ON THE GREEN-KUBO FORMULA AND THE GRADIENT CONDITION ON CURRENTS

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In the diffusive hydrodynamic limit for a symmetric interacting particle system (such as the exclusion process, the zero range process, the stochastic Ginzburg-Landau model, the energy exchange model), a possibly non-linear diffusion equation is derived as the hydrodynamic equation. The bulk diffusion coefficient of the limiting equation is given by Green-Kubo formula and it can be characterized by a variational formula. In the case the system satisfies the gradient condition, the variational problem is explicitly solved and the diffusion coefficient is given from the Green-Kubo formula through a static average only. In other words, the contribution of the dynamical part of Green-Kubo formula is 0. In this paper, we consider the converse, namely if the contribution of the dynamical part of Green-Kubo formula is 0, does it imply the system satisfies the gradient condition or not. We show that if the equilibrium measure $\mu$ is product and $L^2$ space of its single site marginal is separable, then the converse also holds. The result gives a new physical interpretation of the gradient condition.

As an application of the result, we consider a class of stochastic models for energy transport studied by Gaspard and Gilbert in [1, 2], where the exact problem is discussed for this specific model.

1. Introduction. In the study of the hydrodynamic limit for a large scale of interacting particle systems, the system is said to satisfy the gradient condition, if the current of the conserved quantity is given by a linear sum of the difference of a local function and its space-shift. If the system satisfies the gradient condition, the diffusion coefficient of the hydrodynamic equation has an explicit expression and the proof of the scaling limit becomes much simpler than the general case (cf. [4]). The underlying structure for this simplification is that the gradient condition implies that the contribution of the dynamical part of Green-Kubo formula is 0. Then, it might be natural to ask whether the converse statement holds or not. Namely, if the contribution of the dynamical part of Green-Kubo formula is 0, does it imply the system...
satisfies the gradient condition? Though the question sounds very natural, we could not find any explicit answer in the literature. In this paper, we give the answer under the assumption that the equilibrium measure is a product measure. The result reveals that the gradient condition is not just a technical condition but has some physical interpretation. Our motivation originally comes from the series of papers by Gaspard and Gilbert [1, 2] where the relation of the gradient condition and the contribution of the dynamic part of GK-formula was discussed. In Section 4, we show an application of our result to this model.

The proof of our main result relies on very fundamental observations for non-dynamical problems. More precisely, the key theorem (Theorem 1 below) concerns only about the properties of the equilibrium measure.

In the next section, we give our general setting and state the main result. In Section 3, we give a proof of Theorem 1. For simplicity we first discuss about the one-dimensional case and then generalize it to the higher dimensional case. In Section 4, we explain an application to the model studied by Gaspard and Gilbert in [1, 2]. In the last section, we discuss on the extension of our result to general Gibbs measures which are not necessarily product.

2. Setting and main result. We consider a general interacting particle system with stochastic dynamics, whose state space is given by a product space \( \Omega = X^{Z^d} \) where \( X \), the single component space, is a measurable space. We suppose that \( \Omega \) is the product measurable space equipped with a translation invariant probability measure \( \mu \) and denote the expectation with respect to \( \mu \) by \( \langle \cdot \rangle \) and the inner product of \( L^2(\mu) \) by \( \langle \cdot, \cdot \rangle \). We denote by \( (\eta_x)_{x \in Z^d} \) the element of \( \Omega \). A measurable function \( f : \Omega \to \mathbb{R} \) is called local if it depends only on a finite number of coordinates, and for a local function \( f \), we define \( s_f := \min\{n \geq 0; f \text{ does not depend on } (\eta_x)_{|x| \geq n+1} \} \) where \( |x| = \max\{|x_1|, |x_2|, \ldots, |x_d|\} \) for \( x \in Z^d \). Shift operators \( \tau_z \) are defined for each \( z \in Z^d \) as \( (\tau_z \eta)_x = \eta_{x-z} \) and \( (\tau_z f)(\eta) = f(\tau_z \eta) \).

Let \( D := \{f \in L^2(\mu); f : \text{local}\} \). If an operator \( T : D \to D \) satisfies that there exists \( r \geq 0 \) such that \( s_{Tf} \leq \max\{s_f, r\} \), then we call it a local operator.

We consider a set of local operators \( (L_{x,y})_{x,y \in Z^d} \) satisfies \( L_{x,y} 1 = 0 \) and the following conditions with convention \( L_{x,x} \equiv 0 \):

- Translation invariance : \( L_{x,y} = \tau_x L_{0,y-x} \tau_{-x} \)
- Finite range : There exists \( R > 0 \) such that \( L_{0,z} \equiv 0 \) for \( |z| > R \)
- Symmetry : \( L_{0,z} = L_{z,0} \) for any \( z \in Z^d \)
- Reversibility : \( \langle L_{0,x} f, g \rangle = \langle f, L_{0,x} g \rangle \) for \( f, g \in D \)
- Non-positivity : \( D_{0,x}(f) := \langle -L_{0,x} f, f \rangle \geq 0 \) for \( f \in D \)

We suppose that \( L = \sum_{x,y \in Z^d} L_{x,y} \) defines the Markov process \( \{\eta_x(t)\}_{x \in Z^d} \).
whose (formal) generator is $L$ with initial distribution $\mu$. We do not attempt here at a justification of this setting in full generality but rather refer to the examples for full rigor. By the reversibility, $\mu$ is the stationary measure for the process.

