{M}, φ) which are t -free with respect to φ . Consider $\sigma \in \mathfrak{S}_n$. For all $a_1, \dots, a_k \in \mathcal{A}$ and all $b_1, \dots, b_l \in \mathcal{B}$,

$$(19) \quad \begin{aligned} &\varphi_\sigma(a_1, \dots, a_k, b_1, \dots, b_l) \\ &= \sum_{\pi \in \mathcal{P}(k), \pi' \in \mathcal{P}(l)} C(\sigma, \pi, \pi') \varphi_\pi(a_1, \dots, a_k) \varphi_{\pi'}(b_1, \dots, b_l). \end{aligned}$$

(2) *The coefficient $C(\sigma, \pi, \pi')$ can be nonzero only if every block of π is contained in a block of $\{\{\sigma\}\}$.*

(3) *If two elements i and j of $\{1, \dots, k\}$ are consecutive in a cycle of σ , then $C(\sigma, \pi, \pi')$ can be nonzero only if i and j are in the same block of π .*

(4) *The parts (2) and (3) of this lemma are also valid for π' (modulo a translation of k of the indices, since π' is a partition of $\{1, \dots, l\}$ and not of $\{k + 1, \dots, k + l\}$).*

With the notation of the lemma above, we associate to every collection $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ the following family of coefficients:

$$(20) \quad \forall \pi \in \mathcal{P}(k), \pi' \in \mathcal{P}(l) \quad D_c(\pi, \pi') = \sum_{\sigma \in \mathfrak{S}_n} c(\sigma) C(\sigma, \pi, \pi'),$$

which will play an important role in the proofs of Theorems 4.4 and 4.3.

PROOF OF THEOREM 4.4. Let us choose $t > 0$ a positive real. We prove by contradiction that there exists no t -free cumulant of order 7. So, let us assume that there exists one and let $(c(\sigma))_{\sigma \in \mathfrak{S}_7}$ be one of them, which we choose to be conjugation-invariant thanks to Lemma 4.5. Thus, we denote it also by $(c_\lambda)_{\lambda \vdash 7}$. Since $c_7 \neq 0$, we may and will assume that $c_7 = 1$. Then, we proceed as follows.

Let us consider a noncommutative probability space (\mathcal{M}, φ) and two commutative sub-algebras \mathcal{A} and \mathcal{B} of \mathcal{M} which are t -free with respect to φ . Let us choose $a_1, a_2, a_3 \in \mathcal{A}$ and $b_1, b_2, b_3, b_4, b_5 \in \mathcal{B}$, which we assume to be all centered. Set $k_7 = \sum_{\sigma \in \mathfrak{S}_7} c(\sigma) \varphi_\sigma$. By using the t -freeness of \mathcal{A} and \mathcal{B} , we will express $k_7(a_1, a_2, b_1, b_2, b_3, b_4, b_5)$ and $k_7(a_1, a_2, a_3, b_1, b_2, b_3, b_4)$ in terms of the coefficients $(c_\lambda)_{\lambda \vdash 7}$, the joint moments of a_1, a_2, a_3 , and the joint moments of b_1, b_2, b_3, b_4, b_5 . By the assumption that k_7 is a t -free cumulant, the two expressions that we get must vanish. Since the joint distributions of the a 's and of the b 's are both arbitrary among those of families of centered elements, every coefficient of a given product of moments of the a 's and b 's must vanish. This gives us linear relations on the coefficients $(c_\lambda)_{\lambda \vdash 7}$, which will turn out to be incompatible.

