# DETERMINANTAL POINT PROCESSES WITH J-HERMITIAN CORRELATION KERNELS 

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#### Abstract

Let $X$ be a locally compact Polish space and let $m$ be a reference Radon measure on $X$. Let $\Gamma_{X}$ denote the configuration space over $X$, that is, the space of all locally finite subsets of $X$. A point process on $X$ is a probability measure on $\Gamma_{X}$. A point process $\mu$ is called determinantal if its correlation functions have the form $k^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}$. The function $K(x, y)$ is called the correlation kernel of the determinantal point process $\mu$. Assume that the space $X$ is split into two parts: $X=X_{1} \sqcup X_{2}$. A kernel $K(x, y)$ is called $J$-Hermitian if it is Hermitian on $X_{1} \times X_{1}$ and $X_{2} \times X_{2}$, and $K(x, y)=-\overline{K(y, x)}$ for $x \in X_{1}$ and $y \in X_{2}$. We derive a necessary and sufficient condition of existence of a determinantal point process with a $J$-Hermitian correlation kernel $K(x, y)$.


## 1. Introduction and preliminaries.

1.1. Macchi-Soshnikov theorem. Let $X$ be a locally compact Polish space, let $\mathcal{B}(X)$ be the Borel $\sigma$-algebra on $X$, and let $\mathcal{B}_{0}(X)$ denote the collection of all sets from $\mathcal{B}(X)$ which are pre-compact. The configuration space over $X$ is defined as the set of all locally finite subsets of $X$ :

$$
\Gamma:=\Gamma_{X}:=\left\{\gamma \subset X \mid \text { for all } \Delta \in \mathcal{B}_{0}(X)|\gamma \cap \Delta|<\infty\right\} .
$$

Here, for a set $\Lambda,|\Lambda|$ denotes its capacity. Elements $\gamma \in \Gamma$ are called configurations. The space $\Gamma$ can be endowed with the vague topology, that is, the weakest topology on $\Gamma$ with respect to which all maps $\Gamma \ni \gamma \mapsto \sum_{x \in \gamma} f(x), f \in C_{0}(X)$, are continuous. Here $C_{0}(X)$ is the space of all continuous real-valued functions on $X$ with compact support. The configuration space $\Gamma$ equipped with the vague topology is a Polish space. We will denote by $\mathcal{B}(\Gamma)$ the Borel $\sigma$-algebra on $\Gamma$. A probability measure $\mu$ on $(\Gamma, \mathcal{B}(\Gamma))$ is called a point process on $X$. For more detail, see, for example, [9, 11, 13, 16].

A point process $\mu$ can be described with the help of correlation functions, if they exist. Let $m$ be a reference Radon measure on $(X, \mathcal{B}(X))$. The $n$th correlation function of $\mu(n \in \mathbb{N})$ is an $m^{\otimes n}$-a.e. nonnegative measurable symmetric

[^0]function $k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right)$ on $X^{n}$ such that, for any measurable symmetric function $f^{(n)}: X^{n} \rightarrow[0, \infty]$,
\[

$$
\begin{align*}
& \int_{\Gamma_{\left\{x_{1}, \ldots, x_{n}\right\} \subset \gamma}} \sum^{(n)}\left(x_{1}, \ldots, x_{n}\right) \mu(d \gamma) \\
& \quad=\frac{1}{n!} \int_{X^{n}} f^{(n)}\left(x_{1}, \ldots, x_{n}\right) k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right) m\left(d x_{1}\right) \cdots m\left(d x_{n}\right) . \tag{1}
\end{align*}
$$
\]

Under a mild condition on the growth of correlation functions as $n \rightarrow \infty$, they determine a point process uniquely [13].

A point process $\mu$ is called determinantal if there exists a function $K(x, y)$ on $X^{2}$, called the correlation kernel, such that

$$
\begin{equation*}
k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1}^{n}, \quad n \in \mathbb{N} \tag{2}
\end{equation*}
$$

see, for example, [21]. The integral operator $K$ in $L^{2}(X, m)$ which has integral kernel $K(x, y)$ is called the correlation operator of $\mu$.

Assume that the correlation operator $K$ is self-adjoint and bounded on the (real or complex) Hilbert space $L^{2}(X, m)$. In particular, the integral kernel $K(x, y)$ is Hermitian (symmetric in the real case). If the correlation functions $\left(k_{\mu}^{(n)}\right)_{n \in \mathbb{N}}$ in (2) are pointwisely nonnegative, then $K(x, y)$ is a positive definite kernel. Hence, if additionally the function $K(x, y)$ is continuous (it being possible to weaken the latter condition), then the operator $K$ must be nonnegative ( $K \geq 0$ ).

A bounded linear operator $K$ on $L^{2}(X, m)$ is called a locally trace-class operator if, for each $\Delta \in \mathcal{B}_{0}(X)$, the operator $K^{\Delta}:=P^{\Delta} K P^{\Delta}$ is trace-class. Here $P^{\Delta}$ denotes the operator of multiplication by $\chi_{\Delta}$, the indicator function of the set $\Delta$. [Thus, $P^{\Delta}$ is the orthogonal projection of $L^{2}(X, m)$ onto $L^{2}(\Delta, m)$.] If the operator $K$ is self-adjoint and nonnegative, then we can and will choose its integral kernel, $K(x, y)$, so that

$$
\operatorname{Tr} K^{\Delta}=\int_{\Delta} K(x, x) m(d x) \quad \text { for each } \Delta \in \mathcal{B}_{0}(X)
$$

see [21] and [10]. By (1) and (2), for each $\Delta \in \mathcal{B}_{0}(X)$,

$$
\int_{\Gamma}|\gamma \cap \Delta| \mu(d \gamma)=\int_{\Delta} K(x, x) m(d x)
$$

Hence, in order that the correlation functions of $\mu$ be finite, we must indeed assume that the operator $K$ is locally trace-class.

The following theorem, which is due to Macchi [15] and Soshnikov [21], plays a fundamental role in the theory of point processes.

THEOREM 1 (Macchi-Soshnikov). Let $K$ be a self-adjoint, nonnegative, locally trace-class, bounded linear operator on $L^{2}(X, m)$. Then the integral kernel $K(x, y)$ of the operator $K$ is the correlation kernel of a determinantal point process if and only if $0 \leq K \leq 1$.

Note that, in the above theorem, the condition of boundedness of the operator $K$ is not essential. One may instead initially assume that $K$ is a Hermitian, nonnegative, locally trace-class operator which is defined on a proper domain in $L^{2}(X, m)$.

Determinantal point processes with Hermitian correlation kernels occur in various fields of mathematics and physics; see, for example, the review paper [21] and Chapter 4 in [1].
1.2. Complementation principle (particle-hole duality). Assume that the underlying space $X$ is split into two disjoint parts: $X=X_{1} \sqcup X_{2}$. Hence, we get $L^{2}(X, m)=L^{2}\left(X_{1}, m\right) \oplus L^{2}\left(X_{2}, m\right)$. For $i=1,2$, let $P_{i}$ denote the orthogonal projection of $L^{2}(X, m)$ onto $L^{2}\left(X_{i}, m\right)$. Let us define a bounded linear operator $J$ on $L^{2}(X, m)$ by $J:=P_{1}-P_{2}$. Following, for example, [2], we define an (indefinite) $J$-scalar product on $L^{2}(X, m)$ by

$$
[f, g]:=(J f, g)=\left(P_{1} f, P_{1} g\right)-\left(P_{2} f, P_{2} g\right), \quad f, g \in L^{2}(X, m)
$$

Here $(\cdot, \cdot)$ denotes the usual scalar product in $L^{2}(X, m)$. A bounded linear operator $K$ on $L^{2}(X, m)$ is called $J$-self-adjoint if $[K f, g]=[f, K g]$ for all $f, g \in$ $L^{2}(X, m)$. An integral kernel $K(x, y)$ of a $J$-self-adjoint, integral operator $K$ is called $J$-Hermitian. More precisely, $K(x, y)$ is $J$-Hermitian if $K(x, y)=\overline{K(y, x)}$ if $x, y \in X_{1}$ or $x, y \in X_{2}$, and $K(x, y)=-\overline{K(y, x)}$ if $x \in X_{1}, y \in X_{2}$.

For a bounded linear operator $K$ on $L^{2}(X, m)$, we denote

$$
\begin{equation*}
\widehat{K}:=K P_{1}+(1-K) P_{2} \tag{3}
\end{equation*}
$$

As is easily seen, $K$ is $J$-self-adjoint if and only if $\widehat{K}$ is self-adjoint.
Assume now that the underlying space $X$ is discrete, that is, $X$ is a countable set, and as a topological space $X$ it totally disconnected. Thus, a configuration $\gamma$ in $X$ is an arbitrary subset of $X$. Let $m$ be the counting measure on $X: m(\{x\})=1$ for each $x \in X$. Any linear operator $K$ in $L^{2}(X, m)$ may be identified with its matrix [ $K(x, y)]_{x, y \in X}[K(x, y)$ being the integral kernel of $K$ in this case].

Let $\mu$ be a point process on $X$. By (1),

$$
k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\mu\left(\gamma \in \Gamma:\left\{x_{1}, \ldots, x_{n}\right\} \subset \gamma\right)
$$

for distinct points $x_{1}, \ldots, x_{n} \in X$, otherwise $k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=0$. In particular, the correlation functions uniquely identify the corresponding point process.

Following [4], we will now present a complementation principle (a particlehole duality) for determinantal point processes. (This observation is referred by the authors of [4] to a private communication by S. Kerov.) Assume, as above, that the underlying space $X$ is divided into two disjoint parts: $X=X_{1} \sqcup X_{2}$. Consider the mapping $I: \Gamma \rightarrow \Gamma$ defined by

$$
I \gamma:=\widehat{\gamma}:=\left(\gamma \cap X_{1}\right) \cup\left(X_{2} \backslash \gamma\right) .
$$

Thus, on the $X_{1}$ part of the space, the configuration $\widehat{\gamma}$ coincides with $\gamma$, while on the $X_{2}$ part the configuration $\widehat{\gamma}$ consists of all points from $X_{2}$ which do not belong
to $\gamma$ (holes). The mapping $I$ is clearly an involution, that is, $I^{2}$ is the identity mapping. For a point process $\mu$ on $X$, we denote by $\widehat{\mu}$ the push-forward of $\mu$ under $I$.

Proposition 1 ([4]). Let $\mu$ be an arbitrary determinantal point process on a discrete space $X=X_{1} \sqcup X_{2}$, with a correlation kernel $K(x, y)$. Then $\widehat{\mu}$ is the determinantal point process on $X$ with the correlation kernel $\widehat{K}(x, y)$, the integral kernel of the operator $\widehat{K}$ defined by (3).

Combining the Macchi-Soshnikov theorem with Proposition 1, we immediately get the following:

Proposition 2. Let $X=X_{1} \sqcup X_{2}$ be a discrete space and let $m$ be a counting measure on $X$. Let $K$ be a bounded linear operator on $L^{2}(X, m)$ and let $K$ be $J$-self-adjoint. Then $K(x, y)$ is the correlation kernel of a determinantal point process on $X$ if and only if $0 \leq \widehat{K} \leq 1$.
1.3. Formulation of the problem and the main result. In the case of a discrete underlying space $X$, determinantal point processes with $J$-Hermitian correlation kernels occurred in Borodin and Olshanski's studies on harmonic analysis of both the infinite symmetric group and the infinite-dimensional unitary group; see, for example, [5-8, 17] and the references therein. The paper [7], page 1332, also contains references to some earlier works of mathematical physicists on solvable models of systems with positive and negative charged particles. In these papers, one finds further examples of determinantal point processes with $J$-Hermitian correlation kernels.

Furthermore, in their studies, Borodin and Olshanski derived three classes of determinantal point processes with $J$-Hermitian correlation kernels in the case where the underlying space $X$ is continuous: the Whittaker kernel [6] $\left(X=\mathbb{R}_{-} \sqcup \mathbb{R}_{+}\right)$, its scaling limit-the matrix tail kernel $[17](X=\mathbb{R} \sqcup \mathbb{R})$, and the continuous hypergeometric kernel [7] [ $\left.X=\left(-\frac{1}{2}, \frac{1}{2}\right) \sqcup\left\{x \in \mathbb{R}:|x|>\frac{1}{2}\right\}\right]$. It is important to note that, in all these examples, the self-adjoint operator $\widehat{K}$ appears to be an orthogonal projection. This follows from Proposition 5.1 in [8] and the respective results of [6, 17] (see also [8], Proposition 6.6) and [7]. (It should be, however, noted that, in the case of a continuous hypergeometric kernel, the corresponding projection property was proved only under an additional assumption; see the last two paragraphs of Section 10 in [7].)

The aim of the present paper is to derive, in the case of a general underlying space $X$, a necessary and sufficient condition of existence of a determinantal point process with a J-Hermitian correlation kernel. This problem was formulated to the author by Grigori Olshanski. I am extremely grateful to him for this and for many useful discussions and suggestions.

Our main result may be stated as follows. (We will omit a technical detail related to the choice of an integral kernel of the operator K.)

Main result. Assume that $K$ is a $J$-self-adjoint bounded linear operator on $L^{2}(X, m)$. Assume that the operators $P_{1} K P_{1}$ and $P_{2} K P_{2}$ are nonnegative. Assume that, for any $\Delta_{1}, \Delta_{2} \in \mathcal{B}_{0}(X)$ such that $\Delta_{1} \subset X_{1}$ and $\Delta_{2} \subset X_{2}$, the operators $K^{\Delta_{i}}(i=1,2)$ are trace-class, while $P^{\Delta_{2}} K P^{\Delta_{1}}$ is a Hilbert-Schmidt operator. Then the integral kernel $K(x, y)$ of the operator $K$ is the correlation kernel of a determinantal point process if and only if $0 \leq \widehat{K} \leq 1$.

