# A NOTE ON EQUILIBRIUM GLAUBER AND KAWASAKI DYNAMICS FOR PERMANENTAL POINT PROCESSES 

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#### Abstract

We construct two types of equilibrium dynamics of an infinite particle system in a locally compact metric space $X$ for which a permanental point process is a symmetrizing, and hence invariant measure. The Glauber dynamics is a birth-anddeath process in $X$, while in the Kawasaki dynamics interacting particles randomly hop over $X$. In the case $X=\mathbb{R}^{d}$, we consider a diffusion approximation for the Kawasaki dynamics at the level of Dirichlet forms. This leads us to an equilibrium dynamics of interacting Brownian particles for which a permanental point process is a symmetrizing measure.


## 1. Introduction

Let $X$ be a locally compact Polish space and let $\nu$ be a Radon non-atomic measure on it. Let $\Gamma=\Gamma_{X}$ denote the space of all locally finite subsets (configurations) in $X$.

A Glauber dynamics (a birth-and-death process of an infinite system of particles in $X$ ) is a Markov process on $\Gamma$ whose formal (pre-)generator has the form

$$
\begin{align*}
\left(L_{\mathrm{G}} F\right)(\gamma) & =\sum_{x \in \gamma} d(x, \gamma \backslash x)(F(\gamma \backslash x)-F(\gamma))  \tag{1.1}\\
& +\int_{X} \nu(d x) b(x, \gamma)(F(\gamma \cup x)-F(\gamma)), \quad \gamma \in \Gamma
\end{align*}
$$

Here and below, for simplicity of notation we write $x$ instead of $\{x\}$. The coefficient $d(x, \gamma \backslash x)$ describes the rate at which particle $x$ of configuration $\gamma$ dies, while $b(x, \gamma)$ describes the rate at which, given configuration $\gamma$, a new particle is born at $x$.

A Kawasaki dynamics (a dynamics of hopping particles) is a Markov process on $\Gamma$ whose formal (pre-)generator is

$$
\begin{equation*}
\left(L_{\mathrm{K}} F\right)(\gamma)=\sum_{x \in \gamma} c(x, y, \gamma \backslash x)(F(\gamma \backslash x \cup y)-F(\gamma)), \quad \gamma \in \Gamma \tag{1.2}
\end{equation*}
$$

The coefficient $c(x, y, \gamma \backslash x)$ describes the rate at which particle $x$ of configuration $\gamma$ hops to $y$, taking the rest of the configuration, $\gamma \backslash x$, into account.

Equilibrium Glauber and Kawasaki dynamics which have a standard Gibbs measure as symmetrizing (and hence invariant) measure were constructed in [19, 20]. In [22], this construction was extended to the case of an equilibrium dynamics which has a determinantal (fermion) point process as invariant measure, For further studies of equilibrium and non-equilibrium Glauber and Kawasaki dynamics, we refer to $[3,7,8,9,10,11,13$, $14,15,16,17,18,21,28]$ and the references therein.

The aim of this note is to show that general criteria of existence of Glauber and Kawasaki dynamics which were developed in [22] are applicable to a wide class of $\alpha$-permanental $(\alpha \in \mathbb{N})$ point processes, proposed by Shirai and Takahashi [30]. This

[^0]class includes classical permanental (boson) point processes, see e.g. [5, 30]. We will also consider a diffusion approximation for the Kawasaki dynamics at the level of Dirichlet forms (compare with [15]). This will lead us to an equilibrium dynamics of interacting Brownian particles for which an $\alpha$-permanental point process is a symmetrizing measure. As a by-product of our considerations, we will also extend the result of [30] on the existence of $\alpha$-permanental point process.

## 2. Equilibrium Glauber and Kawasaki dynamics - General results

Let $X$ be a locally compact Polish space. We denote by $\mathcal{B}(X)$ the Borel $\sigma$-algebra on $X$, and by $\mathcal{B}_{0}(X)$ the collection of all sets from $\mathcal{B}(X)$ which are relatively compact. We fix a Radon, non-atomic measure on $(X, \mathcal{B}(X)$ ). (For most applications, the reader may think of $X$ as $\mathbb{R}^{d}$ and $\nu$ as the Lebegue measure.)

The configuration space $\Gamma$ over $X$ is defined as the set of all subsets of $X$ which are locally finite

$$
\Gamma:=\left\{\gamma \subset X:\left|\gamma_{\Lambda}\right|<\infty \text { for each } \Lambda \in \mathcal{B}_{0}(X)\right\}
$$

where $|\cdot|$ denotes the cardinality of a set and $\gamma_{\Lambda}:=\gamma \cap \Lambda$. One can identify any $\gamma \in \Gamma$ with the positive Radon measure $\sum_{x \in \gamma} \varepsilon_{x}$, where $\varepsilon_{x}$ is the Dirac measure with mass at $x$ and $\sum_{x \in \varnothing} \varepsilon_{x}:=$ zero measure. The space $\Gamma$ can be endowed with the vague topology, i.e., the weakest topology on $\Gamma$ with respect to which all maps

$$
\Gamma \ni \gamma \mapsto\langle\varphi, \gamma\rangle:=\int_{X} \varphi(x) \gamma(d x)=\sum_{x \in \gamma} \varphi(x), \quad \varphi \in C_{0}(X)
$$

are continuous. Here, $C_{0}(X)$ is the space of all continuous, real-valued functions on $X$ with compact support. We denote the Borel $\sigma$-algebra on $\Gamma$ by $\mathcal{B}(\Gamma)$. A point process in $X$ is a probability measure on $(\Gamma, \mathcal{B}(\Gamma))$.

We fix a point process $\mu$ which satisfies the so-called condition $\left(\Sigma_{\nu}^{\prime}\right)$ [5, 26], i.e., there exist a measurable function $r: X \times \Gamma \rightarrow[0,+\infty]$, called the Papangelou intensity of $\mu$, such that

$$
\begin{equation*}
\int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) F(x, \gamma)=\int_{\Gamma} \mu(d \gamma) \int_{X} \nu(d x) r(x, \gamma) F(x, \gamma \cup x) \tag{2.1}
\end{equation*}
$$

for any measurable function $F: X \times \Gamma \rightarrow[0,+\infty]$. The condition $\left(\Sigma_{\nu}^{\prime}\right)$ can be thought of as a kind of weak Gibbsianess of $\mu$. Intuitively, we may treat the Papangelou intensity as

$$
\begin{equation*}
r(x, \gamma)=\exp [-E(x, \gamma)] \tag{2.2}
\end{equation*}
$$

where $E(x, \gamma)$ is the relative energy of interaction between particle $x$ and configuration $\gamma$.
To define an equilibrium Glauber dynamics for which $\mu$ is a symmetrizing measure, we fix a death coefficient as a measurable function $d: X \times \Gamma \rightarrow[0,+\infty]$, and then define a birth coefficient $b: X \times \Gamma \rightarrow[0,+\infty]$ by

$$
\begin{equation*}
b(x, \gamma)=d(x, \gamma) r(x, \gamma), \quad(x, \gamma) \in X \times \Gamma \tag{2.3}
\end{equation*}
$$

To define a Kawasaki dynamics, we fix a measurable function $c: X^{2} \times \Gamma^{2} \rightarrow[0,+\infty]$ which satisfies

$$
\begin{equation*}
r(x, \gamma) c(x, y, \gamma)=r(y, \gamma) c(y, x, \gamma), \quad(x, y, \gamma) \in X^{2} \times \Gamma \tag{2.4}
\end{equation*}
$$

Formulas (2.3) and (2.4) are called the balance conditions [13, 14]. We will also assume that the function $c(x, y, \gamma)$ vanishes if at least one of the functions $r(x, \gamma)$ and $r(y, \gamma)$ vanishes, i.e.,

$$
\begin{equation*}
c(x, y, \gamma)=c(x, y, \gamma) \chi_{\{r>0\}}(x, \gamma) \chi_{\{r>0\}}(y, \gamma) \tag{2.5}
\end{equation*}
$$

Here, for a set $A, \chi_{A}$ denotes the indicator function of $A$. We refer to [22, Remark 3.1] for a justification of this assumption, which involves the interpretation of $r(x, \gamma)$ as in (2.2), see also Remark 2.4 below.

We denote by $\mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right)$ the space of all functions of the form

$$
\begin{equation*}
\Gamma \ni \gamma \mapsto F(\gamma)=g\left(\left\langle\varphi_{1}, \gamma\right\rangle, \ldots,\left\langle\varphi_{N}, \gamma\right\rangle\right) \tag{2.6}
\end{equation*}
$$

where $N \in \mathbb{N}, \varphi_{1}, \ldots, \varphi_{N} \in C_{0}(X)$ and $g \in C_{\mathrm{b}}\left(\mathbb{R}^{N}\right)$. Here, $C_{\mathrm{b}}\left(\mathbb{R}^{N}\right)$ denotes the space of all continuous bounded functions on $\mathbb{R}^{N}$. We assume that, for each $\Lambda \in \mathcal{B}_{0}(X)$,

$$
\begin{gather*}
\int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \gamma(d x) d(x, \gamma \backslash x)<\infty  \tag{2.7}\\
\int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) \int_{X} \nu(d y) c(x, y, \gamma \backslash x)\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(y)\right)<\infty \tag{2.8}
\end{gather*}
$$

As easily seen, conditions (2.7) and (2.8) are sufficient in order to define bilinear forms

$$
\begin{aligned}
\mathcal{E}_{\mathrm{G}}(F, G) & :=\int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) d(x, \gamma \backslash x)(F(\gamma \backslash x)-F(\gamma))(G(\gamma \backslash x)-G(\gamma)) \\
\mathcal{E}_{\mathrm{K}}(F, G): & =\frac{1}{2} \int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) \int_{X} \nu(d y) c(x, y, \gamma \backslash x)(F(\gamma \backslash x \cup y)-F(\gamma)) \\
& \times(G(\gamma \backslash x \cup y)-G(\gamma))
\end{aligned}
$$

where $F, G \in \mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right)$.
For the construction of the Kawasaki dynamics, we will also assume that the following technical assumptions holds:

$$
\begin{align*}
& \exists u, v \in \mathbb{R} \quad \forall \Lambda \in \mathcal{B}_{0}(X): \\
& \quad \int_{\Lambda} \gamma(d x) \int_{\Lambda} \nu(d y) r(x, \gamma \backslash x)^{u} r(y, \gamma \backslash x)^{v} c(x, y, \gamma \backslash y) \in L^{2}(\Gamma, \mu)<\infty . \tag{2.9}
\end{align*}
$$