Let us start with $k_7(a_1, a_2, b_1, b_2, b_3, b_4, b_5)$. By Lemma 4.6, this quantity can be written as

$$\begin{aligned}
 (21) \quad & \sum_{\sigma \in \mathfrak{S}_7, \pi \in \mathcal{P}(2), \pi' \in \mathcal{P}(5)} c(\sigma) C(\sigma, \pi, \pi') \varphi_\pi(a_1, a_2) \varphi_{\pi'}(b_1, b_2, b_3, b_4, b_5) \\
 & = \sum_{\pi \in \mathcal{P}(2), \pi' \in \mathcal{P}(5)} D_c(\pi, \pi') \varphi_\pi(a_1, a_2) \varphi_{\pi'}(b_1, b_2, b_3, b_4, b_5).
 \end{aligned}$$

We are thus interested in computing, for each pair (π, π') , the coefficient $D_c(\pi, \pi')$. It turns out to be convenient to think of b_1, \dots, b_5 as occupying the slots 3 to 7 rather than 1 to 5 and to see π' as a partition of the set $\{3, \dots, 7\}$ accordingly. We hope that no ambiguity will result from this change in our conventions.

Since we have chosen to consider elements which are centered, the sum (21) can be restricted to pairs of partitions without singletons. This leaves us with the following pairs $(\pi, \pi') : (\{\{1, 2\}\}, \{\{3, 4, 5, 6, 7\}\})$, $(\{\{1, 2\}\}, \{\{3, 4\}, \{5, 6, 7\}\})$ and those which are deduced from the latter by permuting 3, 4, 5, 6, 7.

Let us compute $D_c(\{\{1, 2\}\}, \{\{3, 4, 5, 6, 7\}\})$. By the second assertion of Lemma 4.6, the permutations σ which contribute to this term must have 1, 2 on one hand, and 3, 4, 5, 6, 7 on the other hand, in the same cycle. This can occur if σ is either a 7-cycle or the product of a 2-cycle and a 5-cycle.

Let us first compute the contribution of 7-cycles. The coefficient $C(\sigma, \{1, 2\}, \{3, 4, 5, 6, 7\})$ is not the same for all 7-cycles σ . We must distinguish between those in which 1 and 2 are consecutive and those in which they are not. There are $2!5!$ 7-cycles in which 1 and 2 are consecutive. For each such cycle σ , $C(\sigma, \{1, 2\}, \{3, 4, 5, 6, 7\}) = 1$, thanks to (12). In a cycle where 1 and 2 are not consecutive, there may be one, two, three or four elements of $\{3, 4, 5, 6, 7\}$ between 1 and 2. In each case, there are $5!$ cycles, each contributing a factor e^{-2t} , thanks to (13).

Let us now compute the contribution of products of a transposition and a 5-cycle. There are $1!4!$ permutations with two cycles, one which contains 1, 2 and the other 3, 4, 5, 6, 7. Each such permutation contributes a factor c_{52} .

Altogether, we have found that

$$(22) \quad D_c(\{\{1, 2\}\}, \{\{3, 4, 5, 6, 7\}\}) = 24(c_{52} + 10(1 + 2e^{-2t})).$$

Let us now compute $D_c(\{\{1, 2\}\}, \{\{3, 4\}, \{5, 6, 7\}\})$. By the second assertion of Lemma 4.6, there are five possibilities for the partition $\{\{\sigma\}\}$ underlying a permutation σ which contributes to this coefficient. We study them one after the other.

• $\{\{\sigma\}\} = \{\{1, 2, 3, 4, 5, 6, 7\}\}$. Since, by the third assertion of Lemma 4.6, any two elements of $\{3, 4, 5, 6, 7\}$ which are consecutive in σ must be in the same block of $\pi' = \{\{3, 4\}, \{5, 6, 7\}\}$, no element of $\{3, 4\}$ can be consecutive to an element of $\{5, 6, 7\}$ in σ . Since there are only two a 's, the only possibility is that 3 and 4 on one hand, and 5, 6, and 7 on the other hand, are consecutive in σ and separated by

1 and 2. There are $2!2!3!$ 7-cycles with this property. Each of them contributes to the sum with a factor $1 - e^{-2t}$, according to (13).

- $\{\{\sigma\}\} = \{\{1, 2\}, \{3, 4, 5, 6, 7\}\}$. By the third assertion of Lemma 4.6, these permutations do not contribute.