Let us make two remarks regarding the conditions of the main result. First, we note that, if the correlation operator $K$ of $\mu$ is $J$-self-adjoint, then the restrictions of the point process $\mu$ to $X_{1}$ and $X_{2}$ are determinantal point processes on $X_{1}$ and $X_{2}$ with self-adjoint, correlation operators $P_{1} K P_{1}$ and $P_{2} K P_{2}$, respectively. Therefore, we assume that the latter operators are nonnegative.

Second, choose any $\Delta \in \mathcal{B}_{0}(X)$ such that $m\left(\Delta_{i}\right)>0$, where $\Delta_{i}:=\Delta \cap X_{i}$, $i=1,2$. Then, since the operator $K$ is not self-adjoint, the assumption in the main result is weaker than the requirement that the operator $K^{\Delta}$ be trace-class. In fact, $K$ being locally trace-class seems to be a rather unnatural assumption for $J$-selfadjoint operators. This, of course, leads us to some additional difficulties in the proof.

Clearly, Proposition 2 is the special case of our main result in the case where the underlying space $X$ is discrete. The drastic difference between the discrete and the continuous cases is that the mapping $\gamma \mapsto \widehat{\gamma}$ has no analog in the case of a continuous space $X$. Furthermore, if the space $X$ is not discrete, the self-adjoint operator $\widehat{K}$ is not even an integral operator, so it cannot be a correlation operator of a determinantal point process.

To prove the main result, we follow the strategy of dealing with determinantal point processes through the corresponding Fredholm determinants (compare with [15, 19, 21]), or rather the extension of Fredholm determinant as proposed in [4].

Combining the main result and Proposition 5.1 in [8], we also derive a method of constructing a big class of determinantal point processes with $J$-self-adjoint correlation operators $K$ such that the corresponding operators $\widehat{K}$ are orthogonal projections. This class includes the above mentioned examples of determinantal point processes obtained by Borodin and Olshanski.

The paper is organized as follows. In Section 2 we prove a couple of results related to the mentioned extension of the Fredholm determinant. In Section 3 we prove a series of auxiliary statements regarding $J$-self-adjoint operators and their extended Fredholm determinants. Finally, in Section 4 we formulate and prove the main results of the paper.
2. An extension of the Fredholm determinant. We first recall the classical definition of a Fredholm determinant; see, for example, [20] for further detail. Let
$H$ be a complex, separable Hilbert space with scalar product $(\cdot, \cdot)$ and norm $\|\cdot\|$. We denote by $\mathscr{L}(H)$ the space of all bounded linear operators on $H$. An operator $A \in \mathscr{L}(H)$ is called a trace-class operator if $\|A\|_{1}=\operatorname{Tr}(|A|)<\infty$, where $|A|=$ $\left(A^{*} A\right)^{1 / 2}$. The set of all trace-class operators in $H$ will be denoted by $\mathscr{L}_{1}(H)$. The trace of an operator $A \in \mathscr{L}_{1}(H)$ is given by $\operatorname{Tr}(A)=\sum_{n=1}^{\infty}\left(A e_{n}, e_{n}\right)$, where $\left\{e_{n}\right\}_{n=1}^{\infty}$ is an orthonormal basis of $H . \operatorname{Tr}(A)$ is independent of the choice of a basis. Note also that $|\operatorname{Tr}(A)| \leq \operatorname{Tr}(|A|)$. For any $A \in \mathscr{L}_{1}(H)$ and $B \in \mathscr{L}(H)$, we have $A B, B A \in \mathscr{L}_{1}(H)$ with

$$
\max \left\{\|A B\|_{1},\|B A\|_{1}\right\} \leq\|A\|_{1}\|B\|
$$

where $\|B\|$ denotes the usual operator norm of $B$. In the latter case, we have

$$
\begin{equation*}
\operatorname{Tr}(A B)=\operatorname{Tr}(B A) \tag{4}
\end{equation*}
$$

Denote by $\wedge^{n}(H)$ the $n$th antisymmetric tensor power of the Hilbert space $H$, which is a closed subspace of $H^{\otimes n}$, the $n$th tensor power of $H$. For any $A \in \mathscr{L}(H)$, the operator $A^{\otimes n}$ in $H^{\otimes n}$ acts invariantly on $\wedge^{n}(H)$ and we denote by $\wedge^{n}(A)$ the restriction of $A^{\otimes n}$ to $\wedge^{n}(H)$. If $A \in \mathscr{L}_{1}(H)$, then $\wedge^{n}(A) \in \mathscr{L}_{1}\left(\wedge^{n}(H)\right)$ and

$$
\begin{equation*}
\left\|\wedge^{n}(A)\right\|_{1} \leq \frac{1}{n!}\|A\|_{1}^{n} \tag{5}
\end{equation*}
$$

The Fredholm determinant is then defined by

$$
\begin{equation*}
\operatorname{Det}(1+A)=1+\sum_{n=1}^{\infty} \operatorname{Tr}\left(\wedge^{n}(A)\right) \tag{6}
\end{equation*}
$$

The Fredholm determinant can be characterized as the unique function which is continuous in $A$ with respect to the trace norm $\|A\|_{1}$ and which coincides with the usual determinant when $A$ is a finite-dimensional operator.

One can extend the Fredholm determinant to a wider class on operators. Assume that we are given a splitting of $H$ into two subspaces:

$$
\begin{equation*}
H=H_{1} \oplus H_{2} \tag{7}
\end{equation*}
$$

According to this splitting, we write an operator $A \in \mathscr{L}(H)$ in block form,

$$
A=\left[\begin{array}{ll}
A_{11} & A_{21}  \tag{8}\\
A_{12} & A_{22}
\end{array}\right]
$$

where $A_{i j}: H_{j} \rightarrow H_{i}, i, j=1,2$. We define the even and odd parts of $A$ as follows:

$$
A_{\text {even }}:=\left[\begin{array}{cc}
A_{11} & 0 \\
0 & A_{22}
\end{array}\right], \quad A_{\text {odd }}:=\left[\begin{array}{cc}
0 & A_{21} \\
A_{12} & 0
\end{array}\right]
$$

We denote by $\mathscr{L}_{1 \mid 2}(H)$ the set of all operators $A \in \mathscr{L}(H)$ such that $A_{\text {even }} \in$ $\mathscr{L}_{1}(H)$ and $A_{\text {odd }} \in \mathscr{L}_{2}(H)$. Here $\mathscr{L}_{2}(H)$ denotes the space of all Hilbert-Schmidt operators on $H$, equipped with the norm

$$
\|A\|_{2}=\left(\sum_{n=1}^{\infty}\left\|A e_{n}\right\|^{2}\right)^{1 / 2}
$$

where $\left\{e_{n}\right\}_{n=1}^{\infty}$ is an orthonormal basis of $H$. Since $\mathscr{L}_{1}(H) \subset \mathscr{L}_{2}(H)$, one concludes that

$$
\mathscr{L}_{1}(H) \subset \mathscr{L}_{1 \mid 2}(H) \subset \mathscr{L}_{2}(H)
$$

We endow $\mathscr{L}_{1 \mid 2}(H)$ with the topology induced by the trace norm on the even part and by the Hilbert-Schmidt norm on the odd part.

Proposition 3 ([4]). The function $A \mapsto \operatorname{Det}(1+A)$ admits a unique extension to $\mathscr{L}_{1 \mid 2}(H)$ which is continuous in the topology of $\mathscr{L}_{1 \mid 2}(H)$. This extension is given by the formula

$$
\begin{equation*}
\operatorname{Det}(1+A)=\operatorname{Det}\left((1+A) e^{-A}\right) \cdot e^{\operatorname{Tr}\left(A_{\mathrm{even}}\right)} \tag{9}
\end{equation*}
$$

REmark 1. Note that, for each $A \in \mathscr{L}_{2}(H),(1+A) e^{-A}-1 \in \mathscr{L}_{1}(H)$. Therefore, $\operatorname{Det}\left((1+A) e^{-A}\right)$ is a classical Fredholm determinant.

REMARK 2. It should be noted that a possibility of extension of the Fredholm determinant to $\mathscr{L}_{1 \mid 2}(H)$ was already known to Berezin in the 1960s; see [3], page 8 .

We will now give another useful representation of $\operatorname{Det}(1+A)$ for $A \in \mathscr{L}_{1 \mid 2}(H)$.
Proposition 4. Let $A \in \mathscr{L}_{1 \mid 2}(H)$ have a block form (8). Assume that $\left\|A_{11}\right\|<1$. Then

$$
\begin{equation*}
\operatorname{Det}(1+A)=\operatorname{Det}\left(1+A_{11}\right) \cdot \operatorname{Det}\left(1+A_{22}-A_{21}\left(1+A_{11}\right)^{-1} A_{12}\right) \tag{10}
\end{equation*}
$$

[On the right-hand side of formula (10), both factors are classical Fredholm determinants, as both operators $A_{11}$ and $A_{22}-A_{21}\left(1+A_{11}\right)^{-1} A_{12}$ belong to $\mathscr{L}_{1}(H)$.]

REMARK 3. It should be stressed that the inequality $\left\|A_{11}\right\|<1$ can be achieved by every operator in $\mathscr{L}_{1 \mid 2}(H)$. More generally, for each fixed $\varepsilon>0$, we can always assume that $\left\|A_{11}\right\|<\varepsilon$. Indeed, assume $\left\|A_{11}\right\| \geq \varepsilon$. By the canonical decomposition of a compact (in particular, trace-class) operator (e.g., [20], Theorems 1.1 and 1.2), there exists an orthogonal splitting $H_{1}=H_{1}^{\prime} \oplus R$ such that the operator $A_{11}$ acts invariantly in both subspaces $H_{1}^{\prime}$ and $R$, the subspace $R$ is finite-dimensional, and the norm of the operator $A_{11}$ in the space $H_{1}^{\prime}$ is strictly less than $\varepsilon$. Setting $H_{2}^{\prime}:=H_{2} \oplus R$, we get a new orthogonal splitting $H=H_{1}^{\prime} \oplus H_{2}^{\prime}$. Write the operator $A$ in the block form with respect to this new splitting of $H$. Since $R$ is a finite-dimensional space, the even part of $A$ in the new splitting is still a trace-class operator, while the odd part of $A$ in the new splitting is still a Hilbert-Schmidt operator.

Proof of Proposition 4. For $i=1,2$, let $\left\{P_{i}^{(n)}\right\}_{n=1}^{\infty}$ be an ascending sequence of finite-dimensional orthogonal projections in $H_{i}$ such that $P_{i}^{(n)}$ strongly converges to the identity operator in $H_{i}$ as $n \rightarrow \infty$. Set $P^{(n)}:=P_{1}^{(n)}+P_{2}^{(n)}, n \in \mathbb{N}$. Then, for each $n \in \mathbb{N}, A^{(n)}:=P^{(n)} A P^{(n)}$ is a finite-dimensional operator in $H$, and

$$
\left\|A^{(n)}-A\right\|_{1 \mid 2} \rightarrow 0 \quad \text { as } n \rightarrow \infty
$$

Hence, by Proposition 3,

$$
\begin{equation*}
\operatorname{Det}(1+A)=\lim _{n \rightarrow \infty} \operatorname{Det}\left(1+A^{(n)}\right) \tag{11}
\end{equation*}
$$

In the block form,

$$
A^{(n)}=\left[\begin{array}{ll}
A_{11}^{(n)} & A_{21}^{(n)}  \tag{12}\\
A_{12}^{(n)} & A_{22}^{(n)}
\end{array}\right],
$$

where $A_{i j}^{(n)}=P_{i}^{(n)} A_{i j} P_{j}^{(n)}, i, j=1,2$. For each $n \in \mathbb{N}$, the operator $A^{(n)}$ is finitedimensional, hence, $\operatorname{Det}\left(1+A^{(n)}\right)$ is a classical Fredholm determinant. Therefore,

$$
\operatorname{Det}\left(1+A^{(n)}\right)=\operatorname{Det}\left[\begin{array}{cc}
1+A_{11}^{(n)} & +A_{21}^{(n)}  \tag{13}\\
A_{12}^{(n)} & 1+A_{22}^{(n)}
\end{array}\right]
$$

the latter (in fact, usual) determinant refers to the finite-dimensional Hilbert space $P^{(n)} H$. Since $\left\|A_{11}\right\|<1$, we have $\left\|A_{11}^{(n)}\right\|<1$ for all $n$. Hence, $1+A_{11}^{(n)}$ is invertible in $P_{11}^{(n)} H$. Employing the well-known formula for the determinant of a block matrix, we get from (11) and (13)

$$
\begin{equation*}
\operatorname{Det}(1+A)=\lim _{n \rightarrow \infty} \operatorname{Det}\left(1+A_{11}^{(n)}\right) \cdot \operatorname{Det}\left(1+A_{22}^{(n)}-A_{21}^{(n)}\left(1+A_{11}^{(n)}\right)^{-1} A_{12}^{(n)}\right) \tag{14}
\end{equation*}
$$

We state that

$$
\begin{gather*}
\left\|A_{11}^{(n)}-A_{11}\right\|_{1} \rightarrow 0, \quad\left\|A_{22}^{(n)}-A_{22}\right\|_{1} \rightarrow 0,  \tag{15}\\
\left\|A_{21}^{(n)}\left(1-A_{11}^{(n)}\right)^{-1} A_{12}^{(n)}-A_{21}\left(1-A_{11}\right)^{-1} A_{12}\right\|_{1} \rightarrow 0 \tag{16}
\end{gather*}
$$

as $n \rightarrow \infty$. Formula (15) is evident. In view of the formula

$$
\|B C\|_{1} \leq\|B\|_{2}\|C\|_{2}, \quad B, C \in \mathscr{L}_{2}(H)
$$

(see, e.g., [20], Theorem 2.8), the proof of (16) is routine, so we skip it. Thus, (10) follows from (14)-(16).