Note that in formula (2.9) and below, we use the convention $\frac{0}{0}:=0$.
The following theorem was essentially proved in [22].
Theorem 2.1. (i) Assume that a point process $\mu$ satisfies (2.1). Assume that conditions (2.3), (2.7), respectively (2.4), (2.5), (2.8), and (2.9) are satisfied. Let $\sharp=\mathrm{G}, \mathrm{K}$. Then the bilinear form $\left(\mathcal{E}_{\sharp}, \mathcal{F} C_{\mathrm{b}}\left(C_{0}(x), \Gamma\right)\right)$ is closable in $L^{2}(\Gamma, \mu)$ and its closure will be denoted by $\left(\mathcal{E}_{\sharp}, D\left(\mathcal{E}_{\sharp}\right)\right)$. Further there exists a conservative Hunt process (Glauber, respectively Kawasaki dynamics)

$$
M^{\sharp}=\left(\Omega^{\sharp}, \mathcal{F}^{\sharp},\left(\mathcal{F}_{t}^{\sharp}\right)_{t \geq 0},\left(\Theta_{t}^{\sharp}\right)_{t \geq 0},\left(X^{\sharp}(t)\right)_{t \geq 0},\left(P_{\gamma}^{\sharp}\right)_{\gamma \in \Gamma}\right)
$$

on $\Gamma$ which is properly associated with $\left(\mathcal{E}_{\sharp}, D\left(\mathcal{E}_{\sharp}\right)\right)$, i.e., for all ( $\mu$-version of) $F \in L^{2}(\Gamma, \mu)$ and $t>0$

$$
\Gamma \ni \gamma \mapsto p_{t}^{\sharp} F(\gamma):=\int_{\Omega^{\sharp}} F\left(X^{\sharp}(t)\right) d P_{\gamma}^{\sharp}
$$

is an $\mathcal{E}^{\sharp}$-quasi continuous version of $\exp \left(t L_{\sharp}\right) F$, where $\left(-L_{\sharp}, D\left(L_{\sharp}\right)\right)$ is the generator of $\left(\mathcal{E}_{\sharp}, D\left(\mathcal{E}_{\sharp}\right)\right) . M^{\sharp}$ is up-to $\mu$-equivalence unique. In particular, $M^{\sharp}$ is $\mu$-symmetric and has $\mu$ as invariant measure.
(ii) $M^{\sharp}$ from (i) is up to $\mu$-equivalence unique between all Hunt processes

$$
M^{\prime}=\left(\Omega^{\prime}, \mathcal{F}^{\prime},\left(\mathcal{F}_{t}^{\prime}\right)_{t \geq 0},\left(\Theta_{t}^{\prime}\right)_{t \geq 0},\left(X^{\prime}(t)\right)_{t \geq 0},\left(P_{\gamma}^{\prime}\right)_{\gamma \in \Gamma}\right)
$$

on $\Gamma$ having $\mu$ as invariant measure and solving a martingale problem for $\left(L_{\sharp}, D\left(L_{\sharp}\right)\right)$, i.e., for all $G \in D\left(H_{\sharp}\right)$

$$
\widetilde{G}\left(X^{\prime}(t)\right)-\widetilde{G}\left(X^{\prime}(0)\right)-\int_{0}^{t}\left(L_{\sharp} G\right)\left(X^{\prime}(s)\right) d s, \quad t \geq 0
$$

is an $\left(\mathcal{F}_{t}^{\prime}\right)$-martingale under $P_{\gamma}^{\prime}$ for $\mathcal{E}_{\sharp}$-q.e. $\gamma \in \Gamma$. Here, $\widetilde{G}$ denotes an $\mathcal{E}_{\sharp}$-quasi-continuous version of $G$.
(iii) Further assume that, for each $\Lambda \in \mathcal{B}_{0}(X)$,

$$
\begin{equation*}
\int_{\Lambda} \gamma(d x) d(x, \gamma \backslash x) \in L^{2}(\Gamma, \mu), \quad \int_{\Lambda} \nu(d x) b(x, \gamma) \in L^{2}(\Gamma, \mu) \tag{2.10}
\end{equation*}
$$

in the Glauber case, and

$$
\begin{equation*}
\int_{X} \gamma(d x) \int_{X} \nu(d y) c(x, y, \gamma \backslash x)\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(y)\right) \in L^{2}(\Gamma, \mu) \tag{2.11}
\end{equation*}
$$

in the Kawasaki case. Then $\mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right) \subset D\left(L_{\sharp}\right)$, and for each $F \in \mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right)$, $L_{\sharp} F$ is given by formulas (1.1) and (1.2), respectively.

Remark 2.1. We refer to [24] for an explanation of notions appearing in Theorem 2.1, see also a brief explanation of them in [22].

Proof of Theorem 2.1. The statement follows from Theorems 3.1 and 3.2 in [22]. Note that, although these theorems are formulated for determinantal point processes only, their proof only uses the $\left(\Sigma_{\nu}^{\prime}\right)$ property of these point processes. Note also that condition (2.9) is formulated in [22] only for $v=1$, however the proof of Lemma 3.2 in [22] admits a straightforward generalization to the case of an arbitrary $v \in \mathbb{R}$.
Remark 2.2. Part (iii) of Theorem 2.1 states that the operator $\left(-L_{\sharp}, D\left(L_{\sharp}\right)\right)$ is the Friedrichs' extension of the operator $\left(-L_{\sharp}, \mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right)\right)$ defined by formulas (1.1), (1.2), respectively.

Let us fix a parameter $s \in[0,1]$ and define

$$
\begin{gather*}
d(x, \gamma):=r(x, \gamma)^{s-1} \chi_{\{r>0\}}(x, \gamma), \quad(x, \gamma) \in X \times \Gamma,  \tag{2.12}\\
b(x, \gamma):=r(x, \gamma)^{s} \chi_{\{r>0\}}(x, \gamma), \quad(x, \gamma) \in X \times \Gamma,  \tag{2.13}\\
c(x, y, \gamma):=a(x, y) r(x, \gamma)^{s-1} r(y, \gamma)^{s} \chi_{\{r>0\}}(x, \gamma) \chi_{\{r>0\}}(y, \gamma),  \tag{2.14}\\
(x, y, \gamma) \in X^{2} \times \Gamma .
\end{gather*}
$$

Here the function $a: X^{2} \rightarrow[0,+\infty)$ is bounded, measurable, symmetric (i.e., $a(x, y)=$ $a(y, x))$, and satisfies

$$
\begin{equation*}
\sup _{x \in X} \int_{X} a(x, y) \nu(d y)<\infty \tag{2.15}
\end{equation*}
$$

Note that the balance conditions (2.3) and (2.4) are satisfied for these coefficients, and so is condition (2.5).
Remark 2.3. Note that, if $X=\mathbb{R}^{d}$ and $a(x, y)$ has the form $a(x-y)$ for a function $a: \mathbb{R}^{d} \rightarrow[0, \infty)$, then condition (2.15) means that $a \in L^{1}\left(\mathbb{R}^{d}, d x\right)$. (Here and below, in the case $X=\mathbb{R}^{d}$, we use an obvious abuse of notation.)
Remark 2.4. Using representation (2.2), we can rewrite formulas (2.12)-(2.14) as follows:

$$
\begin{aligned}
d(x, \gamma \backslash x) & =\exp [(1-s) E(x, \gamma \backslash x)] \chi_{\{E<+\infty\}}(x, \gamma \backslash x), \\
b(x, \gamma \backslash x) & =\exp [-s E(x, \gamma \backslash x)] \chi_{\{E<+\infty\}}(x, \gamma \backslash x), \\
c(x, y, \gamma \backslash x) & =a(x, y) \exp [(1-s) E(x, \gamma \backslash x)-s E(y, \gamma \backslash x)] \\
& \times \chi_{\{E<+\infty\}}(x, \gamma \backslash x) \chi_{\{E<+\infty\}}(y, \gamma \backslash x) .
\end{aligned}
$$

So, if the corresponding dynamics exist, one can give the following heuristic description of them: Both dynamics are concentrated on configurations $\gamma \in \Gamma$ such that, for each $x \in \gamma$, the relative energy of interaction between $x$ and the rest of configuration, $\gamma \backslash x$, is finite; those particles tend to die, respectively hop, which have a high energy of interaction
with the rest of the configuration, while it is more probable that a new particle is born at $y$, respectively $x$ hops to $y$, if the energy of interaction between $y$ and the rest of the configuration is low.

Let us assume that the point process $\mu$ satisfies:

$$
\forall \Lambda \in \mathcal{B}_{0}(X): \quad \int_{\Lambda} \gamma(d x) \in L^{2}(\Gamma, \mu)
$$

Then, by choosing $u=1-s$ and $v=-s$ in (2.9), we conclude that the coefficient $c$ given by (2.14) satisfies (2.9).

We will construct below a class of point processes $\mu$ for which the coefficients $d, b$ and $c$ given above satisfy the other conditions of Theorem 2.1.