Our interest is in the case where the conservation quantity exists. Actually, we also suppose that there exists a measurable function $\xi: X \to \mathbb{R}$ such that $\xi(\eta_0) \in L^2(\mu)$ and $L_{x,y}(\xi_x + \xi_y) = 0$ and $L_{x,y}\xi_z = 0$ for $z \neq x, y$ where $\xi_x := \xi(\eta_x)$.

**Example 2.1** The exclusion process with a proper jump rate $c$ is in our setting with $X = \{0, 1\}$, $L_{x,y}f = \frac{1}{2}c(x, y, \eta)(f(\eta^{x-y}) - f(\eta))$, $\mu$: a product Bernoulli measure and $\xi(\eta) = \eta$ where $c(x, y, \eta) = c(y, x, \eta)$ and $\eta^{x-y}$ is the configuration obtained from $\eta$ by exchanging the occupation variables at $x$ and $y$.

**Example 2.2** The generalized exclusion process is in our setting with $X = \{0, 1, \ldots, \kappa\}$,

$L_{x,y}f$

$$= 1_{\{x-y=1\}} \frac{1}{2} (1_{\{\eta_x \geq 1, \eta_y \leq \kappa - 1\}} (f(\eta^{x-y}) - f(\eta)) + 1_{\{\eta_y \geq 1, \eta_x \leq \kappa - 1\}} (f(\eta^{y-x}) - f(\eta))),$$

$\mu$: a translation invariant Gibbs measure and $\xi(\eta) = \eta$ where $\eta^{x-y}$ is the configuration obtained from $\eta$ by letting a particle jump from $x$ to $y$ (cf. [4]).

**Example 2.3** The zero-range process with a proper jump rate $g$ is in our setting with $X = \{0, 1, 2, \ldots, \} = \mathbb{Z}_{\geq 0}$, $L_{x,y}f = 1_{\{x-y=1\}} \frac{1}{2} (g(\eta_x)(f(\eta^{x-y}) - f(\eta)) + g(\eta_y)(f(\eta^{y-x}) - f(\eta))),$ $\mu$: a product Gibbs measure and $\xi(\eta) = \eta$ (cf. [4]).

**Example 2.4** The stochastic Ginzburg-Landau process with proper functions $a$ and $V$ is in our setting with $X = \mathbb{R}$, $L_{x,y} = (\partial_{\eta_x} - \partial_{\eta_y})(a(\eta_x, \eta_y)(\partial_{\eta_x} - \partial_{\eta_y})) + a(\eta_x, \eta_y)(V'(\eta_x) - V'(\eta_y))(\partial_{\eta_x} - \partial_{\eta_y}),$ $\mu$: a product Gibbs measure given by the potential $V$ and $\xi(\eta) = \eta$ (cf. [5, 9]).

**Example 2.5** The stochastic energy exchange model is also in our setting as shown in Section 4.

Let $\langle \xi_0 \rangle = \rho$, $S(x, t) := \mathbb{E}[\xi_x(t)\xi_0(0)] - \rho^2$ and $\chi := \sum_{x \in \mathbb{Z}^d} S(x, t)$, which we suppose finite. As studied in [8] (Section 2.2 of Part II) for exclusion processes, the bulk diffusion coefficient matrix $D = (D_{\alpha\beta})$ for the conserved quantity $\xi$ is defined as

$$D_{\alpha\beta} := \lim_{t \to \infty} \frac{1}{\ell} \frac{1}{2\chi} \sum_{x \in \mathbb{Z}^d} x_{\alpha} x_{\beta} S(x, t),$$
\[ D_{\alpha\beta} = \frac{1}{\chi} \left( 2 \sum_x x_\alpha x_\beta D_{0,x}(\xi_0) - 2 \int_0^\infty \sum_x E[j_\alpha e^{Lt} \tau_x j_\beta] dt \right) \]

where \( j_\alpha = \sum_{x \in \mathbb{Z}^d} x_\alpha j_{0,x} \) and \( j_{0,x} = 2L_{0,x} \xi_0 = (L_{0,x} + L_{x,0}) \xi_0 \), which is a current between 0 and \( x \).}

Note that \( \xi_x(t) - \int_0^t \sum_z j_{x,x+s}(s) ds \) is a Martingale and \( 2D_{0,x}(\xi_0) = -\langle \xi_0, j_{0,x} \rangle \).