We will now derive an analog of formula (6) for $A \in \mathscr{L}_{1 \mid 2}(H)$. As follows from the proof of [19], Theorem 2.4, for each $A \in \mathscr{L}_{1}(H)$, we have

$$
\begin{equation*}
\operatorname{Tr}\left(\wedge^{n}(A)\right)=\frac{1}{n!} \sum_{\xi \in S_{n}} \operatorname{sign}(\xi) \prod_{\eta \in \operatorname{Cycle}(\xi)} \operatorname{Tr}\left(A^{|\eta|}\right) \tag{17}
\end{equation*}
$$

Here $S_{n}$ denotes the set of all permutations of $1, \ldots, n$, the product in (17) is over all cycles $\eta$ in permutation $\xi$, and $|\eta|$ denotes the length of cycle $\eta$. For $A \in \mathscr{L}_{1}(H)$, we clearly have $\operatorname{Tr}(A)=\operatorname{Tr}\left(A_{\text {even }}\right)$. We further note that, for each $A \in \mathscr{L}_{2}(H)$, we have $A^{k} \in \mathscr{L}_{1}(H)$ for $k \geq 2$. Thus, for each $A \in \mathscr{L}_{1 \mid 2}(H)$, we set $C_{n}(A)$ to be equal to the right-hand side of (17) in which we set

$$
\operatorname{Tr}(A):=\operatorname{Tr}\left(A_{\text {even }}\right), \quad A \in \mathscr{L}_{1 \mid 2}(H)
$$

Hence, $C_{n}(A)$ is well defined for each $A \in \mathscr{L}_{1 \mid 2}(H)$, and $C_{n}(A)=\operatorname{Tr}\left(\wedge^{n}(A)\right)$ for each $A \in \mathscr{L}_{1}(H)$.

Proposition 5. For each $A \in \mathscr{L}_{1 \mid 2}(H)$, we have

$$
\begin{equation*}
\operatorname{Det}(1+A)=1+\sum_{n=1}^{\infty} C_{n}(A) \tag{18}
\end{equation*}
$$

Proof. We know that formula (18) holds for all $A \in \mathscr{L}_{1}(H)$. Next, for each $A \in \mathscr{L}_{2}(H)$,

$$
\left\|A^{k}\right\|_{1} \leq\|A\|^{k-2}\left\|A^{2}\right\|_{1} \leq\|A\|^{k-2}\|A\|_{2}^{2} \leq\|A\|_{2}^{k}, \quad k \geq 2 .
$$

Hence, by the definition of $C_{n}(A)$,

$$
\begin{equation*}
\left|C_{n}(A)\right| \leq\|A\|_{1 \mid 2}^{n}, \tag{19}
\end{equation*}
$$

where

$$
\|A\|_{1 \mid 2}:=\max \left\{\|A\|_{2},\left\|A_{\text {even }}\right\|_{1}\right\}
$$

[Note that $\|\cdot\|_{1 \mid 2}$ is a norm on $\mathscr{L}_{1 \mid 2}(H)$ which determines its topology.] Hence, if $\|A\|_{1 \mid 2}<1$, the series on the right-hand side of (18) converges absolutely. We fix any $A \in \mathscr{L}_{1 \mid 2}(H)$ with $\|A\|_{1 \mid 2}<1$. For $i=1$, 2 , let $\left\{P_{i}^{(k)}\right\}_{k=1}^{\infty}$ be an ascending sequence of finite-dimensional orthogonal projections as in the proof of Proposition 4. Then, for each $k \in \mathbb{N}, A^{(k)}:=P^{(k)} A P^{(k)}$ is a finite-dimensional operator in $H$, and

$$
\left\|A^{(k)}-A\right\|_{1 \mid 2} \rightarrow 0 \quad \text { as } k \rightarrow \infty
$$

Hence, by (19) and the dominated convergence theorem,

$$
\operatorname{Det}\left(1+A^{(k)}\right)=1+\sum_{n=1}^{\infty} C_{n}\left(A^{(k)}\right) \rightarrow 1+\sum_{n=1}^{\infty} C_{n}(A)
$$

Therefore, by Proposition 3, formula (18) holds in this case.
Now we fix an arbitrary $A \in \mathscr{L}_{1 \mid 2}(H)$. Then, by (19), the function

$$
z \mapsto 1+\sum_{n=1}^{\infty} C_{n}(z A)=1+\sum_{n=1}^{\infty} z^{n} C_{n}(A)
$$

is analytic on the set $\left\{z \in \mathbb{C}:|z|<\|A\|_{1 \mid 2}^{-1}\right\}$. Thus, by the uniqueness of analytic continuation, to prove the proposition, it suffices to show that the function

$$
\mathbb{C} \ni z \mapsto \operatorname{Det}(1+z A)
$$

is entire. But this can be easily deduced from Proposition 4 and Remark 3.
Let us now assume that $H=L^{2}(X, m)$, where $X$ is a locally compact Polish space and $m$ is a Radon measure on $(X, \mathcal{B}(X))$. We fix any $X_{1}, X_{2} \in \mathcal{B}(X)$ such that $X=X_{1} \sqcup X_{2}$. By setting $H_{i}:=L^{2}\left(X_{i}, m\right), i=1$, 2, we get a splitting (7) of $H$.

Proposition 6. Let $K \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right.$ ) be an integral operator with integral kernel $K(x, y)$ such that $\int_{X}|K(x, x)| m(d x)<\infty$ and

$$
\begin{equation*}
\operatorname{Tr}\left(K_{\text {even }}\right)=\int_{X} K(x, x) m(d x) \tag{20}
\end{equation*}
$$

Then

$$
\begin{equation*}
\operatorname{Det}(1+K)=1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{X^{n}} \operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right) \tag{21}
\end{equation*}
$$

Proof. For each $l=2,3, \ldots$, we have

$$
\begin{equation*}
\operatorname{Tr}\left(K^{l}\right)=\int_{X^{l}} K\left(x_{1}, x_{2}\right) K\left(x_{2}, x_{3}\right) \cdots K\left(x_{l}, x_{1}\right) m\left(d x_{1}\right) \cdots m\left(d x_{l}\right) \tag{22}
\end{equation*}
$$

Note that the integral in (22) is independent of the choice of a version of the integral kernel of $K$. Hence, by the definition of $C_{n}(K)$ and formulas (20) and (22), we conclude that

$$
C_{n}(K)=\frac{1}{n!} \int_{X^{n}} \operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)
$$

Now formula (21) follows from Proposition 5.
3. $\boldsymbol{J}$-self-adjoint operators. We again assume that a Hilbert space $H$ is split into two subspaces; see (7). According to this splitting, we write a vector $f \in H$ as $f=\left(f_{2}, f_{2}\right)$ and an operator $A \in \mathscr{L}(H)$ in the block form (8). Denote by $P_{1}$ and $P_{2}$ the orthogonal projections of the Hilbert space $H$ onto $H_{1}$ and $H_{2}$, respectively. Setting $J:=P_{1}-P_{2}$, we introduce an (indefinite) $J$-scalar product on $H$ by

$$
[f, g]:=(J f, g)=\left(f_{1}, g_{1}\right)-\left(f_{2}, g_{2}\right), \quad f, g \in H
$$

An operator $A \in \mathscr{L}(H)$ is called self-adjoint in the indefinite scalar product $[\cdot, \cdot]$, or $J$-self-adjoint, if

$$
[A f, g]=[f, A g], \quad f, g \in H
$$

see, for example, [2]. In terms of the block form (8), an operator $A \in \mathscr{L}(H)$ is $J$-self-adjoint if and only if

$$
\begin{equation*}
A_{11}^{*}=A_{11}, \quad A_{22}^{*}=A_{22}, \quad A_{21}^{*}=-A_{12} \tag{23}
\end{equation*}
$$

REmARK 4. Assume $A$ is a usual matrix which has a block form (8). If the blocks of $A$ satisfy (23), then we will call $A$ a $J$-Hermitian matrix.

For any $A \in \mathscr{L}(H)$, we denote by $\widehat{A}$ the operator from $\mathscr{L}(H)$ given by

$$
\begin{equation*}
\widehat{A}:=A P_{1}+(1-A) P_{2} \tag{24}
\end{equation*}
$$

or, equivalently, in the block form,

$$
\widehat{A}=\left[\begin{array}{cc}
A_{11} & A_{21} \\
-A_{12} & 1-A_{22}
\end{array}\right]
$$

Clearly, if the operator $A$ is self-adjoint, then $\widehat{A}$ is $J$-self-adjoint, while if $A$ is $J$-self-adjoint, then $\widehat{A}$ is self-adjoint. Also $\widehat{\widehat{A}}=A$.

We will use below the following results.
Lemma 1. Let $A \in \mathscr{L}(H)$ be $J$-self-adjoint. Then $\|A\|=\left\|\widehat{A}-P_{2}\right\|$.
Proof. We have $A=\widehat{A} P_{1}+(1-\widehat{A}) P_{2}$, hence,

$$
A^{*}=P_{1} \widehat{A}+P_{2}(1-\widehat{A})
$$

Denote by $B_{A^{*}}$ the quadratic form on $H$ with generator $A^{*}$. For any $f, g \in H$,

$$
\begin{aligned}
B_{A^{*}}(f, g)= & \left(A^{*} f, g\right) \\
= & \left(\widehat{A}_{11} f_{1}, g_{1}\right)+\left(\widehat{A}_{12} f_{2}, g_{1}\right) \\
& +\left(\left(1-\widehat{A}_{22}\right) f_{2}, g_{2}\right)+\left(-\widehat{A}_{21} f_{1}, g_{2}\right)
\end{aligned}
$$

Denote $\tilde{g}=\left(g_{1},-g_{2}\right)=\left(\tilde{g}_{1}, \tilde{g}_{2}\right)$. Then

$$
\begin{aligned}
B_{A^{*}}(f, g) & =\left(\widehat{A}_{11} f_{1}, \tilde{g}_{1}\right)_{H}+\left(\widehat{A}_{12} f_{2}, \tilde{g}_{1}\right)_{H}+\left(\left(\widehat{A}_{22}-1\right) f_{2}, \tilde{g}_{2}\right)_{H}+\left(\widehat{A}_{21} f_{1}, \tilde{g}_{2}\right) \\
& =(\widehat{A} f, \tilde{g})-\left(f_{2}, \tilde{g}\right) \\
& =\left(\left(\widehat{A}-P_{2}\right) f, \tilde{g}\right) \\
& =B_{\widehat{A}-P_{2}}(f, \tilde{g})
\end{aligned}
$$

From here

$$
\left\|\widehat{A}-P_{2}\right\|=\left\|B_{\widehat{A}-P_{2}}\right\|=\left\|B_{A^{*}}\right\|=\left\|A^{*}\right\|=\|A\|
$$

Proposition 7. Let $A \in \mathscr{L}(H)$ be $J$-self-adjoint and assume that $0 \leq \widehat{A} \leq$ 1. Then $\|A\| \leq 1$.

Proof. By Lemma 1, it suffices to show that $\left\|\widehat{A}-P_{2}\right\| \leq 1$. Note that $\widehat{A}-P_{2}$ is self-adjoint. For each $f \in H$,

$$
\left(\left(\widehat{A}-P_{2}\right) f, f\right)=(\widehat{A} f, f)-\left(f_{2}, f_{2}\right) \leq(\widehat{A} f, f) \leq(f, f)
$$

Hence, $\widehat{A}-P_{2} \leq 1$. Next,

$$
\left(\left(\widehat{A}-P_{2}\right) f, f\right)=(\widehat{A} f, f)-\left(f_{2}, f_{2}\right) \geq-\left(f_{2}, f_{2}\right) \geq-(f, f)
$$

and so $\widehat{A}-P_{2} \geq-1$. Thus, $-1 \leq \widehat{A}-P_{2} \leq 1$, which implies the statement.
Proposition 8. Let $A \in \mathscr{L}(H)$ be $J$-self-adjoint and assume that $0 \leq \widehat{A} \leq$ 1. Then $\|A\|=1$ if and only if $\left\|A_{\text {even }}\right\|=1$.

Proof. By Lemma 1, it suffices to prove that $\left\|\widehat{A}-P_{2}\right\|=1$ if and only if $\left\|A_{\text {even }}\right\|=1$. Let us first assume that $\left\|\widehat{A}-P_{2}\right\|=1$.

Since $0 \leq \widehat{A} \leq 1$, we have $0 \leq \widehat{A}_{11} \leq 1$ and $0 \leq \widehat{A}_{22} \leq 1$. Hence, $0 \leq A_{11} \leq 1$ and $0 \leq A_{22} \leq 1$, and so $0 \leq A_{\text {even }} \leq 1$, which in turn implies that $\left\|A_{\text {even }}\right\| \leq 1$. We have to consider two cases.