## 3. Permanental point processes and corresponding equilibrium dynamics

Let $K$ be a linear, bounded, self-adjoint operator on the real space $L^{2}(X, \nu)$. Further assume that $K \geq 0$ and $K$ is locally of trace class, i.e., $\operatorname{Tr}\left(P_{\Lambda} K P_{\Lambda}\right)<\infty$ for all $\Lambda \in$ $\mathcal{B}_{0}(X)$, where $P_{\Lambda}$ denotes the operator of multiplication by $\chi_{\Lambda}$. Hence, each operator $P_{\Lambda} \sqrt{K}$ is of Hilbert-Schmidt class. Following [23] (see also [12, Lemma A.4]), we conclude that $\sqrt{K}$ is an integral operator whose integral kernel, $\varkappa(x, y)$, satisfies

$$
\begin{equation*}
\int_{\Lambda} \int_{X} \nu(d x) \nu(d y) \varkappa(x, y)^{2}<\infty \quad \text { for all } \quad \Lambda \in \mathcal{B}_{0}(X) \tag{3.1}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\varkappa(x, \cdot) \in L^{2}(X, \nu) \quad \text { for } \nu \text {-a.a. } \quad x \in X \tag{3.2}
\end{equation*}
$$

Hence, $K$ is an integral operator whose integral kernel can be chosen as

$$
\begin{align*}
k(x, y) & =\int_{X} \varkappa(x, z) \varkappa(z, y) \nu(d z) \\
& =\int_{X} \varkappa(x, z) \varkappa(y, z) \nu(d z)=(\varkappa(x, \cdot), \varkappa(y, \cdot))_{L^{2}(X, \nu)} \tag{3.3}
\end{align*}
$$

We also have, for each $\Lambda \in \mathcal{B}_{0}(X)$,

$$
\begin{align*}
\operatorname{Tr}\left(P_{\Lambda} K P_{\Lambda}\right) & =\left\|\sqrt{K} P_{\Lambda}\right\|_{\mathrm{HS}}^{2} \\
& =\int_{\Lambda} \nu(d x) \int_{X} \nu(d y) \varkappa(x, y)^{2}=\int_{\Lambda} k(x, x) \nu(d x) \tag{3.4}
\end{align*}
$$

where $\|\cdot\|_{\text {HS }}$ denotes the Hilbert-Schmidt norm.
Proposition 3.1. There exists a random field $(Y(x))_{x \in X}$ on a probability space $(\Omega, \mathcal{A}, P)$ such that the mapping

$$
\begin{equation*}
X \times \Omega \ni(x, \omega) \mapsto Y(x, \omega) \tag{3.5}
\end{equation*}
$$

is measurable, and for $\nu$-a.a. $x \in X, Y(x)$ is a Gaussian random variable with mean 0 and such that

$$
\begin{equation*}
\mathbb{E}(Y(x) Y(y))=k(x, y) \quad \text { for } \nu^{\otimes 2}-a . a .(x, y) \in X^{2} \text { and } \nu \text {-a.a. } x=y \in X \tag{3.6}
\end{equation*}
$$

Remark 3.1. The statement of Proposition 3.1 is well-known if the integral kernel of the operator $K$ admits a continuous version (see e.g. Theorem 1.8 and p. 456 in [30]). In the latter case, $(Y(x))_{x \in X}$ is a Gaussian random field and formula (3.6) holds for all $(x, y) \in X^{2}$.

Proof of Proposition 3.1. Consider a standard triple of real Hilbert spaces

$$
H_{+} \subset H_{0}=L^{2}(X, \nu) \subset H_{-}
$$

Here the Hilbert space $H_{+}$is densely and continuously embedded into $H_{0}$, the inclusion operator $H_{+} \hookrightarrow H_{0}$ is of Hilbert-Schmidt class, and the Hilbert space $H_{-}$is the dual space of $H_{+}$with respect to the center space $H_{0}$ (see e.g. [2]).

Let $\mathbb{P}$ be the standard Gaussian measure on $H_{-}$, i.e., the probability measure on the Borel $\sigma$-algebra $\mathcal{B}\left(H_{-}\right)$which has Fourier transform

$$
\int_{H_{-}} e^{i\langle\omega, f\rangle} \mathbb{P}(d \omega)=\exp \left[-\frac{1}{2}\|f\|_{H_{0}}^{2}\right], \quad f \in H_{+},
$$

where $\langle\omega, f\rangle$ denotes the dual pairing between $\omega \in H_{-}$and $f \in H_{+}$. Then the mapping $H_{+} \ni f \rightarrow\langle\cdot, f\rangle$ can be extended by continuity to an isometry

$$
\begin{equation*}
I: H_{0} \rightarrow L^{2}\left(H_{-}, \mathbb{P}\right) \tag{3.7}
\end{equation*}
$$

For any $f \in H_{0}$ we denote $\langle\cdot, f\rangle:=I f$. Thus, for each $f \in H_{0},\langle\cdot, f\rangle$ is a (complex) Gaussian random variable with mean 0 and for any $f, g \in H_{0}$

$$
\begin{equation*}
\int_{H_{-}}\langle\omega, f\rangle\langle\omega, g\rangle \mathbb{P}(d \omega)=(f, g)_{L^{2}(X, \nu)} \tag{3.8}
\end{equation*}
$$

Thus, by (3.2), we set for $\nu$-a.a. $x \in X, \widetilde{Y}(x, \omega):=\langle\omega, k(x, \cdot)\rangle$. Hence $\tilde{Y}(x)$ is a Gaussian random variable and by (3.3) and (3.8), (3.6) holds.

Hence, it remains to prove that there exists a random field $Y=(Y(x))_{x \in X}$ for which the mapping (3.5) is measurable and such that $Y(x, \omega)=\widetilde{Y}(x, \omega)$ for $\nu \otimes \mathbb{P}$-a.a. $(x, \omega)$. To this end, we fix any $\Lambda \in \mathcal{B}_{0}(X)$ and denote by $\mathcal{B}(\Lambda)$ the trace $\sigma$-algebra of $\mathcal{B}(X)$ on $\Lambda$. We define a set $\mathcal{D}_{\Lambda}$ of the functions $u: \Lambda \times X \rightarrow \mathbb{R}$ of the form

$$
\begin{equation*}
u(x, y)=\sum_{i=1}^{n} \chi_{\Delta_{i}}(x) f_{i}(y) \tag{3.9}
\end{equation*}
$$

where $\Delta_{i} \in \mathcal{B}(\Lambda), f_{i} \in H_{+}, i=1, \ldots, n$. Define a linear mapping

$$
\begin{equation*}
I_{\Lambda}: \mathcal{D}_{\Lambda} \rightarrow L^{2}\left(\Lambda \times H_{-}, \nu \otimes \mathbb{P}\right) \tag{3.10}
\end{equation*}
$$

by setting, for each $u \in \mathcal{D}_{\Lambda}$ of the form (3.9),

$$
\left(I_{\Lambda} u\right)(x, \omega)=\sum_{i=1}^{n} \chi_{\Delta_{i}}(x)\left\langle\omega, f_{i}\right\rangle, \quad(x, \omega) \in \Lambda \times H_{-}
$$

Clearly, $I_{\Lambda}$ can be extended to an isometry

$$
I_{\Lambda}: L^{2}\left(\Lambda \times X, \nu^{\otimes 2}\right) \rightarrow L^{2}\left(\Lambda \times H_{-}, \nu \otimes \mathbb{P}\right)
$$

and we have $I_{\Lambda}=\mathbf{1}_{\Lambda} \otimes I$, where $\mathbf{1}_{\Lambda}$ is the identity operator in $L^{2}(\Lambda, \nu)$ and the operator $I$ is as in (3.7).

Fix any $u \in L^{2}\left(\Lambda \times X, \nu^{\otimes 2}\right)$. As easily seen, there exist a sequence $\left(u_{n}\right)_{n=1}^{\infty} \subset \mathcal{D}_{\Lambda}$ such that $u_{n} \rightarrow u$ in $L^{2}\left(\Lambda \times X, \nu^{\otimes 2}\right)$ and for $\nu$-a.a. $x \in \Lambda, u_{n}(x, \cdot) \rightarrow u(x, \cdot)$ in $L^{2}(X, \nu)$ Hence, for $\nu$-a.a. $x \in \Lambda, I_{\Lambda} u_{n}(x, \cdot) \rightarrow I_{\Lambda} u(x, \cdot)$ in $L^{2}\left(H_{-}, \mathbb{P}\right)$, which implies

$$
\begin{equation*}
\left(I_{\Lambda} u\right)(x, \omega)=\langle\omega, u(x, \cdot)\rangle \quad \text { for } \mathbb{P} \text {-a.a. } \omega \in H_{-} . \tag{3.11}
\end{equation*}
$$

Now, denote by $\varkappa_{\Lambda}$ the restriction of $\varkappa$ to the set $\Lambda \times X$. For $\nu$-a.a. $x \in \Lambda$, we define $Y_{\Lambda}(x):=\left(I_{\Lambda} \varkappa_{\Lambda}\right)(x, \cdot)$. Hence, by (3.11), for $\nu$-a.a. $x \in \Lambda, Y_{\Lambda}(x)=\widetilde{Y}(x) \mathbb{P}$-a.e. Finally, let $\left(\Lambda_{n}\right)_{n=1}^{\infty} \subset \mathcal{B}_{0}(X)$ be such that $\Lambda_{n} \cap \Lambda_{m}=\emptyset$ if $n \neq m$ and $\bigcup_{n=1}^{\infty} \Lambda_{n}=X$. Setting $Y(x):=Y_{\Lambda_{n}}(x)$ for $\nu$-a.a. $x \in \Lambda_{n}, n \in \mathbb{N}$, we conclude the statement.