We define the matrix \( D^s \) as

\[ D^s_{\alpha\beta} = \frac{1}{\chi} \sum_x x_\alpha x_\beta D_{0,x}(\xi_0). \]

We introduce the Hilbert space \( \mathcal{H} \) of functions on \( \Omega \) as the completion of \( \mathcal{D} \) equipped with the (degenerate) scalar product

\[ \langle f | g \rangle := \sum_{x \in \mathbb{Z}^d} (\langle \tau_x f g \rangle - \langle f \rangle \langle g \rangle). \]

Here, we suppose that the measure \( \mu \) satisfies an enough spatial mixing condition to make \( \langle f | g \rangle \) be well-defined for any \( f, g \in \mathcal{D} \). Actually, in our main theorem, we only consider product measures. We also suppose that \( e^{Lt} \) induces the self-adjoint semigroup \( T_t \) on \( \mathcal{H} \) and denote its generator by \( \tilde{L} \).

The following variational formula also holds under a general condition (cf. [8] (Section 2.2 of Part II), [4]):

\[ \sum_{\alpha, \beta=1}^d \ell_{\alpha\beta} D_{\alpha\beta} = \sum_{\alpha, \beta=1}^d \ell_{\alpha\beta} D^s_{\alpha\beta} + \inf_{f \in \mathcal{D}} \left\{ -2(\sum_{\alpha=1}^d \ell_{\alpha} j_\alpha |f|) - \langle f | \tilde{L} f \rangle \right\} \]

for all \( \ell = (\ell_\alpha) \in \mathbb{R}^d \).

So far, we did not prove anything and just introduce the settings. From now on, under the assumption that relations (1) and (2) hold, we state our main result. For this, we introduce the gradient space

\[ \mathcal{G} := \{ \sum_{\alpha=1}^d (\tau^\alpha g_\alpha - g_\alpha); g_\alpha \in \mathcal{D}, \alpha = 1, 2, \ldots, d \} \]
where $\tau^\alpha = \tau e_\alpha$ and $e_\alpha$ is the unit vector to the $\alpha$-th direction. The stochastic system defined by $L$ is said to satisfy the gradient condition, if $j_\alpha \in \mathfrak{G}$ for all $\alpha = 1, 2, \ldots, d$.

Our main result is that our stochastic system satisfies the gradient condition if and only if $D = D^s$ under the condition that $\mu$ is product and the $L^2$ space of its single site marginal is separable. To show this, we first give two simple lemmas.

**Lemma 2.1** If the stochastic system defined by $L$ satisfies the gradient condition, then $D = D^s$. Namely, the variational formula (2) attains its minimum with $f = 0$.

**Proof.** If $j_\alpha \in \mathfrak{G}$ for all $\alpha = 1, 2, \ldots, d$, then $\langle \sum_{\alpha=1}^d \ell_\alpha j_\alpha | f \rangle = 0$. Since $-\langle f | \bar{L} f \rangle \geq 0$ for any $f \in D$ by the positivity condition, $\sum_{\alpha, \beta=1}^d \ell_\alpha \ell_\beta D_{\alpha \beta} = \sum_{\alpha, \beta=1}^d \ell_\alpha \ell_\beta D^s_{\alpha \beta}$ for all $\ell = (\ell_\alpha) \in \mathbb{R}^d$.

**Lemma 2.2** If $D = D^s$ holds, then $\langle j_\alpha | j_\alpha \rangle = 0$ for all $\alpha = 1, 2, \ldots, d$.

**Proof.** If $D = D^s$, we obtain from (2) with $\ell_\alpha = 1, \ell_\beta = 0$ for $\beta \neq \alpha$, that $\inf_{f \in D} \{-2\langle j_\alpha | f \rangle - \langle f | \bar{L} f \rangle\} = 0$. Since $j_\alpha = 2 \sum_x x_\alpha \xi_0 \in D$, we can take $cj_\alpha$ as $f$ in the above variational formula for any $c \in \mathbb{R}$ and obtain $\inf_{c \in \mathbb{R}} \{-2c\langle j_\alpha | j_\alpha \rangle - c^2 \langle j_\alpha | \bar{L} j_\alpha \rangle\} \geq 0$ which implies $\langle j_\alpha | j_\alpha \rangle = 0$.

**Remark 2.1** If the interaction of our system is nearest-neighbour, namely, $R = 1$, then we have $j_\alpha = j_0, e_\alpha - j_0, -e_\alpha = j_0, e_\alpha + \tau e_\alpha j_0, e_\alpha$ since $j_0, -e_\alpha = 2L_0, -e_\alpha \xi_0 = -2L_0, -e_\alpha \xi_0 = -2L_0, -e_\alpha \xi_0 = -\tau e_\alpha j_0, e_\alpha$. For this case, $\langle j_\alpha | j_\alpha \rangle = 0$ is equivalent to $\langle j_0, e_\alpha | j_0, e_\alpha \rangle = 0$. Next theorem is the most essential result and we give its proof in the next section.