Case 1. $-1 \in \sigma\left(\widehat{A}-P_{2}\right)$. [Here, $\sigma(B)$ denotes the spectrum of an operator $B \in \mathscr{L}(H)$.] Then there exists a sequence $\left(f^{(n)}\right)_{n=1}^{\infty}$ in $H$ such that $\left\|f^{(n)}\right\|=1$ and

$$
\left(\left(\widehat{A}-P_{2}\right) f^{(n)}, f^{(n)}\right) \rightarrow-1
$$

Since $\left(\widehat{A} f^{(n)}, f^{(n)}\right) \geq 0$ and $\left(P_{2} f^{(n)}, f^{(n)}\right) \leq 1$, we get

$$
\left(\widehat{A} f^{(n)}, f^{(n)}\right) \rightarrow 0, \quad\left\|f_{2}^{(n)}\right\| \rightarrow 1
$$

Hence, $f_{1}^{(n)} \rightarrow 0$. From here

$$
\left(\widehat{A}_{11} f_{1}^{(n)}, f_{1}^{(n)}\right)+\left(\widehat{A}_{21} f_{1}^{(n)}, f_{2}^{(n)}\right)+\left(\widehat{A}_{12} f_{2}^{(n)}, f_{1}^{(n)}\right) \rightarrow 0
$$

Thus,

$$
\left(\widehat{A}_{22} f_{2}^{(n)}, f_{2}^{(n)}\right) \rightarrow 0
$$

Hence,

$$
\left(\widehat{A}_{22} \frac{f_{2}^{(n)}}{\left\|f_{2}^{(n)}\right\|}, \frac{f_{2}^{(n)}}{\left\|f_{2}^{(n)}\right\|}\right) \rightarrow 0
$$

Hence, $0 \in \sigma\left(\widehat{A}_{22}\right)$, and so $1 \in \sigma\left(1-\widehat{A}_{22}\right)=\sigma\left(A_{22}\right)$.
Case 2. $1 \in \sigma\left(\widehat{A}-P_{2}\right)$. Then there exists a sequence $\left(f^{(n)}\right)_{n=1}^{\infty}$ in $H$ such that $\left\|f^{(n)}\right\|=1$ and

$$
\left(\left(\widehat{A}-P_{2}\right) f^{(n)}, f^{(n)}\right) \rightarrow 1
$$

Since $\left(\widehat{A} f^{(n)}, f^{(n)}\right) \leq 1$ and $\left(P_{2} f^{(n)}, f^{(n)}\right) \geq 0$, we get

$$
\left(\widehat{A} f^{(n)}, f^{(n)}\right) \rightarrow 1, \quad\left\|f_{2}^{(n)}\right\| \rightarrow 0
$$

From here, analogously to the above, we conclude that $1 \in \sigma\left(\widehat{A}_{11}\right)=\sigma\left(A_{11}\right)$.
Thus, in both cases, we get $\left\|A_{\text {even }}\right\|=1$. By inverting the arguments, we conclude the inverse statement.

Proposition 9. Let $A \in \mathscr{L}(H)$ be $J$-self-adjoint and let $A \in \mathscr{L}_{1 \mid 2}(H)$. Assume that $\|A\|<1$ and $A_{11} \geq 0$. Then $\operatorname{Det}(1-A)>0$.

Proof. Since $\|A\|<1$, we get $\left\|A_{11}\right\|<1$. Hence, by formula (10),

$$
\begin{align*}
\operatorname{Det}(1-A) & =\operatorname{Det}\left(1-A_{11}\right) \cdot \operatorname{Det}\left(1-A_{22}-A_{21}\left(1-A_{11}\right)^{-1} A_{12}\right) \\
& =\operatorname{Det}\left(1-A_{11}\right) \cdot \operatorname{Det}\left(1-A_{22}+A_{12}^{*}\left(1-A_{11}\right)^{-1} A_{12}\right) . \tag{25}
\end{align*}
$$

Note that both operators $-A_{11}$ and $-A_{22}+A_{12}^{*}\left(1-A_{11}\right)^{-1} A_{12}$ are trace-class and self-adjoint. Since $\left\|A_{11}\right\|<1$, we get $\operatorname{Det}\left(1-A_{11}\right)>0$. Further, $\left\|A_{22}\right\|<1$ and, hence, there exists $\varepsilon>0$ such that $1-A_{22} \geq \varepsilon 1$. Clearly, since $A_{11} \geq 0$,

$$
A_{12}^{*}\left(1-A_{11}\right)^{-1} A_{12} \geq 0
$$

which implies

$$
1-A_{22}+A_{12}^{*}\left(1-A_{11}\right)^{-1} A_{12} \geq \varepsilon 1
$$

From here

$$
\operatorname{Det}\left(1-A_{22}+A_{12}^{*}\left(1-A_{11}\right)^{-1} A_{12}\right)>0
$$

and the proposition is proven.
Proposition 10. Let $A \in \mathscr{L}_{1 \mid 2}(H)$ and let $A$ be $J$-self-adjoint. Let $0 \leq \widehat{A} \leq$ 1 and let $\|A\|<1$. Let $L:=A(1-A)^{-1}$. Then $L$ is $J$-self-adjoint, $L \in \mathscr{L}_{1 \mid 2}(H)$, and $L_{11} \geq 0, L_{22} \geq 0$.

Proof. We have $L=A+\sum_{n=2}^{\infty} A^{n}$, and

$$
\sum_{n=2}^{\infty}\left\|A^{n}\right\|_{1} \leq\|A\|_{2}^{2} \sum_{n=0}^{\infty}\|A\|^{n}<\infty
$$

Hence, $\sum_{n=2}^{\infty} A^{n} \in \mathscr{L}_{1}(H)$, so $L \in \mathscr{L}_{1 \mid 2}(H)$.
Let us show that the operator $L$ is $J$-self-adjoint. For any $f, g \in H$, we have

$$
(L f, g)=\sum_{n=1}^{\infty}\left(A^{n} f, g\right)=\sum_{n=1}^{\infty}\left(f,\left(A^{*}\right)^{n} g\right)=\sum_{n=1}^{\infty}\left(f,\left(A_{11}-A_{21}-A_{12}+A_{22}\right)^{n} g\right)
$$

Denoting $A_{11}^{\prime}:=A_{11}, A_{22}^{\prime}:=A_{22}, A_{12}^{\prime}:=-A_{12}, A_{21}^{\prime}:=-A_{21}$, we get

$$
\begin{equation*}
(L f, g)=\sum_{\substack{n=1}}^{\infty} \sum_{\substack{i_{k}, j_{k}=1,2 \\ k=1, \ldots, n}}\left(f, A_{i_{1} j_{1}}^{\prime} A_{i_{2} j_{2}}^{\prime} \cdots A_{i_{n} j_{n}}^{\prime} g\right) \tag{26}
\end{equation*}
$$

Assume that $f=f_{1} \in H_{1}, g=g_{1} \in H_{1}$. Then, in the latter sum, the terms, in which the number of the $A_{12}^{\prime}$ operators is not equal to the number of the $A_{21}^{\prime}$ operators, are equal to zero. Hence,

$$
\begin{align*}
\left(L_{11} f_{1}, g_{1}\right) & =\left(L f_{1}, g_{1}\right)=\sum_{\substack{n=1 \\
i_{k}, j_{k}=1,2 \\
k=1, \ldots, n}}^{\infty}\left(f_{1}, A_{i_{1} j_{1}} A_{i_{2} j_{2}} \cdots A_{i_{n} j_{n}} g_{1}\right)  \tag{27}\\
& =\sum_{n=1}^{\infty}\left(f_{1}, A^{n} g_{1}\right)=\left(f_{1}, L g_{1}\right)=\left(f_{1}, L_{11} g_{1}\right) .
\end{align*}
$$

Thus, $L_{11}^{*}=L_{11}$. Analogously, $L_{22}^{*}=L_{22}$.
In the case where $f=f_{1} \in H_{1}$ and $g=g_{2} \in H_{2}$, those terms in the sum in (26), in which the number of the $A_{21}^{\prime}$ operators is not equal to the number of the $A_{12}^{\prime}$ operators plus one, are equal to zero. Hence, similar to (27), we get

$$
\left(L_{21} f_{1}, g_{2}\right)=\left(f_{1},-L_{12} g_{2}\right)
$$

so $L_{21}^{*}=-L_{12}$. Thus, $L$ is $J$-self-adjoint.
Next, we will show that $L_{11} \geq 0$. Analogously to the proofs of Propositions 4 and 5, we define operators $A^{(n)}, n \in \mathbb{N}$. Thus, each $A^{(n)}$ is $J$-self-adjoint and

$$
\begin{equation*}
\left\|A^{(n)}\right\| \leq\|A\|<1 \tag{28}
\end{equation*}
$$

Let $\widehat{A}^{(n)}$ denote the corresponding transformation of the operator $A^{(n)}$ in the Hilbert space $P^{(n)} H$. Recalling representation (12) of $A^{(n)}$, we thus get

$$
\widehat{A}^{(n)}=\left[\begin{array}{cc}
A_{11}^{(n)} & A_{21}^{(n)} \\
-A_{12}^{(n)} & P_{2}^{(n)}-A_{22}^{(n)}
\end{array}\right]=P^{(n)}\left[\begin{array}{cc}
A_{11} & A_{21} \\
-A_{12} & 1-A_{22}
\end{array}\right] P^{(n)}=P^{(n)} \widehat{A} P^{(n)} .
$$

Since $0 \leq \widehat{A} \leq 1$, we therefore conclude that $0 \leq \widehat{A}^{(n)} \leq 1$. In particular, $A_{11}^{(n)} \geq 0$.
We may assume that the dimension of the Hilbert space $P^{(n)} H$ is $n$. Choose an orthonormal basis $\left(e^{(i)}\right)_{i=1, \ldots, n}$ of $P^{(n)} H$ such that $e^{(i)} \in P_{1}^{(n)} H, i=1, \ldots, k$, and $e^{(i)} \in P_{2}^{(n)} H, i=k+1, \ldots, n$. In terms of this orthonormal basis, we may treat the operator $A^{(n)}$ in $P^{(n)} H$ as an $n \times n J$-Hermitian matrix $\left[A_{i j}^{(n)}\right]_{i, j=1, \ldots, n}$. Let $X^{(n)}:=\{1,2, \ldots, n\}, \quad X_{1}^{(n)}:=\{1,2, \ldots, k\}, \quad X_{2}^{(n)}:=\{k+1, k+2, \ldots, n\}$, so that $X^{(n)}=X_{1}^{(n)} \sqcup X_{2}^{(n)}$. In view of Proposition 2, there exists a determinantal point process $\mu^{(n)}$ on $\Gamma_{X^{(n)}}$ with correlation kernel

$$
K^{(n)}(i, j):=A_{i j}^{(n)}, \quad i, j=1, \ldots, n .
$$

By Proposition 9, we have $\operatorname{det}\left(1-A^{(n)}\right)>0$. Let

$$
L^{(n)}:=A^{(n)}\left(1-A^{(n)}\right)^{-1}
$$

We define a possibly signed measure $\rho^{(n)}$ on the configuration space $\Gamma_{X^{(n)}}$ by setting

$$
\rho^{(n)}(\{\varnothing\}):=\operatorname{Det}\left(1-A^{(n)}\right)
$$

and for each nonempty configuration $\left\{i_{1}, i_{2}, \ldots, i_{m}\right\} \in \Gamma_{X^{(n)}}$,

$$
\rho^{(n)}\left(\left\{i_{1}, i_{2}, \ldots, i_{m}\right\}\right):=\operatorname{det}\left(1-A^{(n)}\right) \cdot \operatorname{det}\left(L^{(n)}\left(i_{u}, i_{v}\right)\right)_{u, v=1, \ldots, m}
$$

Analogously to the proof of Theorem 2 below, we may show that $\rho^{(n)}=\mu^{(n)}$. Hence, for each nonempty configuration $\left\{i_{1}, i_{2}, \ldots, i_{m}\right\} \in \Gamma_{X^{(n)}}$,

$$
\operatorname{det}\left(L^{(n)}\left(i_{u}, i_{v}\right)\right)_{u, v=1, \ldots, m} \geq 0
$$

In particular, for any nonempty configuration $\left\{i_{1}, i_{2}, \ldots, i_{m}\right\} \in \Gamma_{X_{1}^{(n)}}$,

$$
\operatorname{det}\left(L_{11}^{(n)}\left(i_{u}, i_{v}\right)\right)_{u, v=1, \ldots, m} \geq 0
$$

Hence, by the Sylvester criterion, $L_{11}^{(n)} \geq 0$, and so

$$
\begin{equation*}
\left(L_{11}^{(n)} f_{1}, f_{1}\right) \geq 0, \quad f_{1} \in H_{1} \tag{29}
\end{equation*}
$$

By (28) and the dominated convergence theorem, we get

$$
\begin{align*}
\lim _{n \rightarrow \infty}\left(L_{11}^{(n)} f_{1}, f_{1}\right) & =\lim _{n \rightarrow \infty}\left(L^{(n)} f_{1}, f_{1}\right) \\
& =\lim _{n \rightarrow \infty} \sum_{l=1}^{\infty}\left(\left(A^{(n)}\right)^{l} f_{1}, f_{1}\right)  \tag{30}\\
& =\sum_{l=1}^{\infty} \lim _{n \rightarrow \infty}\left(\left(A^{(n)}\right)^{l} f_{1}, f_{1}\right) \\
& =\sum_{l=1}^{\infty}\left(A^{l} f_{1}, f_{1}\right)=\left(L f_{1}, f_{1}\right)=\left(L_{11} f_{1}, f_{1}\right)
\end{align*}
$$

Thus, by (29) and (30), $\left(L_{11} f_{1}, f_{1}\right) \geq 0$ for all $f_{1} \in H_{1}$. Analogously, we get $L_{22} \geq 0$.

The following statement about $J$-Hermitian matrices was proven in [17].
Proposition 11 ([17]). Assume that $A$ is a J-Hermitian matrix and assume that its diagonal blocks, $A_{11}, A_{22}$, are nonnegative definite. Then $\operatorname{det}(A) \geq 0$.

REmark 5. Note that the arguments we used in the proof of Proposition 9 are similar to the arguments Olshanski [17] used to prove Proposition 11.

From now on we will again assume that $H=L^{2}(X, m)$, where $X$ is a locally compact Polish space, $m$ is a Radon measure on $(X, \mathcal{B}(X))$, and $X_{1}, X_{2} \in \mathcal{B}(X)$ are such that $X=X_{1} \sqcup X_{2}$. We also set $H_{i}:=L^{2}\left(X_{i}, m\right), i=1,2$. We further define

$$
\mathcal{B}\left(X_{i}\right):=\left\{\Lambda \in \mathcal{B}(X) \mid \Lambda \subset X_{i}\right\}
$$

and, analogously, $\mathcal{B}_{0}\left(X_{i}\right)$, for $i=1,2$.
For $\Delta \in \mathcal{B}_{0}(X)$, we denote by $P^{\Delta}$ the orthogonal projection of $L^{2}(X, m)$ onto $L^{2}(\Delta, m)$, that is, the operator of multiplication by $\chi_{\Delta}$. For an operator $K \in \mathscr{L}\left(L^{2}(X, m)\right)$, we denote $K^{\Delta}:=P^{\Delta} K P^{\Delta}$. We will say that an operator $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ is locally trace-class on $X_{1}$ and $X_{2}$ if, for each $\Delta \in \mathcal{B}_{0}\left(X_{i}\right)$, $i=1,2$, we have $K^{\Delta} \in \mathscr{L}_{1}\left(L^{2}(X, m)\right)$.