Let $Y$ be a random field as in Proposition 3.1. For each $\Lambda \in \mathcal{B}_{0}(X)$, we have

$$
\begin{aligned}
\mathbb{E}\left(\int_{\Lambda} Y(x)^{2} \nu(d x)\right) & =\int_{\Lambda} \mathbb{E}\left(Y(x)^{2}\right) \nu(d x) \\
& =\int_{\Lambda} \nu(d x) \int_{X} \nu(d y) \varkappa(x, y)^{2}<\infty
\end{aligned}
$$

In particular, the function $Y(x)^{2}$ is locally $\nu$-integrable $\mathbb{P}$-a.s. Let $l \in \mathbb{N}$ and let $(\Omega, \mathcal{A}, \mathbb{P})$ be a probability space on which $l$ independent copies $Y_{1}, Y_{2}, \ldots, Y_{l}$ of a random field $Y$ as in Proposition 3.1 are defined. Denote by $\mu^{(l)}$ the Cox point process on $X$ with random intensity $g^{(l)}(x)=\sum_{i=1}^{l} Y_{i}(x)^{2}$, which is locally $\nu$-integrable $\mathbb{P}$-a.s. Thus, $\mu^{(l)}$ is the probability measure on $(\Gamma, \mathcal{B}(\Gamma))$ which satisfies

$$
\begin{equation*}
\int_{\Gamma} \mu^{(l)}(d \gamma) F(\gamma)=\int_{\Omega} \mathbb{P}(d \omega) \int_{\Gamma} \pi_{g^{(l)}(x, \omega) \nu(d x)}(d \gamma) F(\gamma) \tag{3.12}
\end{equation*}
$$

for each measurable function $F: \Gamma \rightarrow[0,+\infty]$. Here, for a locally $\nu$-integrable function $g: X \rightarrow[0,+\infty)$, we denote by $\pi_{g(x) \nu(d x)}$ the Poisson point process in $X$ with intensity measure $g(x) \nu(d x)$, see e.g [5]. This is the unique point process in $X$ which satisfies the Mecke identity

$$
\begin{equation*}
\int_{\Gamma} \pi_{g(x) \nu(d x)}(d \gamma) \int_{X} \gamma(d x) F(x, \gamma)=\int_{\Gamma} \pi_{g(x) \nu(d x)}(d \gamma) \int_{X} \nu(d x) g(x) F(x, \gamma \cup x) \tag{3.13}
\end{equation*}
$$

for each measurable $F: X \times \Gamma \rightarrow[0,+\infty]$. By (3.12) and (3.13) (compare with e.g. [27]), for each $l \in \mathbb{N}$, the point process $\mu^{(l)}$ satisfies condition $\left(\Sigma_{\nu}^{\prime}\right)$ and its Papangelou intensity is given by

$$
\begin{equation*}
r^{(l)}(x, \gamma)=\widetilde{\mathbb{E}}\left(g^{(l)}(x) \mid \mathcal{F}\right)(\gamma)=\widetilde{\mathbb{E}}\left(\sum_{i=1}^{l} Y_{i}(x)^{2} \mid \mathcal{F}\right)(\gamma) \tag{3.14}
\end{equation*}
$$

Here $\widetilde{\mathbb{E}}$ denotes the (conditional) expectation with respect to the probability measure

$$
\begin{equation*}
\widetilde{\mathbb{P}}(d \omega, d \gamma)=\widetilde{\mathbb{P}}(d \omega) \pi_{g^{(l)}(x, \omega) \nu(d x)}(d \gamma) \tag{3.15}
\end{equation*}
$$

on $\Omega \times \Gamma$, while $\mathcal{F}$ denotes the $\sigma$-algebra on $\Omega \times \Gamma$ generated by the mappings

$$
\Omega \times \Gamma \ni(\omega, \gamma) \rightarrow F(\gamma) \in \mathbb{R}
$$

where $F: \Gamma \rightarrow \mathbb{R}$ is measurable.
Recall that a point process $\mu$ in $X$ is said to have correlation functions if, for each $n \in \mathbb{N}$, there exist a non-negative, measurable, symmetric function $k_{\mu}^{(n)}$ on $X^{n}$ such that, for any measurable, symmetric function $f^{n}: X^{n} \rightarrow[0,+\infty]$,

$$
\begin{align*}
& \int_{\Gamma} \sum_{\left\{x_{1}, \ldots, x_{n}\right\} \subset \gamma} f^{(n)}\left(x_{1}, \ldots, x_{n}\right) \mu(d \gamma)  \tag{3.16}\\
& \quad=\frac{1}{n!} \int_{X^{n}} f^{(n)}\left(x_{1}, \ldots, x_{n}\right) k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right) \nu\left(d x_{1}\right) \cdots \nu\left(d x_{n}\right) .
\end{align*}
$$

As well known (e.g. [5]), for a locally $\nu$-integrable function $g: X \rightarrow[0,+\infty$ ), the Poisson point process $\pi_{g(x) \nu(d x)}$ has correlation functions

$$
\begin{equation*}
k_{\mu}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=g\left(x_{1}\right) \cdots g\left(x_{n}\right) \tag{3.17}
\end{equation*}
$$

Let us recall the notion of $\alpha$-permanent [31], called $\alpha$-determinant in [30]. For a square $\operatorname{matrix} A=\left(a_{i j}\right)_{i, j=1}^{n}$ and $\alpha \in \mathbb{R}$, we set

$$
\operatorname{per}_{\alpha} A:=\sum_{\sigma \in S_{n}} \alpha^{n-m(\sigma)} \prod_{i=1}^{n} a_{i \sigma(i)}
$$

where $S_{n}$ is the group of all permutations of $\{1, \ldots, n\}$ and $m(\sigma)$ denotes the number of cycles in $\sigma$. In particular, $\operatorname{per}_{1} A$ is the usual permanent of $A$, while $\operatorname{per}_{-1} A$ is the usual determinant of $A$. Analogously to [30, subsec. 6.4], we conclude from (3.12), (3.16) and (3.17) that the point process $\mu^{(l)}$ has correlation functions

$$
\begin{equation*}
k_{\mu^{(l)}}^{(n)}\left(x_{1}, \ldots, x_{n}\right)=\operatorname{per}_{\frac{l}{2}}\left(l k\left(x_{i}, x_{j}\right)\right)_{i, j=1}^{n} \quad \text { for } \nu^{\otimes n} \text {-a.a. }\left(x_{1}, \ldots, x_{n}\right) \in X^{n} . \tag{3.18}
\end{equation*}
$$

For $l=2$, the point process $\mu^{(2)}$ is often called a boson point process, see e.g. [5, 23]. Thus, we have proved the following
Proposition 3.2. For each $l \in \mathbb{N}$, there exists a point process $\mu^{(l)}$ in $X$ whose correlation functions are given by (3.18). The $\mu^{(l)}$ satisfies condition $\left(\Sigma_{\nu}^{\prime}\right)$ and its Papangelou intensity is given by (3.14).

Remark 3.2. Recall that in [30], under the same assumptions on the operator $K$, the existence of a point process with correlation functions (3.18) was proved for even $l \in \mathbb{N}$, and for odd $l \in \mathbb{N}$ the statement of Proposition 3.2 was proved under the additional assumption of continuity of the integral kernel $k(\cdot, \cdot)$.

We will now prove that, for a point process $\mu^{(l)}$ as in Proposition 3.2, Glauber and Kawasaki dynamics with coefficients (2.12), (2.13) and (2.14), respectively exist.
Theorem 3.1. (i) For each point process $\mu^{(l)}$ as in Proposition 3.2, the coefficients $d(x, \gamma)$ and $b(x, \gamma)$ defined by (2.12) and (2.13), satisfy conditions (2.3) and (2.7) and so statements (i) and (ii) of Theorem 2.1 hold, in particular, a corresponding Glauber dynamics exists.
(ii) Assume additionally that $k(x, x)$ is bounded outside a set $\Delta \in \mathcal{B}_{0}(X)$. Then for a point process $\mu^{(l)}$ as in Proposition 3.2, the coefficient $c(x, y, \gamma)$ defined by (2.14), satisfies (2.4), (2.5), (2.8) and (2.9), and so statements (i) and (ii) of Theorem 2.1 hold, in particular, a corresponding Kawasaki dynamics exists.

Proof. We start with the following
Lemma 3.1. For each $n \in \mathbb{N}$ and for $\nu$-a.a. $x \in X$

$$
\begin{equation*}
\int_{\Gamma} r(x, \gamma)^{n} \mu(d \gamma) \leq \frac{(2 n)!}{2^{n} n!} k(x, x)^{n} \tag{3.19}
\end{equation*}
$$

Proof. Using Jensen's inequality for conditional expectation and the formula for moments of a Gaussian measure (see e.g. [2, Chapter 2, Section 2, Lemma 2.1]), we have

$$
\begin{aligned}
\int_{\Gamma} r(x, \gamma)^{n} \mu(d \gamma) & =\widetilde{\mathbb{E}}\left(\widetilde{\mathbb{E}}\left(Y(x)^{2} \mid \mathcal{F}\right)^{n}\right) \leq \widetilde{\mathbb{E}}\left(\widetilde{\mathbb{E}}\left(Y(x)^{2 n} \mid \mathcal{F}\right)\right) \\
& =\widetilde{\mathbb{E}}\left(Y(x)^{2 n}\right) \leq \frac{(2 n)!}{2^{n} n!}\|\varkappa(x, \cdot)\|_{L^{2}(X, \nu)}^{2 n}=\frac{(2 n)!}{2^{n} n!} k(x, x)^{n}
\end{aligned}
$$

for $\nu$-a.a. $x \in X$.
We will only prove statement (ii) of Theorem 3.1, as the proof of statement (i) is similar and simper. Also, for simplicity of notation, we will only consider the case $l=1$ (for $l>1$ the proof being similar). We will also omit the upper index (1) from our notation. By (2.1) we have, for each $\Lambda \in \mathcal{B}_{0}(X)$,

$$
\begin{align*}
& \int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) \int_{X} \nu(d y) c(x, y, \gamma \backslash x)\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(y)\right) \\
& \quad=\int_{\Gamma} \mu(d \gamma) \int_{X} \nu(d x) \int_{X} \nu(d y) r(x, \gamma) c(x, y, \gamma)\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(y)\right)  \tag{3.20}\\
& \quad=\int_{\Gamma} \mu(d \gamma) \int_{X} \nu(d x) \int_{X} \nu(d y) a(x, y) r(x, \gamma)^{s} r(y, \gamma)^{s} \chi_{\{r>0\}}(x, \gamma)
\end{align*}
$$

$$
\begin{aligned}
& \times \chi_{\{r>0\}}(y, \gamma)\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(y)\right) \\
& \leq \int_{\Gamma} \mu(d \gamma) \int_{X} \nu(d x) \int_{X} \nu(d y) a(x, y) r(x, \gamma)^{s} r(y, \gamma)^{s}\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(y)\right) \\
& =2 \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) r(x, \gamma)^{s} r(y, \gamma)^{s} \\
& \leq 2 \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y)(1+r(x, \gamma))(1+r(y, \gamma))
\end{aligned}
$$