- $\{\{\sigma\}\} = \{\{1, 2, 3, 4\}, \{5, 6, 7\}\}$. There are two possible structures for the 4-cycle of σ in this case. Either the a 's and the b 's are consecutive, or they are intertwined. In the first situation, there are $2!2!2!$ permutations, each of which contributes c_{43} , thanks to (12). In the second situation, there are $2!2!$ permutations, because of a higher symmetry, each of which contributes $e^{-2t}c_{43}$, thanks to (13).

- $\{\{\sigma\}\} = \{\{1, 2, 5, 6, 7\}, \{3, 4\}\}$. Again, there are two possible structures for the 5-cycle of σ , depending on whether the a 's are consecutive or not. There are $2!3!$ permutations where they are, and each contributes c_{52} . There are also $2!3!$ permutations where they are not, each contributing $e^{-2t}c_{52}$.

- $\{\{\sigma\}\} = \{\{1, 2\}, \{3, 4\}, \{5, 6, 7\}\}$. This is the simplest situation. There are 2 permutations with this cycle structure and each contributes c_{322} .

Finally,

$$(23) \quad \begin{aligned} & D_c(\{\{1, 2\}\}, \{\{3, 4\}, \{5, 6, 7\}\}) \\ &= 2(c_{322} + 2(2 + e^{-2t})c_{43} + 6(1 + e^{-2t})c_{52} + 12(1 - e^{-2t})). \end{aligned}$$

Let us perform the same kind of computations on

$$\begin{aligned} & k_7(a_1, a_2, a_3, b_1, b_2, b_3, b_4) \\ &= \sum_{\pi \in \mathcal{P}(\{1,2,3\}), \pi' \in \mathcal{P}(\{4,5,6,7\})} D_c(\pi, \pi') \varphi_\pi(a_1, a_2, a_3) \varphi_{\pi'}(b_1, b_2, b_3, b_4). \end{aligned}$$

Since our variables are centered, the only pairs of partitions which occur in the sum are $(\{\{1, 2, 3\}\}, \{\{4, 5, 6, 7\}\})$, $(\{\{1, 2, 3\}\}, \{\{4, 5\}, \{6, 7\}\})$ and those which are deduced from the latter by permuting 4, 5, 6, 7.

Let us compute $D_c(\{\{1, 2, 3\}\}, \{\{4, 5, 6, 7\}\})$. We shall again use formulae (12)–(14) several times, without citing them every time. The permutations which contribute to this coefficient are 7-cycles and products of a 3-cycle and a 4-cycle. As before, all 7-cycles do not contribute in the same way. If the a 's are consecutive, which is the case for $3!4!$ 7-cycles, the contribution is simply 1. If two a 's are consecutive and the third is on its own, the 7-cycle contributes e^{-2t} . In this case, there can be one, two or three b 's between the isolated a and the pair of consecutive a 's, in the cyclic order. In each case, there are $3!4!$ possible 7-cycles. Finally, the three a 's can be isolated. This happens in $3!4!$ 7-cycles, and each contributes e^{-3t} , thanks to (14). So far, we have a contribution of $144(1 + 3e^{-2t} + e^{-3t})$. The contribution of products of a 3-cycle and a 4-cycle is much simpler to compute: it is $2!3!c_{43}$. We find

$$(24) \quad D_c(\{\{1, 2, 3\}\}, \{\{4, 5, 6, 7\}\}) = 12(c_{43} + 12(1 + 3e^{-2t} + e^{-3t})).$$

Let us finally compute $D_c(\{\{1, 2, 3\}\}, \{\{4, 5\}, \{6, 7\}\})$. Again, by Lemma 4.6, there are five possibilities for the partition $\{\{\sigma\}\}$, which we examine one after the other.

- $\{\{\sigma\}\} = \{\{1, 2, 3, 4, 5, 6, 7\}\}$. No element of $\{4, 5\}$ can be consecutive with an element of $\{6, 7\}$ in σ . Still, there are several possibilities. Let us first consider the 7-cycles where 4, 5 on one hand and 6, 7 on the other hand are consecutive. These two groups must be separated by a 's. There are $2!2!2!3!$ such 7-cycles, and each contributes for $1 - e^{-2t}$, according to (13). Since there are only three a 's, one at least of the two pairs $\{4, 5\}$ and $\{6, 7\}$ must be consecutive. However, it can happen that one is not. This happens in $2!2!2!3!$ 7-cycles, and according to (14), each contributes for $e^{-2t} - e^{-3t}$.