**Theorem 1** Assume that $\mu$ is product with a single site marginal $\nu$, namely $\mu = \nu^\otimes d$, and $L^2(\nu)$ is separable. Then, if $f \in D$ satisfies $\langle f | f \rangle = 0$, then $f \in \mathfrak{G}$. Equivalently, the intersection of the kernel of $\langle \cdot | \cdot \rangle$ and $D$ is the direct sum of the space of constant functions and $\mathfrak{G}$.

Combining this theorem with the above lemmas, we obtain our main result as a straightforward corollary.
Corollary 2.1 Assume that μ is product with a single site marginal ν, namely μ = ν^d, and L^2(ν) is separable. Then, the stochastic system defined by L satisfies the gradient condition, if and only if D = D^s. Moreover, if R = 1, then D = D^s if and only if j_α = g for all α = 1, 2, ..., d.

Remark 2.2 Separability condition for L^2(ν) is quite mild. In particular, if ν is a probability measure on the measurable space (E, Ω) where E is a Borel set of Euclidean space and Ω = Borel sets on E, then L^2(ν) is separable. Also, if X is a countable set, then L^2(ν) is separable.

Remark 2.3 Theorem 1 is nothing to do with the dynamics, but only concerns the probability measure μ.

3. Proof of Theorem 1. In the first subsection, we give a proof for the case d = 1. In the second subsection, we generalize it to the case d ≥ 2.

3.1. One dimensional setting. We consider the case d = 1. Let (X, F, ν) be a probability space where L^2(ν) is separable and Ω := X^k be the infinite product probability space equipped with the probability measure μ := ν^d.

Let D_0 := {f ∈ L^2(μ): f : local, ∥f∥ = 0}. For f ∈ D_0, we define a semi-norm ||·|| as

||f||^2 := lim_{k→∞} \frac{1}{2k+1} \left(\sum_{x=0}^{k} \tau_x f(x)\right)^2 = \sum_{x∈Z} \langle fτ_x f\rangle = ∥f∥^2.

Theorem 1 concerns the relation between the gradient space G := {τg - g: g ∈ D} = {τg - g: g ∈ D_0} and the kernel of the semi-norm C_0 := {f ∈ D_0: ||f|| = 0}.

It is easy to see that G ⊂ C_0. Theorem 1 claims that G ⊃ C_0 hence G = C_0.

To prove this, we first start with a simple lemma. Let ℓ^2 := {a = (α_x)_{x∈Z} ∈ R^d: |A| < ∞}. Here |A| represents the number of elements for a set A. For a ∈ ℓ^2 satisfying a ≠ 0, define M_a := max{x ∈ Z: α_x ≠ 0} and m_a := min{x ∈ Z: α_x ≠ 0}. As a convention, take M_0 = m_0 = 0. We also define A_a := {x ∈ Z: m_a ≤ x ≤ M_a} and s_a := |A_a|.

Lemma 3.1 Let f ∈ D_0 and assume that there exists (α_x)_{x∈Z} ∈ ℓ^2 and h ∈ D_0 satisfying f = \sum x α_x τ_x h and \sum x α_x = 0. Then, there exists a unique function g ∈ D_0 such that f = τg - g, hence f ∈ G. Moreover, if {τ_x h}_{x∈Z} are orthogonal in L^2(μ), then ∥g∥ ≤ \frac{s_a^2}{2} \sum x∈Z α_x^2 (h)^2.

Proof. Uniqueness: If f = τg_1 - g_1 = τg_2 - g_2 and g_1, g_2 ∈ D_0, then τ(g_1 - g_2) = g_1 - g_2 and g_1 - g_2 ∈ D_0. In particular, g_1 - g_2 is local and shift invariant, so it must be a constant. Also, {g_1 - g_2} = 0, hence g_1 = g_2.
Proof. particular, \( \phi \) Lemma 3.2 For any \((\text{pair}) \) index space \( \Theta \) from the assumption that \( \phi \) satisfying 5.2 of [6].

Remark 3.1 The result and the proof of Lemma 3.1 is similar to Lemma 5.2 of [6].

Now, we consider a generalized Fourier series in the space \( L^2(\mu) \). Let \( \mathbb{N}_0 := \{0, 1, 2, \ldots \} \) and \( \{\phi_n\}_{n \in \mathbb{N}_0} \) be a countable orthonormal basis of \( L^2(\nu) \) satisfying \( \phi_0 = 1 \). The existence of the countable orthonormal basis follows from the assumption that \( L^2(\nu) \) is separable. Let us introduce the multi-index space \( \Theta := \{\mathbf{n} = (n_x)_{x \in \mathbb{Z}} \in \mathbb{N}_0^\mathbb{Z} : |\{x \in \mathbb{Z} : n_x \neq 0\}| < \infty \} \). Then, the set of functions \( \{\phi_n\}_{n \in \Theta} \) is the countable orthonormal basis of \( L^2(\mu) \) where \( \phi_n(\eta) := \Pi_{x \in \mathbb{Z}} \phi_{n_x}(\eta_x) \). In particular, if \( f \in L^2(\mu) \), then \( f = \sum_{n \in \Theta} \hat{f}_n \phi_n \) with \( \hat{f}_n = \langle f \phi_n \rangle \).