By (2.15)

$$
\begin{equation*}
\int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y)<\infty \tag{3.21}
\end{equation*}
$$

Below, $C_{i}, i=1,2,3, \ldots$, will denote positive constants whose explicit values are not important for us. We have, by (2.15)

$$
\begin{align*}
\int_{\Gamma} & \mu(d \gamma) \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) r(x, \gamma) \\
& =\int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) r(x, \gamma)\left(\int_{X} \nu(d y) a(x, y)\right)  \tag{3.22}\\
& \leq C_{1} \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) r(x, \gamma) \\
& =C_{1} \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \gamma(d x)=C_{1} \int_{\Lambda} k(x, x) \nu(d x)<\infty
\end{align*}
$$

Next, by (3.14)

$$
\begin{align*}
\int_{\Gamma} & \mu(d \gamma) \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) r(y, \gamma) \\
& =\int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) \int_{\Gamma} \mu(d \gamma) r(y, \gamma) \\
& =\int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) k(y, y)  \tag{3.23}\\
& =\int_{\Lambda} \nu(d x) \int_{\Delta} \nu(d y) a(x, y) k(y, y)+\int_{\Lambda} \nu(d x) \int_{\Delta^{c}} \nu(d y) a(x, y) k(y, y) \\
& \leq C_{2} \int_{\Lambda} \nu(d x) \int_{\Delta} \nu(d y) k(y, y)+C_{3} \int_{\Lambda} \nu(d x) \int_{\Delta^{c}} \nu(d y) a(x, y)<\infty,
\end{align*}
$$

where we used that the function $a$ is bounded and $k(y, y)$ is bounded on $\Delta^{c}$. Analogously, using Lemma 3.1, we have

$$
\begin{align*}
& \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) r(x, \gamma) r(y, \gamma) \\
& \quad \leq \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y)\|r(x, \cdot)\|_{L^{2}(\mu)}\|r(y, \cdot)\|_{L^{2}(\mu)} \\
& \quad \leq C_{4} \int_{\Lambda} \nu(d x) \int_{X} \nu(d y) a(x, y) k(x, x) k(y, y)  \tag{3.24}\\
& \quad \leq C_{5} \int_{\Lambda} \nu(d x) k(x, x) \int_{\Delta} \nu(d y) k(y, y) \\
& \quad+C_{6} \int_{\Lambda} \nu(d x) k(x, x) \int_{\Delta^{c}} \nu(d y) a(x, y)<\infty
\end{align*}
$$

Thus, by $(3.20)-(3.24)$, the theorem is proved.

Theorem 3.2. (i) Let $s \in\left[\frac{1}{2}, 1\right]$, and let the conditions of Theorem 3.1 (i) be satisfied. Then the coefficients $d(x, \gamma)$ and $b(x, \gamma)$ defined by (2.12) and (2.13), satisfy condition (2.10). Thus, $\mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right) \subset D\left(L_{\mathrm{G}}\right)$, and for each $F \in \mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right), L_{\mathrm{G}} F$ is given by formula (1.1).
(ii) Let $s \in\left[\frac{1}{2}, 1\right]$, and let the conditions of Theorem 3.1 (ii) be satisfied. Further assume that either

$$
\begin{equation*}
\forall \Lambda \in \mathcal{B}_{0}(X) \exists \Lambda^{\prime} \in \mathcal{B}_{0}(X) \forall x \in \Lambda \forall y \in\left(\Lambda^{\prime}\right)^{c}: \quad a(x, y)=0 \tag{3.25}
\end{equation*}
$$

or

$$
\begin{equation*}
\int_{\Delta} k(x, x)^{2} \nu(d x)<\infty \tag{3.26}
\end{equation*}
$$

where $\Delta$ is as in Theorem 3.1 (ii). Then the coefficient $c(x, y, \gamma)$ defined by (2.14), satisfies condition (2.11). Thus, $\mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right) \subset D\left(L_{\mathrm{K}}\right)$, and for each $F \in \mathcal{F} C_{\mathrm{b}}\left(C_{0}(X), \Gamma\right)$, $L_{\mathrm{K}} F$ is given by formula (1.2).

Remark 3.3. If $X=\mathbb{R}^{d}$ and the function $a$ is as in Remark 2.3, then condition (3.25) means that the function $\tilde{a}$ has a compact support.

Proof of Theorem 3.2. We again prove only the part related to Kawasaki dynamics and only in the case $l=1$, omitting the upper index (1) from our notation. We first assume that (3.25) is satisfied. Since the function $a$ is bounded and satisfies (3.25), it suffices to show that, for each $\Lambda \in \mathcal{B}_{0}(X)$,

$$
\begin{equation*}
\int_{\Lambda} \gamma(d x) \int_{\Lambda} \nu(d y) r(x, \gamma \backslash x)^{s-1} r(y, \gamma \backslash x)^{s} \chi_{\{r>0\}}(x, \gamma \backslash x) \chi_{\{r>0\}}(y, \gamma \backslash x) \in L^{2}(\mu) \tag{3.27}
\end{equation*}
$$

We note that, for $s \in\left[\frac{1}{2}, 1\right], 2 s-1 \in[0,1]$. Therefore, by the Cauchy inequality, we have

$$
\begin{align*}
& \int_{\Gamma} \mu(d \gamma)\left(\int_{\Lambda} \gamma(d x) r(x, \gamma \backslash x)^{s-1} \chi_{\{r>0\}}(x, \gamma \backslash x)\right. \\
& \left.\quad \times \int_{\Lambda} \nu(d y) r(y, \gamma \backslash x)^{s} \chi_{\{r>0\}}(y, \gamma \backslash x)\right)^{2} \\
& \quad \leq \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \gamma(d x) r(x, \gamma \backslash x)^{2(s-1)} \chi_{\{r>0\}}(x, \gamma \backslash x) \\
& \quad \times\left(\int_{\Lambda} \nu(d y) r(y, \gamma \backslash x)^{s} \chi_{\{r>0\}}(y, \gamma \backslash x)\right)^{2} \gamma(\Lambda) \\
& \quad=\int_{\Gamma} \mu(d \gamma) \int_{\Lambda} \nu(d x) r(x, \gamma)^{2 s-1} \chi_{\{r>0\}}(x, \gamma)  \tag{3.28}\\
& \quad \times\left(\int_{\Lambda} \nu(d y) r(y, \gamma)^{s} \chi_{\{r>0\}}(y, \gamma)\right)^{2}(\gamma(\Lambda)+1) \\
& \quad \leq \int_{\Gamma} \mu(d \gamma)\left(\int_{\Lambda} \nu(d x)(1+r(x, \gamma))\right)^{3}(\gamma(\Lambda)+1) \\
& \quad \leq\left(\int_{\Gamma} \mu(d \gamma)\left(\int_{\Lambda} \nu(d x)(1+r(x, \gamma))\right)^{6}\right)^{1 / 2}\left(\int_{\Gamma} \mu(d \gamma)(\gamma(\Lambda)+1)^{2}\right)^{1 / 2}
\end{align*}
$$

By Lemma 3.1, we have, for each $n \in \mathbb{N}$,

$$
\begin{align*}
& \int_{\Gamma} \mu(d \gamma)\left(\int_{\Lambda} \nu(d x) r(x, \gamma)\right)^{n}=\int_{\Lambda} \nu\left(d x_{1}\right) \cdots \int_{\Lambda} \nu\left(d x_{n}\right) \int_{\Gamma} \mu(d \gamma) r\left(x_{1}, \gamma\right) \cdots r\left(x_{n}, \gamma\right)  \tag{3.29}\\
& \quad \leq \int_{\Lambda} \nu\left(d x_{1}\right) \cdots \int_{\Lambda} \nu\left(d x_{n}\right)\left\|r\left(x_{1}, \cdot\right)\right\|_{L^{n}(\mu)} \cdots\left\|r\left(x_{n}, \cdot\right)\right\|_{L^{n}(\mu)} \\
& \quad \leq \frac{(2 n)!}{2^{n} n!}\left(\int_{\Lambda} \nu(d x) k(x, x)\right)^{n}<\infty
\end{align*}
$$

Now, (3.27) follows from (3.28) and (3.29).
Next, we assume that (3.26) is satisfied. We fix $\Lambda \in \mathcal{B}_{0}(X)$ and denote

$$
u(x, y):=a(x, y)\left(\chi_{\Lambda}(x)+\chi_{\Lambda}(x)\right)
$$

Then, by the Cauchy inequality,

$$
\begin{aligned}
& \int_{\Gamma} \mu(d \gamma)\left(\int_{X} \gamma(d x) \int_{X} \nu(d y) u(x, y) r(x, \gamma \backslash x)^{s-1} \chi_{\{r>0\}}(x, \gamma \backslash x)\right. \\
& \left.\quad \times r(y, \gamma \backslash x)^{s} \chi_{\{r>0\}}(y, \gamma \backslash x)\right)^{2} \\
& \quad \leq \int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) \int_{X} \nu(d y) u(x, y) r(x, \gamma \backslash x)^{2(s-1)} \chi_{\{r>0\}}(x, \gamma \backslash x) \\
& \quad \times r(y, \gamma \backslash x)^{2 s} \chi_{\{r>0\}}(y, \gamma \backslash x) \int_{X} \gamma\left(d x^{\prime}\right) \int_{X} \nu\left(d y^{\prime}\right) u\left(x^{\prime}, y^{\prime}\right) \\
& \quad=\int_{\Gamma} \mu(d \gamma) \int_{X} \nu(d x) \int_{X} \nu(d y) u(x, y) r(x, \gamma)^{2 s-1} \chi_{\{r>0\}}(x, \gamma) \\
& \quad \times r(y, \gamma)^{2 s} \chi_{\{r>0\}}(y, \gamma) \int_{X}\left(\gamma+\varepsilon_{x}\right)\left(d x^{\prime}\right) \int_{X} \nu\left(d y^{\prime}\right) u\left(x^{\prime}, y^{\prime}\right) \\
& \quad \leq \int_{\Gamma} \mu(d \gamma) \int_{X} \nu(d x) \int_{X} \nu(d y) u(x, y)(1+r(x, \gamma))\left(1+r(y, \gamma)^{2}\right. \\
& \quad \times\left(\int_{X} \gamma\left(d x^{\prime}\right) \int_{X} \nu\left(d y^{\prime}\right) u\left(x^{\prime}, y^{\prime}\right)+\int_{X} \nu\left(d y^{\prime}\right) u\left(x, y^{\prime}\right)\right)
\end{aligned}
$$