Proposition 12. Let $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ be $J$-self-adjoint and a locally trace-class operator on $X_{1}$ and $X_{2}$, and let $0 \leq \widehat{K} \leq 1$. Then, for each $\Delta \in \mathcal{B}_{0}(X)$, $K^{\Delta} \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$.

Proof. For each $\Delta_{1} \in \mathcal{B}_{0}\left(X_{1}\right)$, we have

$$
\begin{equation*}
P^{\Delta_{1}} \widehat{K} P^{\Delta_{1}}=K^{\Delta_{1}} \in \mathscr{L}_{1}\left(L^{2}(X, m)\right) \tag{31}
\end{equation*}
$$

Since $\widehat{K} \geq 0$, we get $P^{\Delta_{1}} \widehat{K} P^{\Delta_{1}} \geq 0$. Hence, by (31), $\sqrt{\widehat{K}} P^{\Delta_{1}} \in \mathscr{L}_{2}\left(L^{2}(X, m)\right)$. Next, for each $\Delta_{2} \in \mathcal{B}_{0}\left(X_{2}\right)$,

$$
P^{\Delta_{2}}(1-\widehat{K}) P^{\Delta_{2}}=K^{\Delta_{2}} \in \mathscr{L}_{1}\left(L^{2}(X, m)\right)
$$

Hence, analogously to the above, $\sqrt{1-\widehat{K}} P^{\Delta_{2}} \in \mathscr{L}_{2}\left(L^{2}(X, m)\right)$. From here

$$
\begin{aligned}
& K P^{\Delta_{1}}=\widehat{K} P^{\Delta_{1}}=\sqrt{\widehat{K}} \sqrt{\widehat{K}} P^{\Delta_{1}} \in \mathscr{L}_{2}\left(L^{2}(X, m)\right), \\
& K P^{\Delta_{2}}=(1-\widehat{K}) P^{\Delta_{2}}=\sqrt{1-\widehat{K}} \sqrt{1-\widehat{K}} P^{\Delta_{2}} \in \mathscr{L}_{2}\left(L^{2}(X, m)\right) .
\end{aligned}
$$

Therefore, $K\left(P^{\Delta_{1}}+P^{\Delta_{2}}\right) \in \mathscr{L}_{2}\left(L^{2}(X, m)\right)$. Thus, for each $\Delta \in \mathcal{B}_{0}(X), K P^{\Delta} \in$ $\mathscr{L}_{2}\left(L^{2}(X, m)\right)$, and so $K^{\Delta} \in \mathscr{L}_{2}\left(L^{2}(X, m)\right)$.

By our assumption, for each $\Delta \in \mathcal{B}_{0}(X)$,

$$
\begin{aligned}
K_{\text {even }}^{\Delta} & =P^{\Delta} K_{\text {even }} P^{\Delta}=P^{\Delta_{1}} K_{11} P^{\Delta_{1}}+P^{\Delta_{2}} K_{22} P^{\Delta_{2}} \\
& =K^{\Delta_{1}}+K^{\Delta_{2}} \in \mathscr{L}_{1}\left(L^{2}(X, m)\right) .
\end{aligned}
$$

(Here $\Delta_{i}:=\Delta \cap X_{i}, i=1,2$.) Thus, $K^{\Delta} \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$.
Proposition 13. Let $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ be $J$-self-adjoint, let $K^{\Delta} \in$ $\mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$ for each $\Delta \in \mathcal{B}_{0}(X)$, and let $K_{11} \geq 0, K_{22} \geq 0$. Then $K$ is an integral operator and its integral kernel $K(x, y)$ can be chosen so that:
(i) The kernel $K(x, y)$ is $J$-Hermitian.
(ii) For $i=1,2$ and any $x_{1}, \ldots, x_{n} \in X_{i}(n \in \mathbb{N})$, the matrix

$$
\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}
$$

is nonnegative definite.
(iii) For each $\Delta \in \mathcal{B}_{0}(X)$,

$$
\begin{equation*}
\operatorname{Tr}\left(K_{\text {even }}^{\Delta}\right)=\int_{\Delta} K(x, x) m(d x) \tag{32}
\end{equation*}
$$

Proof. For any $\Delta_{1} \in \mathcal{B}_{0}\left(X_{1}\right)$ and $\Delta_{2} \in \mathcal{B}_{0}\left(X_{2}\right), P^{\Delta_{2}} K P^{\Delta_{1}}$ is a HilbertSchmidt operator, hence an integral operator. Therefore, we can choose an integral kernel of $K_{21}$, which is a function $K_{21}(x, y)$ on $X_{2} \times X_{1}$. We now define an integral kernel $K_{12}(x, y)$ of the operator $K_{12}$ by setting $K_{12}(x, y):=-\overline{K_{21}(y, x)}$ for $(x, y) \in X_{1} \times X_{2}$. Next, the operators $K_{11}$ and $K_{22}$ are nonnegative, locally traceclass operators. Hence, we can choose their integral kernels according to [10], Lemma A.3; see also [14], Section 3. By combining the integral kernels $K_{i j}(x, y)$, $i, j=1,2$, we obtain an integral kernel $K(x, y)$ of $K$ with needed properties.

From now on, for an operator $K$ as in Proposition 13, we will always assume that its integral kernel satisfies statements (i)-(iii) of this proposition.

We denote by $B_{0}(X)$ the space of all measurable bounded real-valued functions on $X$ with compact support. For each $\varphi \in B_{0}(X)$, we preserve the notation $\varphi$ for the bounded linear operator of multiplication by $\varphi$ in $L^{2}(X, m)$.

Proposition 14. Let $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ be $J$-self-adjoint, let $K_{11} \geq 0$, $K_{22} \geq 0$, and let $K^{\Delta} \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$ for each $\Delta \in \mathcal{B}_{0}(X)$. Fix any $\Delta \in \mathcal{B}_{0}(X)$ and any $\varphi \in B_{0}(X)$ which vanishes outside $\Delta$. Then $K^{\Delta} \varphi, \operatorname{sgn}(\varphi) \sqrt{|\varphi|} K \sqrt{|\varphi|} \in$ $\mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$ and
$\operatorname{Det}\left(1+K^{\Delta} \varphi\right)=\operatorname{Det}(1+\operatorname{sgn}(\varphi) \sqrt{|\varphi|} K \sqrt{|\varphi|})$

$$
\begin{align*}
=1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{X^{n}} & \varphi\left(x_{1}\right) \cdots \varphi\left(x_{n}\right)  \tag{33}\\
& \times \operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)
\end{align*}
$$

Proof. Since $K^{\Delta} \in \mathscr{L}_{2}\left(L^{2}(X, m)\right), K^{\Delta} \varphi \in \mathscr{L}_{2}\left(L^{2}(X, m)\right)$. Since $K_{\text {even }}^{\Delta} \in$ $\mathscr{L}_{1}\left(L^{2}(X, m)\right)$,

$$
\left(K^{\Delta} \varphi\right)_{\mathrm{even}}=K_{\text {even }}^{\Delta} \varphi \in \mathscr{L}_{1}\left(L^{2}(X, m)\right)
$$

Thus, $K^{\Delta} \varphi \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$.
Denote $\psi_{1}:=\operatorname{sgn}(\varphi) \sqrt{|\varphi|}$ and $\psi_{2}:=\sqrt{|\varphi|}, \psi_{1}, \psi_{2} \in B_{0}(X)$. Since $\psi_{1}$ and $\psi_{2}$ vanish outside $\Delta$, we get

$$
\psi_{1} K \psi_{2}=\psi_{1} K^{\Delta} \psi_{2}
$$

and, analogously to the above, we conclude that $\psi_{1} K \psi_{2} \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$.
Since $K_{\text {even }}^{\Delta} \in \mathscr{L}_{1}\left(L^{2}(X, m)\right)$ and since $\psi_{1}, \psi_{2} \in \mathscr{L}\left(L^{2}(X, m)\right)$, by (4),

$$
\begin{align*}
\operatorname{Tr}\left(\left(K^{\Delta} \varphi\right)_{\text {even }}\right) & =\operatorname{Tr}\left(K_{\text {even }}^{\Delta} \psi_{2} \psi_{1}\right)=\operatorname{Tr}\left(\psi_{1} K_{\text {even }}^{\Delta} \psi_{2}\right)  \tag{34}\\
& =\operatorname{Tr}\left(\psi_{1} K_{\text {even }} \psi_{2}\right)=\operatorname{Tr}\left(\left(\psi_{1} K \psi_{2}\right)_{\text {even }}\right)
\end{align*}
$$

Next, for $l=2,3, \ldots$,

$$
\begin{align*}
\operatorname{Tr}\left(\left(\psi_{1} K \psi_{2}\right)^{l}\right) & =\operatorname{Tr}\left(\psi_{1} K^{\Delta} \varphi K^{\Delta} \varphi \cdots K^{\Delta} \varphi K^{\Delta} \psi_{2}\right) \\
& =\operatorname{Tr}\left(K^{\Delta} \varphi K^{\Delta} \varphi \cdots K^{\Delta} \varphi K^{\Delta} \psi_{2} \psi_{1}\right)=\operatorname{Tr}\left(\left(K^{\Delta} \varphi\right)^{l}\right) . \tag{35}
\end{align*}
$$

By (34) and (35), $C_{n}\left(K^{\Delta} \varphi\right)=C_{n}\left(\psi_{1} K \psi_{2}\right)$ for each $n \in \mathbb{N}$, hence, formula (33) holds.

Next, we note that the integral kernel $K^{\Delta}(x, y)$ of the operator $K^{\Delta}$ is the restriction of $K(x, y)$ to $\Delta^{2}$. Clearly, the integral kernel of $K^{\Delta} \varphi$ is $K^{\Delta}(x, y) \varphi(y)$. Using (32), it is not hard to show that

$$
\operatorname{Tr}\left(\left(K^{\Delta} \varphi\right)_{\mathrm{even}}\right)=\int_{X} K^{\Delta}(x, x) \varphi(x) m(d x)
$$

Hence, by Proposition 6,

$$
\begin{aligned}
& \operatorname{Det}\left(1+K^{\Delta} \varphi\right) \\
& \begin{aligned}
&= 1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{X^{n}} \operatorname{det}\left[K^{\Delta}\left(x_{i}, x_{j}\right) \varphi\left(x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right) \\
&= 1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{X^{n}} \varphi\left(x_{1}\right) \cdots \varphi\left(x_{n}\right) \\
& \quad \times \operatorname{det}\left[K^{\Delta}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right) \\
&=1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{X^{n}} \varphi\left(x_{1}\right) \cdots \varphi\left(x_{n}\right) \\
& \quad \times \operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right) .
\end{aligned}
\end{aligned}
$$

4. Main results. We again assume that $X$ is a locally compact Polish space and $m$ is a Radon measure on $(X, \mathcal{B}(X))$. We will also assume that $m$ takes a positive value on each open nonempty set in $X$. Let $\Gamma=\Gamma_{X}$ be the configuration space over $X$. Let $\mu$ be a point process on $X$, that is, a probability measure on $(\Gamma, \mathcal{B}(\Gamma))$. Assume that $\mu$ satisfies

$$
\begin{equation*}
\int_{\Gamma} C^{|\gamma \cap \Delta|} \mu(d \gamma) \quad \text { for all } \Delta \in \mathcal{B}_{0}(X) \text { and all } C>0 \tag{36}
\end{equation*}
$$

Then the Bogoliubov functional of $\mu$ is defined as

$$
\begin{equation*}
B_{\mu}(\varphi):=\int_{\Gamma} \prod_{x \in \gamma}(1+\varphi(x)) \mu(d \gamma), \quad \varphi \in B_{0}(X) \tag{37}
\end{equation*}
$$

Note that, since the function $\varphi$ has compact support, only a finite number of terms in the product $\prod_{x \in \gamma}(1+\varphi(x))$ are not equal to one. Note also that the integrability of the function $\prod_{x \in \gamma}(1+\varphi(x))$ for each $\varphi \in B_{0}(X)$ is equivalent to condition (36). If a point process $\mu$ has correlation functions $\left(k_{\mu}^{(n)}\right)_{n=1}^{\infty}$ [see (1)], then condition (36) is also equivalent to

$$
\sum_{n=1}^{\infty} \frac{C^{n}}{n!} \int_{\Delta^{n}} k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right) m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)<\infty
$$

$$
\text { for all } \Delta \in \mathcal{B}_{0}(X) \text { and all } C>0
$$

and the Bogoliubov functional of $\mu$ is given by

$$
\begin{equation*}
B_{\mu}(\varphi)=1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{X^{n}} \varphi\left(x_{1}\right) \cdots \varphi\left(x_{n}\right) k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right) m\left(d x_{1}\right) \cdots m\left(d x_{n}\right) \tag{38}
\end{equation*}
$$

for each $\varphi \in B_{0}(X)$. The Bogoliubov functional of $\mu$ uniquely determines this point process. For more detail about the Bogoliubov functional see, for example, [12].