We define the shift operator \( (\tau_x \mathbf{n})_x = n_{x-z} \) and \( \Theta_s := \{\mathbf{n} = (n_x)_{x \in \mathbb{Z}} \in \Theta : n_x = 0 \ (\forall x < 0), n_0 \neq 0 \} \). Then, for any \( \mathbf{n} \in \Theta \setminus \{0\} \), there exists a unique pair \( (x, \mathbf{n}_s) \in \mathbb{Z} \times \Theta_s \) such that \( \mathbf{n} = \tau_x \mathbf{n}_s \).

The next lemma is about the locality of the Fourier series.

Lemma 3.2 For any \( f \in \mathcal{D}_0 \) and \( \mathbf{n}_s \in \Theta_s \), \( \tilde{f}_{\tau_x \mathbf{n}_s} = 0 \) if \( |x| \geq s_f + 1 \). In particular, \( \tilde{f}_{\tau_x \mathbf{n}_s} \in \ell_2^c \).

Moreover, for \( \mathbf{n}_s \in \Theta_s \) satisfying \( n_{s+y} \neq 0 \) with some \( |y| \geq 2s_f + 1 \), \( \tilde{f}_{\tau_x \mathbf{n}_s} = 0 \) for all \( x \in \mathbb{Z} \).

Proof. For \( |x| \geq s_f + 1 \), \( \phi_{n_{x,0}}(\eta_x) \) and \( f \) are independent and \( \langle \phi_{n_{x,0}}(\eta_x) \rangle = 0 \), so \( \langle \phi_{n_{x,0}}(\eta_x) \Pi_{y \in \mathbb{Z} \setminus \{0\}} \phi_{n_{x,y}}(\eta_{x+y}) f \rangle = 0 \).
Similarly, if \( n_* \in \Theta_* \) satisfying \( n_* y \neq 0 \) with some \( |y| \geq 2s_f + 1 \), then \( \phi_{n_* y}(\eta_x) \) and \( f \) are independent for \( |x| \leq -s_f - 1 \) and \( \phi_{n_* y}(\eta_x + y) \) and \( f \) are independent for \( |x| \geq -s_f \) so we have \( \tilde{f}_{\tau_n} = 0 \) for both cases.

The next lemma is simple but one of the keys of our main result.

**Lemma 3.3** For \( f \in \mathcal{D}_0 \), \( |f|^2 = \sum_{n_* \in \Theta_*} (\sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_*)^2 \).

**Proof.** Since \( \langle f \rangle = 0 \), \( f_0 = 0 \). Then, by the general observation, \( f = \sum_{n \in \Theta} \tilde{f}_n \phi_n = \sum_{n_* \in \Theta_*} \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \phi_{\tau_n} n_* \). Then,

\[
|f|^2 = \sum_{x \in \mathbb{Z}} (f_{\tau_n} f) = \sum_{x \in \mathbb{Z}} \left( \sum_{n_* \in \Theta_*} \sum_{x' \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \phi_{\tau_n} n_* \right) \left( \sum_{n'_* \in \Theta_*} \sum_{x' \in \mathbb{Z}} \tilde{f}_{\tau_n} n'_* \phi_{\tau_n} n'_* \right).
\]

Since \( \sum_{n_* \in \Theta_*} \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* < \infty \) and \( \{\phi_n\}_{n \in \Theta} \) is an orthonormal basis, we have

\[
\langle \left( \sum_{n_* \in \Theta_*} \tilde{f}_{\tau_n} n_* \phi_{\tau_n} n_* \right) \left( \sum_{n'_* \in \Theta_*} \sum_{x' \in \mathbb{Z}} \tilde{f}_{\tau_n} n'_* \phi_{\tau_n} n'_* \right) \rangle = \sum_{n_* \in \Theta_*} \sum_{x' \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \phi_{\tau_n} n_* \phi_{\tau_n} n'_* \phi_{\tau_n} n'_*.
\]

Therefore,

\[
|f|^2 = \sum_{x \in \mathbb{Z}} \sum_{n_* \in \Theta_*} \sum_{x' \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \tilde{f}_{\tau_n} n'_* = \sum_{n_* \in \Theta_*} \left( \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \right)^2.
\]