By (2.15), it suffices to prove that

$$
\begin{align*}
\int_{\Gamma} \mu(d \gamma) & \left(\int_{X} \nu(d x) \int_{X} \nu(d y) u(x, y)(1+r(x, \gamma))\left(1+r(y, \gamma)^{2}\right)\right)^{2}<\infty  \tag{3.30}\\
& \int_{\Gamma} \mu(d \gamma)\left(\int_{X} \gamma(d x) \int_{X} \nu(d y) u(x, y)\right)^{2}<\infty \tag{3.31}
\end{align*}
$$

We first to prove (3.31). We have, by Proposition 3.2,

$$
\begin{aligned}
\int_{\Gamma} & \left(\int_{X} \gamma(d x) \int_{X} \nu(d y) u(x, y)\right)^{2} \\
& =\int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right) \int_{\Gamma} \mu(d \gamma) \int_{X} \gamma(d x) \int_{X} \gamma\left(d x^{\prime}\right) u(x, y) u\left(x^{\prime}, y^{\prime}\right) \\
& =\int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right) \int_{\Gamma} \mu(d \gamma)\left(\int_{X} \gamma(d x) u(x, y) u\left(x, y^{\prime}\right)\right. \\
& \left.+\int_{X} \gamma(d x) \int_{X}\left(\gamma-\varepsilon_{x}\right)\left(d x^{\prime}\right) u(x, y) u\left(x^{\prime}, y^{\prime}\right)\right) \\
& =\int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right)\left(\int_{X} \nu(d x) k(x, x) u(x, y) u\left(x, y^{\prime}\right)\right. \\
& \left.+\int_{X} \nu(d x) \int_{X} \nu\left(d x^{\prime}\right)\left(\frac{1}{2} k\left(x, x^{\prime}\right)^{2}+k(x, x) k\left(x^{\prime}, x^{\prime}\right)\right) u(x, y) u\left(x^{\prime}, y^{\prime}\right)\right) \\
& \leq \int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right)\left(\int_{X} \nu(d x) k(x, x) u(x, y) u\left(x, y^{\prime}\right)\right. \\
& \left.+\int_{X} \nu(d x) \int_{X} \nu\left(d x^{\prime}\right) \frac{3}{2} k(x, x) k\left(x^{\prime}, x^{\prime}\right) u(x, y) u\left(x^{\prime}, y^{\prime}\right)\right) \\
& =\int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right) \int_{X} \nu(d x) k(x, x) u(x, y) u\left(x, y^{\prime}\right)
\end{aligned}
$$

$$
\begin{aligned}
& +\frac{3}{2}\left(\int_{X} \nu(d y) \int_{X} \nu(d x) k(x, x) u(x, y)\right)^{2} \\
& \leq \int_{\Delta} \nu(d x) k(x, x)\left(\int_{X} \nu(d y) u(x, y)\right)^{2} \\
& +C_{7} \int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right) \int_{X} \nu(d x) u(x, y) u\left(x, y^{\prime}\right) \\
& +\frac{3}{2}\left(\int_{\Delta} \nu(d x) k(x, x) \int_{X} \nu(d y) u(x, y)+C_{7} \int_{X} \nu(d y) \int_{X} \nu(d x) u(x, y)\right)^{2}<\infty .
\end{aligned}
$$

Next, we prove (3.30). By Lemma 3.1 and (3.26), we have

$$
\begin{aligned}
\int_{\Gamma} & \mu(d \gamma)\left(\int_{X} \nu(d x) \int_{X} \nu(d y) u(x, y)(1+r(x, \gamma))\left(1+r(y, \gamma)^{2}\right)\right)^{2} \\
& =\int_{X} \nu(d x) \int_{X} \nu\left(d x^{\prime}\right) \int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right) u(x, y) u\left(x^{\prime}, y^{\prime}\right) \\
& \times \int_{\Gamma} \mu(d \gamma)(1+r(x, \gamma))\left(1+r\left(x^{\prime}, \gamma\right)\right)\left(1+r(y, \gamma)^{2}\right)\left(1+r\left(y^{\prime}, \gamma\right)^{2}\right) \\
& \leq \int_{X} \nu(d x) \int_{X} \nu\left(d x^{\prime}\right) \int_{X} \nu(d y) \int_{X} \nu\left(d y^{\prime}\right) u(x, y) u\left(x^{\prime}, y^{\prime}\right)\left(1+\|r(x, \cdot)\|_{L^{4}(\mu)}\right) \\
& \times\left(1+\left\|r\left(x^{\prime}, \cdot\right)\right\|_{L^{4}(\mu)}\right)\left(1+\left\|r(y, \cdot)^{2}\right\|_{L^{4}(\mu)}\right)\left(1+\left\|r\left(y^{\prime}, \cdot\right)^{2}\right\|_{L^{4}(\mu)}\right) \\
& \leq C_{8}\left(\int_{X} \nu(d x) \int_{X} \nu(d y) u(x, y)(1+k(x, x))\left(1+k(y, y)^{2}\right)\right)^{2}<\infty
\end{aligned}
$$

Thus, the theorem is proved.

## 4. Diffusion approximation

¿From now on, we set $X=\mathbb{R}^{d}, d \in \mathbb{N}$, and $\nu$ to be Lebesgue measure. We will show that, under an appropriate scaling, the Dirichlet form of the Kawasaki dynamics converges to a Dirichlet form which identifies a diffusion process on $\Gamma$ having a permanental point process $\mu^{(l)}$ as a symmetrizing measure. The way we scale the Kawasaki dynamics will be similar to the ansatz of [15].

We denote by $\mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)$ the space of all functions of the form (2.6) where $N \in \mathbb{N}, \varphi_{1}, \ldots, \varphi_{N} \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$ and $g \in C_{\mathrm{b}}^{\infty}\left(\mathbb{R}^{N}\right)$. Here, $C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$ denotes the space of smooth functions on $\mathbb{R}^{d}$ with compact support, and $C_{\mathrm{b}}^{\infty}\left(\mathbb{R}^{N}\right)$ denotes the space of all smooth bounded functions on $\mathbb{R}^{N}$ whose all derivatives are bounded. Clearly,

$$
\mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right) \subset \mathcal{F} C_{\mathrm{b}}\left(C_{0}\left(\mathbb{R}^{d}\right), \Gamma\right)
$$

and the set $\mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)$ is a core for the Dirichlet form $\left(\mathcal{E}_{\mathrm{K}}, D\left(\mathcal{E}_{\mathrm{K}}\right)\right)$.
We fix $s=1 / 2$. Let us assume that the function $a(x, y)$ is as in Remark 2.3. Thus, the coefficient $c(x, y, \gamma)$ has the form

$$
\begin{equation*}
c(x, y, \gamma)=a(x-y) r(x, \gamma)^{-1 / 2} r(y, \gamma)^{1 / 2} \chi_{\{r>0\}}(x, \gamma) \chi_{\{r>0\}}(y, \gamma) \tag{4.1}
\end{equation*}
$$

Note that $y-x$ describes the change of the position of a particle which hops from $x$ to $y$. We now scale the function $a$ as follows: for each $\varepsilon>0$, we denote

$$
\begin{equation*}
a_{\varepsilon}(x):=\varepsilon^{-d-2} a(x / \varepsilon), \quad x \in \mathbb{R}^{d} . \tag{4.2}
\end{equation*}
$$

The Dirichlet form $\left(\mathcal{E}_{\mathrm{K}}, D\left(\mathcal{E}_{\mathrm{K}}\right)\right)$ which corresponds to the choice of function $a$ as in (4.2) will be denoted by $\left(\mathcal{E}_{\varepsilon}, D\left(\mathcal{E}_{\varepsilon}\right)\right)$.

Theorem 4.1. Assume that the function a has compact support, and the value a $(x)$ only depends on $|x|$, i.e., $a(x)=\tilde{a}(|x|)$ for some function $\tilde{a}:[0, \infty) \rightarrow \mathbb{R}$. Further assume
that the function $\varkappa(x, y)$ has the form $\varkappa(x-y)$ for some $\varkappa: \mathbb{R}^{d} \rightarrow \mathbb{C}$, and

$$
\begin{equation*}
\lim _{y \rightarrow 0} \int_{\mathbb{R}^{d}}(\varkappa(x)-\varkappa(x+y))^{2} d x=0 \tag{4.3}
\end{equation*}
$$

For each $l \in \mathbb{N}$, define a bilinear form $\left(\mathcal{E}_{0}, \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)\right)$ by

$$
\begin{equation*}
\mathcal{E}_{0}(F, G):=c \int_{\Gamma} \mu^{(l)}(d \gamma) \int_{\mathbb{R}^{d}} d x r(x, \gamma)\left\langle\nabla_{x} F(\gamma \cup x), \nabla_{x} G(\gamma \cup x)\right\rangle \tag{4.4}
\end{equation*}
$$

Here

$$
c:=\frac{1}{2} \int_{\mathbb{R}^{d}} a(x) x_{1}^{2} d x
$$

( $x_{1}$ denoting the first coordinate of $x \in \mathbb{R}^{d}$ ), $\nabla_{x}$ denotes the gradient in the $x$ variable, and $\langle\cdot, \cdot\rangle$ stands for the scalar product in $\mathbb{R}^{d}$. Then, for any $F, G \in \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)$,

$$
\mathcal{E}_{\varepsilon}(F, G) \rightarrow \mathcal{E}_{0}(F, G) \quad \text { as } \varepsilon \rightarrow 0
$$

Remark 4.1. Assume that the function $\varkappa$ is differentiable on $\mathbb{R}^{d}$. Denote

$$
K(x, \delta):=\sup _{y \in B(x, \delta)}|\nabla \varkappa(y)|, \quad x \in \mathbb{R}^{d}, \quad \delta>0
$$

Here $B(x, \delta)$ denotes the closed ball in $\mathbb{R}^{d}$ centered at $x$ and of radius $\delta$. Assume that, for some $\delta>0$,

$$
\begin{equation*}
K(\cdot, \delta) \in L^{2}\left(\mathbb{R}^{d}, d x\right) \tag{4.5}
\end{equation*}
$$

Then condition (4.3) is clearly satisfied. Note that condition (4.5) is slightly stronger than the condition $|\nabla \varkappa| \in L^{2}\left(\mathbb{R}^{d}, d x\right)$.