Let us now briefly recall some known facts about configuration spaces and point processes; see, for example, $[9,16]$ for further details. The $\sigma$-algebra $\mathcal{B}(\Gamma)$ coincides with the minimal $\sigma$-algebra on $\Gamma$ with respect to which all mappings of the form $\Gamma \ni \gamma \mapsto|\gamma \cap \Lambda|$ with $\Lambda \in \mathcal{B}_{0}(X)$ are measurable. For a fixed set $\Delta \in \mathcal{B}(X)$, we denote by $\mathcal{B}_{\Delta}(\Gamma)$ the minimal $\sigma$-algebra on $\Gamma$ with respect to which all mappings of the form $\Gamma \ni \gamma \mapsto|\gamma \cap \Lambda|$ with $\Lambda \in \mathcal{B}_{0}(X), \Lambda \subset \Delta$, are measurable. In particular, $\mathcal{B}_{\Delta}(\Gamma)$ is a sub- $\sigma$-algebra of $\mathcal{B}(\Gamma)$. The $\sigma$-algebras $\mathcal{B}\left(\Gamma_{\Delta}\right)$ and $\mathcal{B}_{\Delta}(\Gamma)$ can be identified in the sense that, for each $A \in \mathcal{B}\left(\Gamma_{\Delta}\right),\{\gamma \in \Gamma \mid \gamma \cap \Delta \in A\} \in$ $\mathcal{B}_{\Delta}(\Gamma)$ and each set from $\mathcal{B}_{\Delta}(\Gamma)$ has a unique such representation. Hence, the restriction of a point process $\mu$ on $X$ to the $\sigma$-algebra $\mathcal{B}_{\Delta}(\Gamma)$-denoted by $\mu_{\Delta}$-can be identified with a point process on $\Delta$, that is, with a probability measure on $\left(\Gamma_{\Delta}, \mathcal{B}\left(\Gamma_{\Delta}\right)\right)$.

Let $\Delta$ be a compact subset of $X$. Then the configuration space $\Gamma_{\Delta}$ consists of all finite configurations in $\Delta$, that is, $\Gamma_{\Delta}=\bigsqcup_{n=0}^{\infty} \Gamma_{\Delta}^{(n)}$, where $\Gamma_{\Delta}^{(0)}:=\{\varnothing\}$ and for $n \in \mathbb{N}, \Gamma_{\Delta}^{(n)}$ consists of all $n$-point configurations in $\Delta$. Denote

$$
\widetilde{\Delta}^{n}:=\left\{\left(x_{1}, \ldots, x_{n}\right) \in \Delta^{n} \mid x_{i} \neq x_{j} \text { if } i \neq j\right\} .
$$

Let $\mathcal{B}\left(\Gamma_{\Delta}^{(n)}\right)$ denote the image of the $\sigma$-algebra $\mathcal{B}\left(\widetilde{\Delta}^{n}\right)$ under the mapping

$$
\widetilde{\Delta}^{n} \ni\left(x_{1}, \ldots, x_{n}\right) \mapsto\left\{x_{1}, \ldots, x_{n}\right\} \in \Gamma_{\Delta}^{(n)} .
$$

Then $\mathcal{B}\left(\Gamma_{\Delta}\right)$ is the minimal $\sigma$-algebra on $\Gamma_{\Delta}$ which contains all $\mathcal{B}\left(\Gamma_{\Delta}^{(n)}\right), n \in \mathbb{N}$. A point process $\mu$ on $X$ has local densities in $\Delta$ if, for each $n \in \mathbb{N}$, there exists a nonnegative measurable symmetric function $d_{\mu}^{(n)}[\Delta]\left(x_{1}, \ldots, x_{n}\right)$ on $\widetilde{\Delta}^{n}$ such that

$$
\begin{aligned}
\int_{\Gamma_{\Delta}^{(n)}} & f^{(n)}(\gamma) \mu_{\Delta}(d \gamma) \\
& =\frac{1}{n!} \int_{\widetilde{\Delta}^{n}} f^{(n)}\left(\left\{x_{1}, \ldots, x_{n}\right\}\right) d_{\mu}^{(n)}[\Delta]\left(x_{1}, \ldots, x_{n}\right) m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)
\end{aligned}
$$

for each measurable function $f^{(n)}: \Gamma_{\Delta}^{(n)} \rightarrow[0, \infty)$. We also denote $d_{\mu}^{(0)}[\Delta]:=$ $\mu_{\Delta}(\{\varnothing\})$. In the case where $X=\Delta$ (so that $X$ is a compact Polish space), we will write $d_{\mu}^{(n)}$ instead of $d_{\mu}^{(n)}[\Delta]$.

THEOREM 2. Let $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ be $J$-self-adjoint. Let $K$ be a locally trace-class operator on $X_{1}$ and $X_{2}$, and let $0 \leq \widehat{K} \leq 1$. Then there exists a unique point process $\mu$ on $X$ which has correlation functions

$$
\begin{equation*}
k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} \tag{39}
\end{equation*}
$$

The Bogoliubov functional of $\mu$ is given by

$$
\begin{equation*}
B_{\mu}(\varphi)=\operatorname{Det}(1+\operatorname{sgn}(\varphi) \sqrt{|\varphi|} K \sqrt{|\varphi|}), \quad \varphi \in B_{0}(X) \tag{40}
\end{equation*}
$$

If additionally $\|K\|<1$, then for each $\Delta \in \mathcal{B}_{0}(X)$, the point process $\mu$ has local densities in $\Delta$ :

$$
\begin{align*}
d_{\mu}^{(0)}[\Delta] & =\operatorname{Det}\left(1-K^{\Delta}\right),  \tag{41}\\
d_{\mu}^{(n)}[\Delta]\left(x_{1}, \ldots, x_{n}\right) & =\operatorname{Det}\left(1-K^{\Delta}\right) \operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n},
\end{align*}
$$

where $L[\Delta]:=K^{\Delta}\left(1-K^{\Delta}\right)^{-1}$.
Proof. By Proposition $7,\|K\| \leq 1$. We first assume that $\|K\|<1$. We fix any compact $\Delta \subset X$. By Proposition $12, K^{\Delta} \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$, hence, $K^{\Delta} \in$ $\mathscr{L}_{1 \mid 2}\left(L^{2}(\Delta, m)\right)$. Clearly, $K^{\Delta}$ is $J$-self-adjoint. Setting $\Delta_{i}:=\Delta \cap X_{i}, i=1,2$, we get

$$
\begin{align*}
P^{\Delta} \widehat{K} P^{\Delta} & =P^{\Delta}\left(K P_{1}+(1-K) P_{2}\right) P^{\Delta} \\
& =K^{\Delta} P^{\Delta_{1}}+\left(1-K^{\Delta}\right) P^{\Delta_{2}}=\widehat{K}^{\Delta}, \tag{42}
\end{align*}
$$

where the latter operator is understood as the transformation (24) of the operator $K^{\Delta}$ in the Hilbert space $L^{2}(\Delta, m)=L^{2}\left(\Delta_{1}, m\right) \oplus L^{2}\left(\Delta_{2}, m\right)$. As $0 \leq \widehat{K} \leq 1$, we conclude from (42) that $0 \leq \widehat{K}^{\Delta} \leq 1$. Since $\|K\|<1$, we have $\left\|K^{\Delta}\right\|<1$. Hence, by Proposition $9, \operatorname{Det}\left(1-K^{\Delta}\right)>0$.

Furthermore, by Proposition 10, the operator $L[\Delta]$ is $J$-self-adjoint and

$$
L[\Delta] \in \mathscr{L}_{1 \mid 2}\left(L^{2}(\Delta, m)\right), \quad L[\Delta]_{11} \geq 0, \quad L[\Delta]_{22} \geq 0
$$

Hence, we can choose an integral kernel $L[\Delta](x, y)$ of the operator $L[\Delta]$ according to Proposition 13. Therefore, for any $x_{1}, \ldots, x_{n} \in \Delta_{1}, x_{n+1}, \ldots, x_{n+m} \in \Delta_{2}$, the matrix $\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n+m}$ is $J$-Hermitian and the diagonal blocks

$$
\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}, \quad\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=n+1, \ldots, n+m}
$$

are nonnegative definite. Hence, by Proposition 11,

$$
\operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n+m} \geq 0
$$

Therefore, for each $n \in \mathbb{N}$, the function

$$
\widetilde{\Delta}^{n} \ni\left(x_{1}, \ldots, x_{n}\right) \mapsto \operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}
$$

is symmetric and takes nonnegative values.
Hence, we can define a positive measure $\mu_{\Delta}$ on $\left(\Gamma_{\Delta}, \mathcal{B}\left(\Gamma_{\Delta}\right)\right)$ for which

$$
\begin{align*}
d_{\mu_{\Delta}}^{(0)} & =\operatorname{Det}\left(1-K^{\Delta}\right), \\
d_{\mu_{\Delta}}^{(n)}\left(x_{1}, \ldots, x_{n}\right) & =\operatorname{Det}\left(1-K^{\Delta}\right) \operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}, \quad n \in \mathbb{N} . \tag{43}
\end{align*}
$$

Note that

$$
\operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}=0 \quad \text { for all }\left(x_{1}, \ldots, x_{n}\right) \in \Delta^{n} \backslash \widetilde{\Delta}_{n}, \quad n \in \mathbb{N}
$$

Hence, the Bogoliubov functional of $\mu_{\Delta}$ is given by

$$
B_{\mu_{\Delta}}(\varphi)=\operatorname{Det}\left(1-K^{\Delta}\right)
$$

$$
\begin{equation*}
\times\left(1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{\Delta^{n}}\left(1+\varphi\left(x_{1}\right)\right) \cdots\left(1+\varphi\left(x_{n}\right)\right)\right. \tag{44}
\end{equation*}
$$

$$
\begin{array}{r}
\left.\times \operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)\right) \\
\varphi \in B(\Delta)
\end{array}
$$

Here $B(\Delta)$ denotes the set of all bounded measurable functions on $\Delta$. It follows from Proposition 14 and (44) that

$$
\begin{equation*}
B_{\mu_{\Delta}}(\varphi)=\operatorname{Det}\left(1-K^{\Delta}\right) \operatorname{Det}(1+L[\Delta](1+\varphi)), \quad \varphi \in B(\Delta) \tag{45}
\end{equation*}
$$

Hence, by [4], Corollary A.3, and Proposition 14,

$$
\begin{align*}
B_{\mu_{\Delta}}(\varphi) & =\operatorname{Det}\left(1-K^{\Delta}\right)(1+L[\Delta](1+\varphi)) \\
& =\operatorname{Det}\left(1+K^{\Delta} \varphi\right)  \tag{46}\\
& =\operatorname{Det}(1+\operatorname{sgn}(\varphi) \sqrt{|\varphi|} K \sqrt{|\varphi|}), \quad \varphi \in B(\Delta)
\end{align*}
$$

Now we take any sequence of compact subsets of $X,\left\{\Delta_{n}\right\}_{n=1}^{\infty}$, such that

$$
\Delta_{1} \subset \Delta_{2} \subset \cdots, \quad \bigcup_{n=1}^{\infty} \Delta_{n}=X
$$

By (46), the probability measures $\mu_{\Delta_{n}}$ on $\left(\Gamma, \mathcal{B}_{\Delta_{n}}(\Gamma)\right)$ form a consistent family of probability measures. Therefore, by the Kolmogorov theorem, there exists a unique probability measure on $(\Gamma, \mathcal{B}(\Gamma))$ such that the restriction of $\mu$ to each $\mathcal{B}_{\Delta_{n}}(\Gamma)$ coincides with $\mu_{\Delta_{n}}$. By (46), the Bogoliubov functional of $\mu$ is given by (40), while the statement about the local densities of $\mu$ follows from (43). The determinantal form of the correlation functions of $\mu$-formula (39)-follows from (38), (40) and Proposition 14.

Let us now consider the case where $\|K\|=1$. For each $\varepsilon \in(0,1)$, set $K_{\varepsilon}:=\varepsilon K$. Hence, $\left\|K_{\varepsilon}\right\|<1$. We have

$$
\begin{equation*}
\widehat{K}_{\varepsilon}=\varepsilon K P_{1}+(1-\varepsilon K) P_{2}=\varepsilon \widehat{K}+(1-\varepsilon) P_{2} \tag{47}
\end{equation*}
$$

Since $\widehat{K} \geq 0$ and $P_{2} \geq 0$, we get $\widehat{K}_{\varepsilon} \geq 0$, and since $\widehat{K} \leq 1$ and $P_{2} \leq 1$, we get $\widehat{K}_{\varepsilon} \leq 1$. Hence, by the proved above, there exists a point process $\mu_{\varepsilon}$ which has correlation functions

$$
\begin{equation*}
k_{\mu_{\varepsilon}}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\varepsilon^{n} \operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} \tag{48}
\end{equation*}
$$

Hence, the corresponding correlation measure is $\star$-positive definite in the sense of [11]; see also [14]. By taking the limit as $\varepsilon \rightarrow 0$, we therefore conclude that the functions

$$
\begin{equation*}
k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right):=\operatorname{det}\left[K\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}, \quad n \in \mathbb{N} \tag{49}
\end{equation*}
$$

determine a $\star$-positive definite correlation measure. By Proposition 14, for each $\Delta \in \mathcal{B}_{0}(X)$ and $C>0$,

$$
1+\sum_{n=1}^{\infty} \frac{C^{n}}{n!} \int_{\Delta^{n}} k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right) m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)=\operatorname{Det}\left(1+C K^{\Delta}\right)<\infty
$$

Hence, by [14], Corollary 1, we conclude that there exists a unique probability measure $\mu$ on ( $\Gamma, \mathcal{B}(\Gamma))$ which has correlation functions (49). By Proposition 14 and formula (38), the Bogoliubov functional of $\mu$ is given by (40).

The following corollary easily follows from Theorem 2 and Proposition 5.1 in [8] and its proof.