**Proposition 3.1** If \( f \in \mathcal{C}_0 \), then \( f \in \mathfrak{E} \).

**Proof.** By Lemma 3.3, \( |f| = 0 \) implies \( \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* = 0 \) for any \( n_* \in \Theta_* \). Then, combining the fact that \( \phi_n \in \mathcal{D}_0 \) for each \( n \in \Theta_* \) with Lemma 3.1, for each fixed \( n_* \in \Theta_* \), there exists \( g_{n_*} \in \mathcal{D}_0 \) such that \( \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} \phi_{\tau_n} n_* = \tau g_{n_*} - g_{n_*} \). Moreover, since \( \{\phi_{\tau_n} n_*\}_{x \in \mathbb{Z}} \) are orthogonal and \( \tilde{f}_{\tau_n} n_* = 0 \) for \( |x| \geq s_f + 1 \) by Lemma 3.2,

\[
\langle g_{n_*} \rangle \leq (2s_f + 1)^2 \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \left( \phi_{n_*} \right)^2 = (2s_f + 1)^2 \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_*^2.
\]

By the construction, \( \{g_{n_*}\}_{n_*} \) are orthogonal in \( L^2(\mu) \) and so \( g := \sum_{n_* \in \Theta_*} g_{n_*} \in L^2(\mu) \) since

\[
\langle g \rangle = \sum_{n_* \in \Theta_*} \langle g_{n_*} \rangle \leq (2s_f + 1)^2 \sum_{n_* \in \Theta_*} \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_*^2 = (2s_f + 1)^2 \langle f \rangle.
\]

Also, \( \langle g \rangle = 0 \). The locality of \( g \) follows from the following two facts: (i) \( g_{n_*} = 0 \) if \( n_* \in \Theta_* \) satisfying \( n_* y \neq 0 \) with some \( |y| \geq 2s_f + 1 \) by Lemma 3.2, (ii) the support of \( g_{n_*} \) is included in the union of the support of \( \{\tau \phi_{\tau_n} n_*\}_{-s_f \leq x \leq s_f} \). Therefore, \( g \in \mathcal{D}_0 \).

Finally, we see that

\[
f = \sum_{n_* \in \Theta_*} \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \phi_{\tau_n} n_* = \sum_{n_* \in \Theta_*} \sum_{x \in \mathbb{Z}} \tilde{f}_{\tau_n} n_* \tau \phi_{\tau_n} n_* = \sum_{n_* \in \Theta_*} (\tau g_{n_*} - g_{n_*}) = \tau g - g
\]

which implies \( f \in \mathfrak{E} \).
3.2. Multi-dimensional setting. In this subsection, we generalize our result to the multi-dimensional setting.

Let \((X, \mathcal{F}, \nu)\) be an probability space where \(L^2(\nu)\) is separable and \(\Omega := X^{\mathbb{Z}^d}\) be the infinite product probability space equipped with the probability measure \(\mu := \nu^{\mathbb{Z}^d}\).

Let \(D_0 := \{ f \in L^2(\mu); f \text{ local}, \langle f \rangle = 0 \}. \) For \( f \in D_0 \), we define a semi-norm \(\| \cdot \|\) as

\[
\| f \|^2 := \lim_{k \to \infty} \frac{1}{(2k + 1)^d} \left\langle \left( \sum_{|x| \leq k} \tau_x f \right)^2 \right\rangle = \sum_{x \in \mathbb{Z}^d} \langle f \tau_x f \rangle = \langle f | f \rangle.
\]

Recall that

\[
\mathcal{G} := \left\{ \sum_{\alpha = 1}^d (\tau^\alpha g_\alpha - g_\alpha); g_\alpha \in D, \alpha = 1, 2, \ldots, d \right\}
\]

\[
= \left\{ \sum_{\alpha = 1}^d (\tau^\alpha g_\alpha - g_\alpha); g_\alpha \in D_0, \alpha = 1, 2, \ldots, d \right\}
\]

and \(C_0 := \{ f \in D_0; \| f \| = 0 \}. \)

Let \(\ell^2_c := \{ a = (a_x)_{x \in \mathbb{Z}^d} \in \mathbb{R}^{\mathbb{Z}^d}; |\{ x \in \mathbb{Z}^d; a_x \neq 0 \}| < \infty \}. \) For \( a \in \ell^2_c \) satisfying \( a \neq 0 \), define \( M_a = \max_{1 \leq \alpha \leq d} \max \{ x \in \mathbb{Z}; \exists y \neq 0 \text{ s.t. } a_y = x \} \) and \( m_a := \min_{1 \leq \alpha \leq d} \min \{ x \in \mathbb{Z}; \exists y \neq 0 \text{ s.t. } a_y = x \}. \) As a convention, take \( M_0 = m_0 = 0. \)

We also define \( \Lambda_a := \{ x \in \mathbb{Z}^d; m_a \leq x \leq M_a, \alpha = 1, 2, \ldots, d \} \) and \( s_a := |\Lambda_a| \). Here, the only essential property is that the hypercube \( \Lambda_a \) satisfies \( \Lambda_a \supset \{ x \in \mathbb{Z}^d; a_x \neq 0 \} \).