- $\{\{\sigma\}\} = \{\{1, 2, 3\}, \{4, 5, 6, 7\}\}$. These permutations do not contribute.

- $\{\{\sigma\}\} = \{\{1, 2, 3, 4, 5\}, \{6, 7\}\}$. As usual by now, there are two possibilities in the 5-cycle of σ . Either the two b 's are consecutive, which happens in $2!3!$ cases with the contribution c_{52} , or they are not. This happens in $2!3!$ cases, and each case contributes for $e^{-2t}c_{52}$.

- $\{\{\sigma\}\} = \{\{1, 2, 3, 6, 7\}, \{4, 5\}\}$. By symmetry, this contribution is equal to the one above.

- $\{\{\sigma\}\} = \{\{1, 2, 3\}, \{4, 5\}, \{6, 7\}\}$. There are 2 permutations, each contributing for c_{322} .

Finally,

$$(25) \quad D_c(\{\{1, 2, 3\}\}, \{\{4, 5\}, \{6, 7\}\}) = 2(c_{322} + 12(1 + e^{-2t})c_{52} + 24(1 - e^{-3t})).$$

Let us summarize our results. We have proved that, if there exists a t -free cumulant of order 7, whose associated 7-linear form is denoted by k_7 , then for all centered $a_1, a_2, a_3 \in \mathcal{A}$ and $b_1, \dots, b_5 \in \mathcal{B}$, the following equalities hold:

$$\begin{aligned} &k_7(a_1, a_2, b_1, b_2, b_3, b_4, b_5) \\ &= 24(c_{52} + 10(1 + 2e^{-2t}))\varphi(a_1a_2)\varphi(b_1b_2b_3b_4b_5) \\ &\quad + 2(c_{322} + 2(2 + e^{-2t})c_{43} \\ &\quad\quad + 6(1 + e^{-2t})c_{52} + 12(1 - e^{-2t}))\varphi(a_1a_2)\varphi(b_1b_2)\varphi(b_3b_4b_5) \\ &\quad + 2(c_{322} + 2(2 + e^{-2t})c_{43} \\ &\quad\quad + 6(1 + e^{-2t})c_{52} + 12(1 - e^{-2t}))\varphi(a_1a_2)\varphi(b_1b_3)\varphi(b_2b_4b_5) + \dots, \end{aligned}$$

where all partitions of $\{b_1, b_2, b_3, b_4, b_5\}$ into a pair and a triple appear, and

$$\begin{aligned} &k_7(a_1, a_2, a_3, b_1, b_2, b_3, b_4) \\ &= 12(c_{43} + 12(1 + 3e^{-2t} + e^{-3t}))\varphi(a_1a_2a_3)\varphi(b_1b_2b_3b_4) \\ &\quad + 2(c_{322} + 12(1 + e^{-2t})c_{52} + 24(1 - e^{-3t}))\varphi(a_1a_2a_3)\varphi(b_1b_2)\varphi(b_3b_4) \\ &\quad + 2(c_{322} + 12(1 + e^{-2t})c_{52} + 24(1 - e^{-3t}))\varphi(a_1a_2a_3)\varphi(b_1b_3)\varphi(b_2b_4) \\ &\quad + 2(c_{322} + 12(1 + e^{-2t})c_{52} + 24(1 - e^{-3t}))\varphi(a_1a_2a_3)\varphi(b_1b_4)\varphi(b_2b_3). \end{aligned}$$

Since k_7 is a t -free cumulant, these two expressions must be zero for all choices of a 's and b 's. Since the joint distributions of the a 's and of the b 's are both arbitrary among those of families of centered elements, this implies that the coefficients which appear in these equalities in front of the various products of moments of a 's and b 's must vanish. This implies the following relations:

$$\begin{aligned} c_{52} &= -10(1 + 2e^{-2t}), \\ c_{43} &= -12(1 + 3e^{-2t} + e^{-3t}), \\ c_{322} &= -2(2 + e^{-2t})c_{43} - 6(1 + e^{-2t})c_{52} - 12(1 - e^{-2t}), \\ c_{322} &= -12(1 + e^{-2t})c_{52} - 24(1 - e^{-3t}). \end{aligned}$$