COROLLARY 1. Let $G: L^{2}\left(X_{1}, m\right) \rightarrow L^{2}\left(X_{2}, m\right)$ be a bounded linear operator such that, for any $\Delta_{1} \in \mathcal{B}_{0}\left(X_{1}\right)$ and $\Delta_{2} \in \mathcal{B}_{0}\left(X_{2}\right)$, the operators $G P^{\Delta_{1}}$ and $P^{\Delta_{2}} G$ are Hilbert-Schmidt. Let an operator $L \in \mathscr{L}\left(L^{2}(X, m)\right)$ be defined by

$$
L:=\left[\begin{array}{cc}
0 & G \\
-G^{*} & 0
\end{array}\right] .
$$

Then operator $1+L$ is invertible, and we set $K:=L(1+L)^{-1}$. We further have the following:
(i) The operator $K$ is $J$-self-adjoint.
(ii) The operator $K$ is locally trace-class on $X_{1}$ and $X_{2}$.
(iii) The operator $\widehat{K}$ is the orthogonal projection of $L^{2}(X, m)$ onto the subspace

$$
\left\{h \oplus G h \mid h \in L^{2}\left(X_{1}, m\right)\right\} .
$$

Thus, by Theorem 2, there exists a unique determiminantal point process with correlation kernel $K(x, y)$.

REmARK 6. As we mentioned in Section 1, the Whittaker kernel [6], the matrix tail kernel [17] and the continuous hypergeometric kernel [7] have their $L$ operators as in Corollary 1, and so their $\widehat{K}$ operators are orthogonal projections.

Proof of Corollary 1. That the operator $1+L$ is invertible is shown in [8], Section 5. Statement (iii) is just [8], Proposition 5.1. By the proof of Proposition 5.1 [8], the operator $L$ has the following block form:

$$
\begin{aligned}
& K_{11}=G^{*} G\left(1+G^{*} G\right)^{-1}, \\
& K_{22}=G G^{*}\left(1+G G^{*}\right)^{-1}, \\
& K_{21}=G\left(1+G^{*} G\right)^{-1}, \\
& K_{21}=-G^{*}\left(1+G^{*} G\right)^{-1} .
\end{aligned}
$$

Hence, statement (i) obviously follows. So we only need to prove statement (ii). To this end, we fix any $\Delta_{1} \in \mathcal{B}_{0}\left(X_{1}\right)$ and $\Delta_{2} \in \mathcal{B}_{0}\left(X_{2}\right)$. By the assumption of the corollary, $P^{\Delta_{2}} G$ is a Hilbert-Schmidt operator. Therefore,

$$
P^{\Delta_{2}} K_{21} P^{\Delta_{1}}=P^{\Delta_{2}} G\left(1+G^{*} G\right)^{-1} P^{\Delta_{1}}
$$

is a Hilbert-Schmidt operator, hence so is the operator $P^{\Delta_{1}} K_{12} P^{\Delta_{2}}$. Again by the assumption of the corollary, $G P^{\Delta_{1}}$ is a Hilbert-Schmidt operator, hence,

$$
\left(G P^{\Delta_{1}}\right)^{*}\left(G P^{\Delta_{1}}\right)=P^{\Delta_{1}} G^{*} G P^{\Delta_{1}}
$$

is a trace-class operator. Let $\left\{e^{(n)}\right\}_{n \geq 1}$ be an orthonormal basis in $L^{2}\left(\Delta_{1}, m\right)$. Then, by the spectral theorem,

$$
\begin{aligned}
\sum_{n \geq 1}\left(K_{11} e^{(n)}, e^{(n)}\right)_{L^{2}\left(\Delta_{1}, m\right)} & =\sum_{n \geq 1}\left(G^{*} G\left(1+G^{*} G\right)^{-1} e^{(n)}, e^{(n)}\right)_{L^{2}\left(\Delta_{1}, m\right)} \\
& \leq \sum_{n \geq 1}\left(G^{*} G e^{(n)}, e^{(n)}\right)_{L^{2}\left(\Delta_{1}, m\right)}<\infty
\end{aligned}
$$

Therefore, the operator $P^{\Delta_{1}} K_{11} P^{\Delta_{1}}$ is trace-class. Analogously, we may also show that the operator $P^{\Delta_{2}} K_{22} P^{\Delta_{2}}$ is trace-class. Thus, statement (ii) is proven.

Corollary 2. Let an operator $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ be $J$-self-adjoint and locally trace-class on $X_{1}$ and $X_{2}$. Let $0 \leq \widehat{K} \leq 1$ and let $\|K\|=1$. Let $\mu$ be the corresponding determinantal point process. Assume that $\Delta \in \mathcal{B}_{0}(X)$ is such that $\left\|K^{\Delta}\right\|=1$. Then

$$
\mu_{\Delta}(\{\varnothing\})=\operatorname{Det}\left(1-K^{\Delta}\right)=0,
$$

that is, the $\mu$ probability of the event that there are no particles in $\Delta$ is equal to zero.

PROOF. By (40), for each $\Delta \in \mathcal{B}_{0}(X)$ and $z>0$,

$$
\begin{aligned}
\int_{\Gamma} e^{-z|\gamma \cap \Delta|} \mu(d \gamma) & =\int_{\Gamma} \prod_{x \in \gamma}\left(1+\left(e^{-z}-1\right) \chi_{\Delta}\right) \mu(d \gamma) \\
& =\operatorname{Det}\left(1-\left(1-e^{-z}\right) K^{\Delta}\right)
\end{aligned}
$$

Letting $z \rightarrow \infty$ and using the dominated convergence theorem, we get

$$
\mu_{\Delta}(\{\varnothing\})=\operatorname{Det}\left(1-K^{\Delta}\right)
$$

Since $\left\|K^{\Delta}\right\|=1$, by Proposition 8 , at least one of the operators $K^{\Delta_{1}}=K_{11}^{\Delta_{1}}$, $K^{\Delta_{2}}=K_{22}^{\Delta_{2}}$ must have norm 1. (Here, as above, $\Delta_{i}=\Delta \cap X_{i}, i=1,2$.) Assume $\left\|K^{\Delta_{1}}\right\|=1$ (the other case is analogous). As $\operatorname{Det}\left(1-K^{\Delta_{1}}\right)$ is a classical Fredholm determinant and the operator $K^{\Delta_{1}}$ is self-adjoint, we get $\operatorname{Det}\left(1-K^{\Delta_{1}}\right)=0$. Thus, we have $\mu_{\Delta_{1}}(\{\varnothing\})=0$, that is, the $\mu$ probability of the event that there are no particles in the set $\Delta_{1}$ is equal to 0 . From here the statement follows.

REmARK 7. Note that, for a determinantal point process $\mu$ with a $J$-selfadjoint correlation operator $K$, the restriction of $\mu$ to the $\sigma$-algebra $\mathcal{B}_{X_{i}}(\Gamma)$ ( $i=1,2$ ) may be identified with the determinantal point process on $X_{i}$ whose correlation operator is the self-adjoint operator $K_{i i}$.

We will now show that the conditions on a $J$-self-adjoint operator $K$ in Theorem 2 are, in fact, necessary for a determinantal point process with correlation kernel $K(x, y)$ to exist.

Theorem 3. Let $K \in \mathscr{L}\left(L^{2}(X, m)\right)$ be $J$-self-adjoint, let $K_{11} \geq 0, K_{22} \geq 0$, and let $K^{\Delta} \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$ for each $\Delta \in \mathcal{B}_{0}(X)$. Let an integral kernel $K(x, y)$ of the operator $K$ be chosen so that statements (i)-(iii) of Proposition 13 are satisfied. Then there exists a unique point process $\mu$ on $X$ which has correlation functions (39) if and only if $0 \leq \widehat{K} \leq 1$.

Proof. We only have to prove that, if a point process $\mu$ exists, then $0 \leq \widehat{K} \leq$ 1. We divide the proof into several steps.
(1) Fix any compact $\Delta \subset X$. By Proposition 14 and (38), the Bogoliubov functional of $\mu$ is given by (40). Hence, analogously to the proof of Corollary 2, we get

$$
\mu_{\Delta}(\{\varnothing\})=\operatorname{Det}\left(1-K^{\Delta}\right)
$$

In particular,

$$
\begin{equation*}
\operatorname{Det}\left(1-K^{\Delta}\right) \geq 0 \tag{50}
\end{equation*}
$$

(2) From now on we will additionally assume that $\|K\|<1$. Then $\left\|K^{\Delta}\right\|<1$ and we set $L[\Delta]:=K^{\Delta}\left(1-K^{\Delta}\right)^{-1}$. Just as in the proof of Proposition 10, we derive that $L[\Delta]$ is $J$-self-adjoint and $L[\Delta] \in \mathscr{L}_{1 \mid 2}\left(L^{2}(X, m)\right)$. To choose an integral kernel of the operator $L[\Delta]$, we represent it in the form

$$
L[\Delta]=K^{\Delta}+K^{\Delta} L[\Delta] .
$$

As $L[\Delta] \in \mathscr{L}_{2}\left(L^{2}(\Delta, m)\right)$, we first choose an arbitrary $J$-Hermitian integral kernel of this operator, which we denote by $\tilde{L}[\Delta](x, y)$. Now we set

$$
L[\Delta](x, y):=K(x, y)+\int_{\Delta} K(x, z) \widetilde{L}[\Delta](z, y) m(d y), \quad x, y \in \Delta
$$

As is easily seen, this integral kernel satisfies

$$
\begin{align*}
& \operatorname{Tr}\left(L[\Delta]_{\text {even }}^{\Lambda}\right) \\
& \quad=\int_{\Lambda} L[\Delta](x, x) m(d x) \quad \text { for each } \Lambda \in \mathcal{B}_{0}(X), \Lambda \subset \Delta \tag{51}
\end{align*}
$$

(3) By (51),
(52) $\operatorname{Tr}\left((L[\Delta](1+\varphi))_{\text {even }}\right)=\int_{\Delta} L[\Delta](x, x)(1+\varphi(x)), \quad \varphi \in B_{0}(X)$.

Since formula (45) clearly holds for the Boliubov functional of $\mu_{\Delta}$, using (52) and Proposition 6, we get, for each $\varphi \in B(\Delta)$,

$$
\begin{aligned}
& B_{\mu_{\Delta}}(\varphi)=\operatorname{Det}\left(1-K^{\Delta}\right) \\
& \times\left(1+\sum_{n=1}^{\infty} \int_{\Delta^{n}}\left(1+\varphi\left(x_{1}\right)\right) \cdots\left(1+\varphi\left(x_{n}\right)\right)\right. \\
& \left.\times \operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)\right) .
\end{aligned}
$$

This implies that the measure $\mu_{\Delta}$ has densities (43). Hence, by (50), for $m^{\otimes n}$-a.a. $\left(x_{1}, \ldots, x_{n}\right) \in \Delta^{n}$,

$$
\operatorname{det}\left[L[\Delta]\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} \geq 0
$$

In particular, for $i=1,2$, for $m^{\otimes n}-$ a.a. $\left(x_{1}, \ldots, x_{n}\right) \in \Delta_{i}^{n}$,

$$
\begin{equation*}
\operatorname{det}\left[L[\Delta]_{i i}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} \geq 0 \tag{53}
\end{equation*}
$$

Here $\Delta_{i}:=\Delta \cap X_{i}, i=1,2$.
(4) Following [17], Proposition 1.5, let us find a representation of $L[\Delta]_{11}$ in terms of the blocks of the operator $K^{\Delta}$. Since $L[\Delta]\left(1-K^{\Delta}\right)=K^{\Delta}$, we have

$$
\begin{align*}
L[\Delta]_{11}\left(1-K_{11}^{\Delta}\right)-L[\Delta]_{12} K_{21}^{\Delta} & =K_{11}^{\Delta},  \tag{54}\\
-L[\Delta]_{11} K_{12}^{\Delta}+L[\Delta]_{12}\left(1-K_{22}^{\Delta}\right) & =K_{12}^{\Delta} . \tag{55}
\end{align*}
$$

From (55),

$$
-L[\Delta]_{11} K_{12}^{\Delta}\left(1-K_{22}^{\Delta}\right)^{-1}+L[\Delta]_{12}=K_{12}^{\Delta}\left(1-K_{22}^{\Delta}\right)^{-1}
$$

hence,

$$
-L[\Delta]_{11} K_{12}^{\Delta}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}+L[\Delta]_{12} K_{21}^{\Delta}=K_{12}^{\Delta}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}
$$

Adding this to (54) yields

$$
L[\Delta]_{11}\left(1-Q[\Delta]_{11}\right)=Q[\Delta]_{11}
$$

where

$$
\begin{align*}
Q[\Delta]_{11} & :=K_{11}^{\Delta}+K_{12}^{\Delta}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta} \\
& =K_{11}^{\Delta}-\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta} . \tag{56}
\end{align*}
$$

Since the operator $1-K_{11}^{\Delta}$ is strictly positive and the operator $\left(K_{21}^{\Delta}\right)^{*}(1-$ $\left.K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}$ is nonnegative, the operator $1-Q[\Delta]_{11}$ is strictly positive, hence invertible. Therefore,

$$
\begin{equation*}
L[\Delta]_{11}=Q[\Delta]_{11}\left(1-Q[\Delta]_{11}\right)^{-1} \tag{57}
\end{equation*}
$$

(5) By (56), the operator $Q[\Delta]_{11}$ is self-adjoint and trace-class. Since the operator $1-Q[\Delta]_{11}$ is strictly positive, we therefore get

$$
\operatorname{Det}\left(1-Q[\Delta]_{11}\right)>0
$$

[Note that $\operatorname{Det}\left(1-Q[\Delta]_{11}\right)$ is a usual Fredholm determinant.] Therefore, by (53), we can define a nonnegative, finite measure $v\left[\Delta_{1}\right]$ on $\left(\Gamma_{\Delta_{1}}, \mathcal{B}\left(\Gamma_{\Delta_{1}}\right)\right)$ whose local densities are

$$
\begin{align*}
& d_{\nu\left[\Delta_{1}\right]}^{(0)}=\operatorname{Det}\left(1-Q[\Delta]_{11}\right) \\
& d_{\nu\left[\Delta_{1}\right]}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\operatorname{Det}\left(1-Q[\Delta]_{11}\right) \operatorname{det}\left[L[\Delta]_{11}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n},  \tag{58}\\
& n \in \mathbb{N} .
\end{align*}
$$