**Lemma 3.4** Assume that \((a_x)_{x \in \mathbb{Z}^d} \in \ell^2_c \) and \(\sum_{x \in \mathbb{Z}^d} a_x = 0. \) Then, there exists a \(d\)-tuple of functions \((b^1, b^2, \ldots, b^d) \in (\ell^2_c)^d\) such that \( a_x = \sum_{\alpha = 1}^d (b^x_\alpha - b_x^\alpha). \)

In particular, we can choose the \(d\)-tuple to satisfy that the support of \(b^\alpha\) are included in \(\Lambda_a\) and \(\sum_{\alpha = 1}^d \sum_{x \in \mathbb{Z}^d} (b^x_\alpha)^2 \leq s_a^2 \sum_{x \in \mathbb{Z}^d} a_x^2. \)

**Proof.** It is a classical discrete problem. Consider \(\Lambda_a\) as a finite graph, and let \(\Delta\) be the graph Laplacian of \(\Lambda_a\). Then, since the graph is connected and \(\sum_{x \in \Lambda_a} a_x = 0, \) there exists a solution \((q_x)_{x \in \Lambda_a}\) of the Poisson equation \(\Delta q = a. \)

By the definition of graph Laplacian, for any \(x \in \Lambda_a, \) we have

\[
a_x = \sum_{\alpha = 1}^d \sum_{x, \mathcal{E} \in \Lambda_a} (q_{x + e_\alpha} - q_x) + \sum_{x - e_\alpha \in \Lambda_a} (q_{x - e_\alpha} - q_x).
\]

Let \(\tilde{b}_{x,y} = q_y - q_x \) for \(x, y \in \Lambda_a\) and \(\tilde{b}_{x,y} = 0\) otherwise. Then,

\[
a_x = \sum_{i = 1}^d (\tilde{b}_{x,x+e_\alpha} - \tilde{b}_{x-e_\alpha, x})
\]
for any $x \in \mathbb{Z}^d$. Taking $b^0_x = -\tilde{b}_{x,x+e_\alpha}$, we have

$$a_x = \sum_{\alpha=1}^{d} (b^0_{x,e_\alpha} - b^2_x)$$

for any $x \in \mathbb{Z}^d$. The last estimate

$$\sum_{\alpha=1}^{d} \sum_{x \in \mathbb{Z}^d} (b^0_x)^2 = \sum_{\alpha=1}^{d} \sum_{x, x+e_\alpha \in \Lambda_a} (q_{x+e_\alpha} - q_x)^2$$

$$= \sum_{x \in \Lambda_a} q_x \Delta q_x \leq s^2_a \sum_{x \in \mathbb{Z}^d} (\Delta q_x)^2 = s^2_a \sum_{x \in \mathbb{Z}^d} a_x^2$$

follows from the simple spectral gap estimate of $\Delta$.

**Lemma 3.5** Let $f \in \mathcal{D}_0$ and assume that there exists $(a_x)_{x \in \mathbb{Z}^d} \in \ell^2_c$ and $h \in \mathcal{D}_0$ satisfying $f = \sum_{x \in \mathbb{Z}^d} a_x \tau_x h$ and $\sum_{x \in \mathbb{Z}^d} a_x = 0$. Then, there exists a $d$-tuple of functions $(g_1, g_2, \ldots, g_d) \in (\mathcal{D}_0)^d$ such that $f = \sum_{\alpha=1}^{d} (\tau^\alpha g_\alpha - g_\alpha)$, hence $f \in \Theta$. Moreover, if $\{\tau_x h\}_{x \in \mathbb{Z}^d}$ are orthogonal in $L^2(\mu)$, then we can find such a $d$-tuple of functions $(g_1, g_2, \ldots, g_d)$ which satisfy also $\sum_{\alpha=1}^{d} (g_\alpha^2) \leq s^2_a \sum_{x \in \mathbb{Z}^d} a_x^2(h^2)$.

**Proof.** From Lemma 3.4, we have a $d$-tuple of functions $(b^1, b^2, \ldots, b^d) \in (\ell^2_c)^d$ such that $a_x = \sum_{\alpha=1}^{d} (b^\alpha_{x-e_\alpha} - b^\alpha_x)$ whose support is in $\Lambda_a$. Define $g_\alpha = \sum_{x \in \mathbb{Z}^d} b^\alpha_x \tau_x h$. We show it is the desired set of functions.