Proof of Theorem 4.1. Again we will only present the proof in the case $l=1$, omitting the upper index (1). We start with the following

Lemma 4.1. Fix any $\Lambda \in \mathcal{B}_{0}\left(\mathbb{R}^{d}\right)$ and $\alpha \in(0,1]$. Then, under the conditions of Theorem 4.1,

$$
r(x+\varepsilon y, \gamma)^{\alpha} \rightarrow r(x, \gamma)^{\alpha} \quad \text { in } \quad L^{2}\left(\Gamma \times \Lambda \times \mathbb{R}^{d}, \mu(d \gamma) d x d y a(y)\right) \quad \text { as } \quad \varepsilon \rightarrow 0
$$

Proof. We first prove the statement for $\alpha=1$. Thus, equivalently we have to prove that (4.6) $r(x+\varepsilon y, \gamma) \rightarrow r(x, \gamma) \quad$ in $\quad L^{2}\left(\Omega \times \Gamma \times \Lambda \times \mathbb{R}^{d}, \tilde{\mathbb{P}}(d \omega, d \gamma) d x d y a(y)\right) \quad$ as $\quad \varepsilon \rightarrow 0$.

We have, using Jensen's inequality for conditional expectation,

$$
\begin{align*}
& \int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) \int_{\Omega \times \Gamma} \tilde{\mathbb{P}}(d \omega, d \gamma)(r(x+\varepsilon y)-r(x, \gamma))^{2} \\
& \quad=\int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) \int_{\Omega \times \Gamma} \tilde{\mathbb{P}}(d \omega, d \gamma) \tilde{\mathbb{E}}\left(Y(x+\varepsilon y)^{2}-Y(x)^{2} \mid \mathcal{F}\right)^{2}  \tag{4.7}\\
& \quad \leq \int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) \int_{\Omega \times \Gamma} \tilde{\mathbb{P}}(d \omega, d \gamma)\left(Y(x+\varepsilon y)^{2}-Y(x)^{2}\right)^{2} \\
& \quad=\int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) \int_{\Omega} d \mathbb{P}\left(Y(x+\varepsilon y)^{4}+Y(x)^{4}-2 Y(x+\varepsilon y)^{2} Y(x)^{2}\right)
\end{align*}
$$

Using the formula for moments of a Gaussian measure, we have

$$
\begin{align*}
& \int_{\Omega} Y(x+\varepsilon y)^{4} d \mathbb{P} \\
& \quad=3\left(\int_{\mathbb{R}^{d}} \varkappa(x+\varepsilon y-u)^{2} d u\right)^{2}  \tag{4.8}\\
& \quad=3\left(\int_{\mathbb{R}^{d}} \varkappa(x-u)^{2} d u\right)^{2} \\
& \quad=\int_{\Omega} Y(x)^{4} d \mathbb{P}
\end{align*}
$$

Analogously, using condition (4.3) and the dominated convergence theorem, we get

$$
\begin{align*}
\int_{\Lambda} & d x \int_{\mathbb{R}^{d}} d y a(y) \int_{\Omega} d \mathbb{P} Y(x+\varepsilon y)^{2} Y(x)^{2} \\
& =\int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y)\left[\int_{\mathbb{R}^{d}} \varkappa(x+\varepsilon y-u)^{2} d u \cdot \int_{\mathbb{R}^{d}} \varkappa\left(x-u^{\prime}\right)^{2} d u^{\prime}\right. \\
& \left.+2\left(\int_{\mathbb{R}^{d}} \varkappa(x+\varepsilon y-u) \varkappa(x-u) d u\right)^{2}\right]  \tag{4.9}\\
& \rightarrow \int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) \int_{\Omega} d \mathbb{P} Y(x)^{4} \quad \text { as } \quad \varepsilon \rightarrow 0
\end{align*}
$$

By (4.7)-(4.9), statement (4.6) follows.
To prove the result for $\alpha \in(0,1)$, it is now sufficient to show the following
$\operatorname{Claim}$. Let $(\mathbf{A}, \mathcal{A}, m)$ be a measure space and let $m(A)<\infty$. Let $f_{\varepsilon} \in L^{2}(m), f_{\varepsilon} \geq 0$, $\varepsilon \in[-1,1]$, and let $f_{\varepsilon} \rightarrow f_{0}$ in $L^{2}(m)$ as $\varepsilon \rightarrow 0$. Then, for each $\alpha \in(0,1), f_{\varepsilon}^{\alpha} \rightarrow f_{0}^{\alpha}$ in $L^{2}(m)$ as $\varepsilon \rightarrow 0$.

By e.g. [1, Theorems 21.2 and 21.4], $f_{\varepsilon} \rightarrow f_{0}$ in $L^{2}(m)$ implies that
(i) $f_{\varepsilon} \rightarrow f_{0}$ in measure;
(ii) $\sup _{\varepsilon \in[-1,1]} \int f_{\varepsilon}^{2} d m<\infty$;
(iii) For each $\theta>0$ there exist $h \in L^{1}(m)$ and $\delta>0$ such that, for all $0<|\varepsilon| \leq 1$ and for each $A \in \mathcal{A}$

$$
\int_{A} h d m \leq \delta \Rightarrow \int_{A} f_{\varepsilon}^{2} d m \leq \theta
$$

Hence, for $\alpha \in(0,1)$, we get
a) $f_{\varepsilon}^{\alpha} \rightarrow f_{0}^{\alpha}$ in measure;
b) $\sup _{\varepsilon \in[-1,1]} \int f_{\varepsilon}^{2 \alpha} d m \leq \sup _{\varepsilon \in[-1,1]} \int\left(1+f_{\varepsilon}^{2}\right) d m<\infty$;
c) Let $\theta, h$, and $\delta$ be as in (iii). Set $h^{\prime}:=h+\frac{\delta}{\theta}$. Clearly, $h \in L^{1}(m)$. Assume that, for some $A \in \mathcal{A}, \int_{A} h^{\prime} d m \leq \delta$. Hence $\int_{A} h d m \leq \delta$, and therefore $\int_{A} f_{\varepsilon}^{2} d m \leq \delta$ for all $0<|\varepsilon| \leq 1$. Furthermore, we get $\int_{A} \frac{\delta}{\theta} d m \leq \delta$, and therefore $m(A) \leq \theta$. Now

$$
\int_{A} f_{\varepsilon}^{2 \alpha} d m \leq \int_{A}\left(1+f_{\varepsilon}^{2}\right) d m \leq 2 \theta
$$

Applying again [1, Theorems 21.2 and 21.4], we conclude the claim.
Fix any $F \in \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)$. We have

$$
\begin{aligned}
& \mathcal{E}_{\varepsilon}(F, F) \\
& =\frac{1}{2} \int_{\Gamma} \mu(d \gamma) \int_{\mathbb{R}^{d}} d x \int_{\mathbb{R}^{d}} d y \varepsilon^{-d-2} a((x-y) / \varepsilon) r(x, \gamma)^{1 / 2} r(y, \gamma)^{1 / 2}(F(\gamma \cup x)-F(\gamma \cup y))^{2}
\end{aligned}
$$

$$
=\frac{1}{2} \int_{\Gamma} \mu(d \gamma) \int_{\mathbb{R}^{d}} d x \int_{\mathbb{R}^{d}} d y a(y) r(x+\varepsilon y, \gamma)^{1 / 2} r(x, \gamma)^{1 / 2}\left(\frac{F(\gamma \cup\{x+\varepsilon y\})-F(\gamma \cup x)}{\varepsilon}\right)^{2}
$$

Assume that $0<|\varepsilon| \leq 1$. Noting that the function $F$ is local (i.e., there exists $\Delta \in \mathcal{B}_{0}\left(\mathbb{R}^{d}\right)$ such that $F(\gamma)=F\left(\gamma_{\Delta}\right)$ for all $\left.\gamma \in \Gamma\right)$ and that the function $a$ has a compact support, we conclude that there exists $\Lambda \in \mathcal{B}_{0}\left(\mathbb{R}^{d}\right)$ such that

$$
\begin{align*}
\mathcal{E}_{\varepsilon}(F, F) & =\int_{\Gamma} \mu(d \gamma) \int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) r(x+\varepsilon y, \gamma)^{1 / 2} r(x, \gamma)^{1 / 2} \\
& \times\left(\frac{F(\gamma \cup\{x+\varepsilon y\})-F(\gamma \cup x)}{\varepsilon}\right)^{2} \tag{4.10}
\end{align*}
$$