Analogously to (43)-(46), we conclude from (58) that the Bogoliubov transform of the measure $\nu\left[\Delta_{1}\right]$ is given by

$$
\begin{align*}
B_{\nu\left[\Delta_{1}\right]}(\varphi) & =\int_{\Gamma_{\Delta_{1}}} \prod_{x \in \gamma}(1+\varphi(x)) \nu\left[\Delta_{1}\right](d \gamma) \\
& =\operatorname{Det}\left(1+\operatorname{sgn}(\varphi) \sqrt{|\varphi|} Q[\Delta]_{11} \sqrt{|\varphi|}\right), \quad \varphi \in B_{0}\left(\Delta_{1}\right) \tag{59}
\end{align*}
$$

Setting $\varphi \equiv 0$, we see that $\nu\left[\Delta_{1}\right]\left(\Gamma_{\Delta_{1}}\right)=1$, that is, $\nu\left[\Delta_{1}\right]$ is a point process in $\Delta_{1}$.
(6) We can now choose an integral kernel of the operator $Q[\Delta]_{11}$ analogously to [10], Lemma A.3, and [14], Section 3. Indeed, since $K_{21}^{\Delta}$ is a HilbertSchmidt operator, $\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}$ is a nonnegative trace-class operator in $L^{2}\left(\Delta_{1}, m\right)$. The operator $\left(\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}\right)^{1 / 2}$ is Hilbert-Schmidt, hence an integral operator. We choose its integral kernel, denoted by $\theta(x, y)$, so that

$$
\begin{aligned}
\theta(x, y) & =\overline{\theta(y, x)} \quad \text { for all } x, y \in \Delta_{1} \\
\theta(x, \cdot) & \in L^{2}\left(\Delta_{1}, m\right) \quad \text { for all } x \in \Delta_{1}
\end{aligned}
$$

[Recall that $\int_{\Delta_{1}^{2}}|\theta(x, y)|^{2} m(d x) m(d y)=\left\|\left(\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}\right)^{1 / 2}\right\|_{2}^{2}<\infty$.] Now, we set an integral kernel of the operator $\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}$ to be

$$
\begin{aligned}
\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta}(x, y): & =\int_{\Delta_{1}} \theta(x, z) \theta(z, y) m(d z) \\
& =(\theta(x, \cdot), \theta(y, \cdot))_{L^{2}\left(\Delta_{1}, m\right)}, \quad x, y \in \Delta_{1}
\end{aligned}
$$

We similarly construct an integral kernel of the operator $K_{11}^{\Delta}$ :

$$
K_{11}^{\Delta}(x, y)=(\eta(x, \cdot), \eta(y, \cdot))_{L^{2}\left(\Delta_{1}, m\right)}, \quad x, y \in \Delta_{1}
$$

Hence, by virtue of (56), we may choose an integral kernel of the operator $Q[\Delta]_{11}$ as follows:

$$
\begin{equation*}
Q[\Delta]_{11}(x, y)=(\eta(x, \cdot), \eta(y, \cdot))_{L^{2}\left(\Delta_{1}, m\right)}-(\theta(x, \cdot), \theta(y, \cdot))_{L^{2}\left(\Delta_{1}, m\right)} \tag{60}
\end{equation*}
$$

As is easily seen, for each $\Delta \in \mathcal{B}_{0}(X), \Lambda \subset \Delta_{1}$,

$$
\operatorname{Tr}\left(Q[\Delta]_{11}^{\Lambda}\right)=\int_{\Lambda} Q[\Delta]_{11}(x, x) m(d x)
$$

Now, analogously to Proposition 14, we get from (59)

$$
\begin{aligned}
B_{v\left[\Delta_{1}\right]}(\varphi)=1+\sum_{n=1}^{\infty} \frac{1}{n!} \int_{\Delta_{1}^{n}} & \varphi\left(x_{1}\right) \cdots \varphi\left(x_{n}\right) \\
& \times \operatorname{det}\left[Q[\Delta]_{11}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} m\left(d x_{1}\right) \cdots m\left(d x_{n}\right)
\end{aligned}
$$

for each $\varphi \in B_{0}\left(\Delta_{1}\right)$. Hence, the correlation functions of the point process $\nu\left[\Delta_{1}\right]$ are

$$
k_{\nu\left[\Delta_{1}\right]}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\operatorname{det}\left[Q[\Delta]_{11}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n}, \quad n \in \mathbb{N} .
$$

Therefore, for each $n \in \mathbb{N}$,

$$
\begin{equation*}
\operatorname{det}\left[Q[\Delta]_{11}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} \geq 0 \quad \text { for } m^{\otimes n} \text {-a.a. }\left(x_{1}, \ldots, x_{n}\right) \in \Delta_{1}^{n} \tag{61}
\end{equation*}
$$

(7) Obviously, the following two mappings are measurable:

$$
\Delta_{1} \ni x \mapsto \eta(x, \cdot) \in L^{2}\left(\Delta_{1}, m\right), \quad \Delta_{1} \ni x \mapsto \theta(x, \cdot) \in L^{2}\left(\Delta_{1}, m\right)
$$

Therefore, by Lusin's theorem (see, e.g., [18]), for each $\varepsilon>0$, there exists a compact set $\Lambda_{\varepsilon} \subset \Delta_{1}$ such that $m\left(\Delta_{1} \backslash \Lambda_{\varepsilon}\right) \leq \varepsilon$ and the mappings

$$
\Lambda_{\varepsilon} \ni x \mapsto \eta(x, \cdot) \in L^{2}\left(\Delta_{1}, m\right), \quad \Lambda_{\varepsilon} \ni x \mapsto \theta(x, \cdot) \in L^{2}\left(\Delta_{1}, m\right)
$$

are continuous. Therefore, by (60), the function

$$
\begin{equation*}
\Lambda_{\varepsilon}^{2} \ni(x, y) \mapsto Q[\Delta]_{11}(x, y) \in \mathbb{C} \tag{62}
\end{equation*}
$$

is continuous. Hence, by (61),

$$
\operatorname{det}\left[Q[\Delta]_{11}\left(x_{i}, x_{j}\right)\right]_{i, j=1, \ldots, n} \geq 0 \quad \text { for all }\left(x_{1}, \ldots, x_{n}\right) \in \Lambda_{\varepsilon}^{n}
$$

Thus, the continuous kernel (62) is positive definite, and therefore the operator $Q[\Delta]_{11}$ is nonnegative on $L^{2}\left(\Lambda_{\varepsilon}, m\right)$. By letting $\varepsilon \rightarrow 0$, we conclude that $Q[\Delta]_{11} \geq 0$ on $L^{2}\left(\Delta_{1}, m\right)$. Hence, by (56),

$$
\begin{equation*}
K_{11}^{\Delta} \geq\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta} \quad \text { on } L^{2}\left(\Delta_{1}, m\right) \tag{63}
\end{equation*}
$$

(8) We denote by $\widehat{K}^{\Delta}$ the corresponding transformation of the operator $K^{\Delta}$ in the Hilbert space $L^{2}(\Delta, m)=L^{2}\left(\Delta_{1}, m\right) \oplus L^{2}\left(\Delta_{2}, m\right)$. Hence, $\widehat{K}^{\Delta}=P^{\Delta} \widehat{K} P^{\Delta}$ and

$$
\widehat{K}^{\Delta}=\left[\begin{array}{cc}
K_{11}^{\Delta} & K_{21}^{\Delta} \\
\left(K_{21}^{\Delta}\right)^{*} & 1-K_{22}^{\Delta}
\end{array}\right] .
$$

By (63), for each $f=\left(f_{1}, f_{2}\right) \in L^{2}(\Delta, m)$,

$$
\begin{aligned}
\left(\widehat{K}^{\Delta} f, f\right)= & \left(K_{11}^{\Delta} f_{1}, f_{1}\right)+\left(K_{21}^{\Delta} f_{1}, f_{2}\right) \\
& +\left(\left(K_{21}^{\Delta}\right)^{*} f_{2}, f_{1}\right)+\left(\left(1-K_{22}^{\Delta}\right) f_{2}, f_{2}\right) \\
\geq & \left(\left(K_{21}^{\Delta}\right)^{*}\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta} f_{1}, f_{1}\right) \\
& +\left(\left(1-K_{22}^{\Delta}\right) f_{2}, f_{2}\right)-2\left|\left(K_{21}^{\Delta} f_{1}, f_{2}\right)\right| \\
= & \left(\left(1-K_{22}^{\Delta}\right)^{-1} K_{21}^{\Delta} f_{1}, K_{21}^{\Delta} f_{1}\right) \\
& +\left(\left(1-K_{22}^{\Delta}\right) f_{2}, f_{2}\right)-2\left|\left(K_{21}^{\Delta} f_{1}, f_{2}\right)\right|
\end{aligned}
$$

Since $K_{22}^{\Delta}$ is a compact self-adjoint operator in $L^{2}\left(\Delta_{2}, m\right)$, we can choose an orthnormal basis of $L^{2}\left(\Delta_{2}, m\right)$ which consists of eigenvectors of the operator $K_{22}^{\Delta}$,
and we denote by $\lambda_{n}$ the eigenvalue belonging to eigenvector $e_{n}, n \geq 1$. Clearly, $\lambda_{n}<1$ for all $n$. Then, by (64),

$$
\begin{aligned}
\left(\widehat{K}^{\Delta} f, f\right) \geq & \sum_{n=1}^{\infty}\left(1-\lambda_{n}\right)^{-1}\left|\left(K_{21}^{\Delta} f_{1}, e_{n}\right)\right|^{2}+\sum_{n=1}^{\infty}\left(1-\lambda_{n}\right)\left|\left(f_{2}, e_{n}\right)\right|^{2} \\
& -\sum_{n=1}^{\infty} 2\left|\left(K_{21}^{\Delta} f_{1}, e_{n}\right)\left(f_{2}, e_{n}\right)\right| \\
= & \sum_{n=1}^{\infty}\left(\left(1-\lambda_{n}\right)^{-1 / 2}\left|\left(K_{21}^{\Delta} f_{1}, e_{n}\right)\right|-\left(1-\lambda_{n}\right)^{1 / 2}\left|\left(f_{2}, e_{n}\right)\right|\right)^{2} \geq 0
\end{aligned}
$$

Thus, for each compact $\Delta \subset X$, the operator $\widehat{K}^{\Delta}=P_{\Delta} \widehat{K} P_{\Delta}$ is nonnegative. Hence, $\widehat{K} \geq 0$. Exchanging the role of the sets $X_{1}$ and $X_{2}$ and using instead of the operator $K$ the operator $1-K$, we therefore get $1-\widehat{K} \geq 0$. Thus, $0 \leq \widehat{K} \leq 1$.
(9) We now assume that $\|K\|=1$. Using the procedure of thinning of the point process $\mu$ (see, e.g., [9], Example 8.2(a)), we conclude that, for each $\varepsilon \in(0,1)$, there exists a point process $\mu_{\varepsilon}$ which has correlation functions as in formula (48), that is, a determinantal point process corresponding to the operator $K_{\varepsilon}:=\varepsilon K$. By the proved above $0 \leq \widehat{K}_{\varepsilon} \leq 1$. Hence, by (47),

$$
0 \leq \varepsilon \widehat{K}+(1-\varepsilon) P_{2} \leq 1
$$

Letting $\varepsilon \rightarrow 1$, we get $0 \leq \widehat{K} \leq 1$.
(10) Finally, we assume that $\|K\|>1$ and we have to show that a determinantal point process does not exist in this case. Assume the contrary, that is, assume that there exists a determinantal point process with correlation kernel $K(x, y)$. Since $\|K\|>1$, there exists a compact set $\Delta \subset X$ such that $\left\|K^{\Delta}\right\|>1$. Analogously to part (9), using the procedure of thinning, we conclude that, for each $\varepsilon \in(0,1)$, there exists a determinantal point process with correlation kernel $K_{\varepsilon}(x, y):=\varepsilon K(x, y)$. We choose $\varepsilon:=\left\|K^{\Delta}\right\|^{-1}$, so that $\left\|K_{\varepsilon}^{\Delta}\right\|=1$. We take the restriction of the corresponding probability measure to the $\sigma$-algebra $\mathcal{B}_{\Delta}(\Gamma)$, that is, a point process on $\left(\Gamma_{\Delta}, \mathcal{B}\left(\Gamma_{\Delta}\right)\right)$. We denote this point process by $\mu_{\varepsilon, \Delta}$. By part (9), we have $0 \leq \widehat{K}_{\varepsilon}^{\Delta} \leq 1$. Then, by Corollary 2,

$$
\operatorname{Det}\left(1-K_{\varepsilon}^{\Delta}\right)=0
$$

Next, following the idea of [21], Remark 4, we consider

$$
\begin{aligned}
\int_{\Gamma}(1-\varepsilon)^{|\gamma \cap \Delta|} \mu(d \gamma) & =\int_{\Gamma} \prod_{x \in \gamma}\left(1-\varepsilon \chi_{\Delta}(x)\right) \mu(d \gamma) \\
& =\operatorname{Det}\left(1-\varepsilon K^{\Delta}\right)=\operatorname{Det}\left(1-K_{\varepsilon}^{\Delta}\right)=0 .
\end{aligned}
$$

On the other hand, $(1-\varepsilon)^{|\gamma \cap \Delta|}>0$ for all $\gamma \in \Gamma$. Hence,

$$
\int_{\Gamma}(1-\varepsilon)^{|\gamma \cap \Delta|} \mu(d \gamma)>0
$$

which is a contradiction.

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