It does not take a long computation to see that the two expressions of c_{322} are different, since the first involves e^{-5t} , whereas the second does not. We leave it to the reader to check that the difference between the two values of c_{322} that we have obtained is equal to $24e^{-3t}(1 - e^{-t})^2$. This quantity vanishes only for $t = 0$ or $t = +\infty$. \square

In order to prove that t -free cumulants of order at most 6 exist, we are going to construct them. We prove first a lemma which settles the problem of the coefficients c_λ for the partitions λ whose smallest part is 1.

Let us introduce some notation. Let $\mu = (\mu_1 \geq \dots \geq \mu_r)$ be a partition of some integer n . We denote by $\ell(\mu)$ the number of nonzero parts of μ and we write $\mu \vdash n$ if $\mu_{\ell(\mu)} \geq 2$, that is, if μ has no part equal to 1. Let $i \geq 1$ an integer. We denote by $\mu + \delta_i$ the partition of $n + 1$ whose parts are $\mu_1, \dots, \mu_{i-1}, \mu_i + 1, \mu_{i+1}, \dots, \mu_r$, rearranged in nonincreasing order. If $i > \ell(\mu)$, then $\mu + \delta_i$ is simply the partition μ to which a part equal to 1 has been appended.

PROPOSITION 4.7. *Let $n \geq 2$ be an integer. Choose $t \in [0, +\infty]$. A collection $(c_\lambda)_{\lambda \vdash n}$ is a t -free cumulant if and only if the following two conditions hold:*

- (1) *The relation (17) is satisfied for all m_1, \dots, m_n which are centered.*
- (2) *For all $\mu \vdash n - 1$,*

$$(26) \quad c_{\mu + \delta_{\ell(\mu)+1}} = - \sum_{i=1}^{\ell(\mu)} \mu_i c_{\mu + \delta_i}.$$

Moreover, a collection of complex numbers $(c_\lambda)_{\lambda \vdash n}$ which satisfies (17) for all m_1, \dots, m_n which are centered can be extended in a unique way into a t -free cumulant of order n .

When σ is a permutation of $\{1, \dots, n\}$, let us denote by $[\sigma]$ the partition of the integer n given by the lengths of the cycles of σ .

PROOF OF PROPOSITION 4.7. Let c be a t -free cumulant of order n . In order to check that (26) is satisfied, let us choose m_1, \dots, m_{n-1} in some probability space (\mathcal{M}, φ) and write the fact that $k_n(m_1, \dots, m_{n-1}, 1) = 0$. We find

$$(27) \quad \sum_{\lambda \vdash n} c_\lambda \sum_{\substack{\sigma \in \mathfrak{S}_n \\ [\sigma] = \lambda}} \varphi_\sigma(m_1, \dots, m_{n-1}, 1) = 0.$$

Let $r_n : \mathfrak{S}_n \rightarrow \mathfrak{S}_{n-1}$ denote the following function: for all $\sigma \in \mathfrak{S}_n$, $r_n(\sigma)$ is the permutation of $\{1, \dots, n - 1\}$ obtained by removing n from the cycle of σ which contains it. For each σ , we have the equality $\varphi_\sigma(m_1, \dots, m_{n-1}, 1) = \varphi_{r_n(\sigma)}(m_1, \dots, m_{n-1})$. Now a permutation $\tau \in \mathfrak{S}_{n-1}$ has exactly n preimages by r_n . Moreover, if $[\tau] = \mu = (\mu_1 \geq \dots \geq \mu_{\ell(\mu)} > 0) \vdash n - 1$, then all preimages of τ belong to one of the conjugacy classes $\mu + \delta_i$ for $i = 1, \dots, \ell(\mu) + 1$. Finally, $r_n^{-1}(\tau)$ contains exactly one element of $\mu + \delta_{\ell(\mu)+1}$ and μ_i elements of $\mu + \delta_i$ for $i = 1, \dots, \ell(\mu)$. We can thus rewrite (27) as follows:

$$(28) \quad \sum_{\mu \vdash n-1} \left(\sum_{i=1}^{\ell(\mu)} \mu_i c_{\mu+\delta_i} + c_{\mu+\delta_{\ell(\mu)+1}} \right) \sum_{\substack{\tau \in \mathfrak{S}_{n-1} \\ [\tau] = \mu}} \varphi_\tau(m_1, \dots, m_{n-1}) = 0.$$