By the dominated convergence theorem

$$
\begin{equation*}
r(x, \gamma)^{1 / 2}\left(\frac{F(\gamma \cup\{x+\varepsilon y\})-F(\gamma \cup x)}{\varepsilon}\right)^{2} \rightarrow r(x, \gamma)^{1 / 2}\left\langle\nabla_{x} F(\gamma \cup x), y\right\rangle^{2} \tag{4.11}
\end{equation*}
$$

in $L^{2}\left(\Gamma \times \Lambda \times \mathbb{R}^{d}, \mu(d \gamma) d x d y a(y)\right)$ as $\varepsilon \rightarrow 0$. By Lemma 4.1 with $\alpha=1 / 2$, (4.10) and

$$
\begin{equation*}
\mathcal{E}_{\varepsilon}(F, F) \rightarrow \frac{1}{2} \int_{\Gamma} \mu(d \gamma) \int_{\Lambda} d x \int_{\mathbb{R}^{d}} d y a(y) r(x, \gamma)\left\langle\nabla_{x} F(\gamma \cup x), y\right\rangle^{2} \tag{4.11}
\end{equation*}
$$

Since $a(y)=\tilde{a}(|y|)$, for any $i, j \in\{1, \ldots, d\}, i \neq j$, we have

$$
\int_{\mathbb{R}^{d}} a(y) y_{i} y_{j} d y=0
$$

and

$$
c=\frac{1}{2} \int_{\mathbb{R}^{d}} a(y) y_{i}^{2} d y, \quad i=1, \ldots, d
$$

Therefore, by (4.12),

$$
\mathcal{E}_{\varepsilon}(F, F) \rightarrow c \int_{\Gamma} \mu(d \gamma) \int_{\mathbb{R}^{d}} d x r(x, \gamma)\left|\nabla_{x} F(\gamma \cup x)\right|^{2}
$$

From here the theorem follows by the polarization identity.
We will now show that the limiting form $\left(\mathcal{E}_{0}, \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)\right)$ is closable and its closure identifies a diffusion process.

In what follows, we will assume that the conditions of Theorem 4.1 are satisfied. We have

$$
\begin{aligned}
k(x, y) & =\int_{\mathbb{R}^{d}} \varkappa(x-u) \varkappa(y-u) d u \\
& =\int_{\mathbb{R}^{d}} \varkappa(u-y) \varkappa(u-x) d u=\int_{\mathbb{R}^{d}} \varkappa(u) \varkappa(u+y-x) d u .
\end{aligned}
$$

Hence, by (4.3), the function $k(x, y)$ is continuous on $\left(\mathbb{R}^{d}\right)^{2}$. Thus, by Remark 3.1, $(Y(x))_{x \in X}$ is a Gaussian random field and formula (3.6) holds for all $(x, y) \in\left(\mathbb{R}^{d}\right)^{2}$.

Consider the semimetric

$$
\begin{align*}
D(x, y): & =\frac{1}{2}\left(\int_{\Omega}(Y(x)-Y(y))^{2} d \mathbb{P}\right)^{1 / 2} \\
& =\frac{1}{2}(k(x, x)+k(y, y)-2 k(x, y))^{1 / 2}  \tag{4.13}\\
& =\left(\int_{\mathbb{R}^{d}} \varkappa(u)(\varkappa(u)-\varkappa(u+y-x)) d u\right)^{1 / 2}, \quad x, y \in \mathbb{R}^{d} .
\end{align*}
$$

The associated metric entropy $H(D, \delta)$ is defined as $H(D, \delta):=\log N(D, \delta)$, where $N(D, \delta)$ is the minimal number of points in a $\delta$-net in $B(0,1)=\left\{x \in \mathbb{R}^{d}| | x \mid \leq 1\right\}$
with respect to the semimetric $D$, i.e., points $x_{i}$ such that the open balls centered at $x_{i}$ and of radius $\delta$ (with respect to $D$ ) cover $B(0,1)$. The expression

$$
J(D):=\int_{0}^{1} \sqrt{H(D, \delta)} d \delta
$$

is called the Dudley integral. The following result holds, see e.g. [4, Corollary 7.1.4] and the references therein.
Theorem 4.2. Assume that $J(D)<\infty$. Then the Gaussian random field $(Y(x))_{x \in \mathbb{R}^{d}}$ has a continuous modification.

Remark 4.2. Let $\varkappa$ be as in Remark 4.1. Then, by (4.13), for any $x, y \in B(0,1)$

$$
\begin{aligned}
D(x, y)^{2} & \leq\|\varkappa(\cdot)\|_{L^{2}\left(\mathbb{R}^{d}, d x\right)}\left(\int_{\mathbb{R}^{d}}(\varkappa(u)-\varkappa(u+y-x))^{2} d u\right)^{1 / 2} \\
& \leq\|\varkappa(\cdot)\|_{L^{2}\left(\mathbb{R}^{d}, d x\right)}\|K(\cdot, 2)\|_{L^{2}\left(\mathbb{R}^{d}, d x\right)}|y-x|
\end{aligned}
$$

where we assumed that $K(\cdot, 2) \in L^{2}\left(\mathbb{R}^{d}, d x\right)$. Then $J(D)<\infty$, see e.g. [4, Example 7.1.5].
Denote by $\ddot{\Gamma}$ the space of all multiple configurations in $\mathbb{R}^{d}$ : Thus, $\ddot{\Gamma}$ is the set of all Radon $\mathbb{Z}_{+} \cup\{+\infty\}$-valued measures on $\mathbb{R}^{d}$, In particular, $\Gamma \subset \ddot{\Gamma}$. Analogously to the case of $\Gamma$, we define the vague topology on $\ddot{\Gamma}$ and the corresponding Borel $\sigma$-algebra $\mathcal{B}(\ddot{\Gamma})$.

Theorem 4.3. Let $\varkappa(x, y)$ be of the form $\varkappa(x-y)$ for some $\varkappa \in L^{2}\left(\mathbb{R}^{d}, d x\right)$. Let $J(D)<\infty$. Let $l \in \mathbb{N}$ and $c>0$. Then
(i) The bilinear form $\left(\mathcal{E}_{0}, \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)\right)$ defined by (4.4) is closable on $L^{2}\left(\Gamma, \mu^{(l)}\right)$ and its closure will be denoted by $\left(\mathcal{E}_{0}, D\left(\mathcal{E}_{0}\right)\right)$.
(ii) There exists a conservative diffusion process

$$
M^{0}=\left(\Omega^{0}, \mathcal{F}^{0},\left(\mathcal{F}_{t}^{0}\right)_{t \geq 0},\left(\Theta_{t}^{0}\right)_{t \geq 0},\left(X^{0}(t)\right)_{t \geq 0},\left(P_{\gamma}^{0}\right)_{\gamma \in \ddot{\Gamma}}\right)
$$

on $\ddot{\Gamma}$ which is properly associated with $\left(\mathcal{E}_{0}, D\left(\mathcal{E}_{0}\right)\right)$. In particular, $M^{0}$ is $\mu^{(l)}$-symmetric and has $\mu^{(l)}$ as invariant measure. In the case $d \geq 2$, the set $\ddot{\Gamma} \backslash \Gamma$ is $\mathcal{E}^{0}$-exceptional, so that $\ddot{\Gamma}$ may be replaced by with $\Gamma$ in the above statement.

Proof. We again discuss only the case $l=1$, omitting the upper index (1). By (4.4), for any $F, G \in \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)$,

$$
\begin{align*}
\mathcal{E}_{0}(F, G) & =c \int_{\Omega \times \Gamma} \tilde{\mathbb{P}}(d \omega, d \gamma) \int_{\mathbb{R}^{d}} d x \tilde{\mathbb{E}}\left(Y(x, \omega)^{2} \mid \mathcal{F}\right)\left\langle\nabla_{x} F(\gamma \cup x), \nabla_{x} G(\gamma \cup x)\right\rangle \\
& =\int_{\Omega \times \Gamma} \tilde{\mathbb{P}}(d \omega, d \gamma) \int_{\mathbb{R}^{d}} d x Y(x, \omega)^{2}  \tag{4.14}\\
& \times\left\langle\nabla_{x}(F(\gamma \cup x)-F(\gamma)), \nabla_{x}(G(\gamma \cup x)-G(\gamma))\right\rangle .
\end{align*}
$$

Fix $(\omega, \gamma) \in \Omega \times \Gamma$. Denote

$$
f(x):=F(\gamma \cup x)-F(\gamma), \quad g(x):=G(\gamma \cup x)-G(\gamma)
$$

Clearly, $f, g \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$. In view of Theorem 4.2, $Y(x, \omega)^{2}$ is a continuous function of $x \in \mathbb{R}^{d}$. Hence, by [6, Theorem 6.2], the bilinear form

$$
\mathcal{E}(f, g):=\int_{\mathbb{R}^{d}}\langle\nabla f(x), \nabla g(x)\rangle Y(x, \omega)^{2} d x, \quad f, g \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)
$$

is closable on $L^{2}\left(\mathbb{R}^{d},|Y(x, \omega)|^{2} d x\right)$. Now the closability of $\left(\mathcal{E}_{0}, \mathcal{F} C_{\mathrm{b}}^{\infty}\left(C_{0}^{\infty}\left(\mathbb{R}^{d}\right), \Gamma\right)\right)$ on $L^{2}\left(\Gamma, \mu^{(l)}\right)$ follows by a straightforward generalization of the proof of [6, Theorem 6.3]. Part (ii) of the theorem can be shown completely analogously to [25, 29], see also [20].

Remark 4.3. Heuristically, the generator of $\left(\mathcal{E}_{0}, D\left(\mathcal{E}_{0}\right)\right)$ has the form

$$
(L F)(\gamma)=\sum_{x \in \gamma}\left(\Delta_{x} F(\gamma)+\left\langle\frac{\nabla_{x} r(x, \gamma \backslash x)}{r(x, \gamma \backslash x)}, \nabla_{x} F(\gamma)\right\rangle\right)
$$

Here, for $x \in \gamma, \nabla_{x} F(\gamma):=\left.\nabla_{y} F(\gamma \backslash x \cup y)\right|_{y=x}$ and analogously $\Delta_{x}$ is defined. However, we should not expect that $r(x, \gamma)$ is differentiable in $x$.

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