Since the distribution of (m_1, \dots, m_{n-1}) is arbitrary, all the coefficients between the brackets must vanish. It follows that (26) is satisfied.

Conversely, let $(c(\sigma))_{\sigma \in \mathfrak{S}_n}$ be a collection which satisfies (17) for centered elements and (26). Then, by the computation that we have just done, this collection satisfies (18) and hence, by multilinearity, (17) for arbitrary elements.

Let us prove the last assertion. For any $\lambda \vdash n$ with at least one part equal to 1, the relation (26) expresses the value of c_λ in terms of $c_{\lambda'}$ for partitions λ' of n which have strictly less parts equal to 1 than λ . The collection $(c_\lambda)_{\lambda \vdash n}$ is thus completely and uniquely determined by $(c_\lambda)_{\lambda \vdash n}$. The fact that the resulting collection is a t -free cumulant is granted by the first part of the proposition. \square

The last result simplifies greatly the search for t -free cumulants, since it allows one to restrict to centered elements and partitions in parts at least equal to 2. We apply it to find cumulants of order less than 6. These provide a partial obstruction to the t -freeness of two random variables.

PROOF OF THEOREM 4.3. Let us prove that there exist t -free cumulants up to order 6. We proceed by first establishing enough conditions that their coefficients must satisfy, in order to determine these coefficients. Then, we check that we actually have a t -free cumulant.

We will always normalize our cumulants by the condition $c_n = 1$.

• $n = 2$. By Proposition 4.7, the condition $c_2 = 1$ suffices to determine the whole cumulant, and $c_{11} = -1$. The relation (12) implies that we have indeed got a t -free cumulant.

• $n = 3$. Again, the condition $c_3 = 1$ determines completely the cumulant. Using (26), we find $c_{21} = -2$ and $c_{111} = 4$. The relation (12) implies again that the collection thus obtained is a t -free cumulant. Indeed, by (12), the product of any three centered elements, one being t -free with the two others, is centered. Hence, our collection satisfies (17) on centered elements.

• $n = 4$. This is the first case where the relation $c_4 = 1$ does not suffice to determine the cumulant. Indeed, we must compute c_{22} . For this, let us choose in some probability space elements a_1, a_2, \dots and b_1, b_2, \dots , such that $\{a_1, a_2, \dots\}$ and $\{b_1, b_2, \dots\}$ are t -free. We will use this notation again in this proof without redefining it. Let us assume that a t -free cumulant c of order 4 is given and let us compute $D_c(\{\{1, 2\}\}, \{\{3, 4\}\})$ [see (20)]. Again, we shall use formulae (12) and (13) several times, without citing them every time. There are 4-cycles which contribute to this coefficient. In $2!2!$ of them, 1 and 2 are consecutive and they contribute for 1 each. In $2!$ others, 1 and 2 are not consecutive and each such cycle contributes for e^{-2t} . There is also one product of two 2-cycles, which contributes for c_{22} . Finally, $D_c(\{\{1, 2\}\}, \{\{3, 4\}\}) = c_{22} + 2(2 + e^{-2t})$. The nullity of this coefficient implies $c_{22} = -2(2 + e^{-2t})$. Using (26), we determine the remaining coefficients, and find

$$c_4 = 1, \quad c_{31} = -3, \quad c_{22} = -2(2 + e^{-2t}),$$

$$c_{211} = 2(5 + e^{-2t}), \quad c_{1111} = -6(5 + e^{-2t}).$$

Now let us check that the collection thus defined satisfies (17) for centered elements. Set $k_4 = \sum_{\sigma} c(\sigma)\varphi_{\sigma}$. If we expand $k_4(a_1, b_1, b_2, b_3)$ according to (19), then all terms involve $\varphi(a_1)$ and vanish. Now $k_4(a_1, a_2, b_1, b_2)$ also vanishes, because this is how we have chosen the value of c_{22} . Finally, we do have a t -free cumulant of order 4.

• $n = 5$. Let c be a t -free cumulant of order 5. Let us compute c_{32} by writing the nullity of $D_c(\{\{1, 2, 3\}\}, \{\{4, 5\}\})$ and using formulae (12) and (13). There are $3!2!$ 5-cycles in which 4 and 5 are consecutive, and they contribute for 1 each. There are also $3!2!$ 5-cycles in which they are not consecutive, and each contributes for e^{-2t} . There are finally $2!$ products of a 3-cycle and a 2-cycle, which contribute for c_{32} each. Hence, we must have $c_{32} = -6(1 + e^{-2t})$. Using as usual (26), we find that the other values of c must be

$$c_5 = 1, \quad c_{41} = -4, \quad c_{32} = -6(1 + e^{-2t}), \quad c_{311} = 6(3 + e^{-2t}),$$

$$c_{221} = 12(1 + e^{-2t}), \quad c_{2111} = -12(5 + 2e^{-2t}), \quad c_{11111} = 48(5 + 2e^{-2t}).$$

The fact that $k_5 \sum_{\sigma} c(\sigma)\varphi_{\sigma}$ is a cumulant is checked just as in the case $n = 4$: the identity $k_5(a_1, b_1, b_2, b_3, b_4) = 0$ is granted by (19) and $k_5(a_1, a_2, b_1, b_2, b_3) = 0$ by the choice of c_{32} .

• $n = 6$. Let c be a t -free cumulant of order 5. The value of c_{42} , deduced as usual from the nullity of $D_c(\{\{1, 2, 3, 4\}\}, \{\{5, 6\}\})$, is $c_{42} = -4(2 + 3e^{-2t})$. Similarly, $D_c(\{\{1, 2, 3\}\}, \{\{4, 5, 6\}\}) = 0$ gives us $c_{33} = -3(3 + 6e^{-2t} + e^{-3t})$. Finally,

$D_c(\{\{1, 2\}\}, \{\{3, 4\}, \{5, 6\}\}) = 0$ implies $c_{222} = 8(7 + 17e^{-2t} + 6e^{-4t})$. The other coefficients follow as usual from (26) and we find

$$\begin{aligned} c_6 &= 1, & c_{51} &= -5, & c_{42} &= -4(2 + 3e^{-2t}), & c_{411} &= 4(7 + 3e^{-2t}), \\ c_{33} &= -3(3 + 6e^{-2t} + e^{-3t}), & c_{321} &= 6(7 + e^{-3t} + 12e^{-2t}), \\ c_{3111} &= -12(14 + 15e^{-2t} + e^{-3t}), & c_{222} &= 8(7 + 17e^{-2t} + 6e^{-4t}), \\ c_{2211} &= -8(28 + 53e^{-2t} + 3e^{-3t} + 6e^{-4t}), \\ c_{21111} &= 48(21 + 34e^{-2t} + 2e^{-3t} + 3e^{-4t}), \\ c_{111111} &= -240(21 + 34e^{-2t} + 2e^{-3t} + 3e^{-4t}). \end{aligned}$$

Let us set $k_6 = \sum_{\sigma} c(\sigma)\varphi_{\sigma}$. The nullity of $k_6(a_1, b_1, \dots, b_5)$ follows as usual from (19). That of $k_6(a_1, a_2, b_1, b_2, b_3, b_4)$ results from the choices of c_{42} and c_{222} . Finally, $k_6(a_1, a_2, a_3, b_1, b_2, b_3) = 0$ is granted by the choice of c_{33} .

Nowhere there has been any freedom in the definition of the cumulants. This shows the uniqueness of conjugation-invariant t -free cumulants of order at most than 6. \square

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