First, we have $g_\alpha \in \mathcal{D}_0$ since $b^\alpha \in \ell^2_c$. Also,

$$f = \sum_{x} a_x \tau_x h = \sum_{x} \sum_{\alpha=1}^{d} (b^\alpha_x - b^\alpha_{x-e_\alpha}) \tau_x h = \sum_{x} \sum_{\alpha=1}^{d} b^\alpha_x (\tau_x h - \tau_{x+e_\alpha} h) = \sum_{\alpha=1}^{d} (\tau^\alpha g_\alpha - g_\alpha).$$

Finally, if $\{\tau_x h\}_{x \in \mathbb{Z}^d}$ are orthogonal in $L^2(\mu)$, then

$$\sum_{\alpha=1}^{d} (g_\alpha^2) = (h^2) \sum_{\alpha=1}^{d} \sum_{x \in \mathbb{Z}^d} (b^\alpha_x)^2$$

and the last estimate also follows from Lemma 3.4.

**Remark 3.2** The uniqueness result in Lemma 3.1 fails for the multi-dimensional case since $\Lambda_a$ may have a non trivial cycle. Note that we do not use the uniqueness result anywhere in the proofs.

For the part of the generalize Fourier series, we do not need to change the strategy. Note that we define $\Theta_s$ as the quotient of $\Theta \setminus \{0\}$ by the equivalence relation $n \sim n'$ if any only if there exists $x \in \mathbb{Z}^d$ such that $\tau_x n = n'$.

To make clear the locality of the Fourier series, we introduce the following notation.

For $n_s \in \Theta_s$, let $\text{rad}(n_s) = \max_{1 \leq \alpha \leq d} \max \{|x_\alpha - x'_\alpha|, n_s x \neq 0, n_s x' \neq 0\}$. 
Lemma 3.6 For $n^{*} \in \Theta_{s}$ satisfying $\text{rad}(n^{*}) \geq 2s_{f} + 1$, $\tilde{f}_{\tau_{x}n^{*}} = 0$ for all $x \in \mathbb{Z}^{d}$.

Moreover, if $\text{rad}(n^{*}) \leq 2s_{f}$, then we can choose the representative $n^{*}$ so as $\{x \in \mathbb{Z}^{d}; n^{*}x \neq 0\} \subset \{x \in \mathbb{Z}^{d}; -s_{f} \leq x_{\alpha} \leq s_{f}, \alpha = 1, 2, \ldots, d\}$. Then, for this representative, $\tilde{f}_{\tau_{x}n^{*}} = 0$ if $|x| \geq s_{f} + 1$.

The next lemma holds in the same way as the one-dimensional case.

Lemma 3.7 For $f \in D_{0}, \|f\| = \sum_{n^{*} \in \Theta_{s}} (\sum_{x \in \mathbb{Z}^{d}} \tilde{f}_{x\tau_{x}n^{*}})^{2}$.

Our main result also follows in the same way. Just note that $\sum_{x \in \mathbb{Z}^{d}} \tilde{f}_{x\tau_{x}n^{*}} \tau_{x} \phi_{n^{*}}$ does not depend on the choice of the representative of $n^{*}$.

Proposition 3.2 If $f \in C_{0}$, then $f \in \mathcal{G}$.

Hence, we prove $C_{0} \subset \mathcal{G}$ and so Theorem 1.

4. Application to the stochastic energy transport model. In this section, we show an application of our result to one specific model called stochastic energy transport model, which is paid much attention from particularly physical point of view. See more detailed background of the model in [1, 2].

The model is heuristically obtained as a mesoscopic energy transport model from a microscopic mechanical dynamics consist of a one-dimensional array of two-dimensional cells, each containing a single hard-disc particle or an array of three-dimensional cells, each containing a single hard-sphere particle.

This mesoscopic model completely fits to our general setting taking $(X, \mathcal{F}, \mu) = ((0, \infty), \mathcal{B}((0, \infty), \nu))$ where

$$\nu(d\eta) = \frac{\eta^{d-1} \exp(-\frac{\eta}{T})}{T^{d} \Gamma(\frac{d}{2})} d\eta$$

with a given model parameter $d$ and the temperature $T$.

The operator $L$ is the generator of the infinite volume dynamics, given as

$$Lf = \sum_{x \in \mathbb{Z}^{d}} (L_{x,x+1} + L_{x+1,x})f$$

where

$$L_{x,x+1}f = \frac{1}{2} \int_{-\eta_{x+1}}^{\eta_{x+1}} du[W(\eta_{x}-\eta, \eta_{x+1}+u|\eta_{x}, \eta_{x+1}) f(\ldots, \eta_{x}-u, \eta_{x+1}+u, \ldots) - f(\eta)]$$

and $L_{x,x+1} = L_{x+1,x}$ where $W(\eta_{a}, \eta_{b} | \eta_{a} - u, \eta_{b} + u)$ describes the rate of exchange of energy $u$ between sites $a$ and $b$ at respective energies $\eta_{a}$ and $\eta_{b}$.

The specific forms of the kernel should be found in [1, 2]. The dynamics obviously conserves the sum of the energies, hence $\xi(\eta) = \eta$. 

Under the diffusive space-time scaling limit, the time evolution of the local temperature will be given by
\[ \partial_t T = \partial_x (D(T) \partial_x T), \quad T = T(x,t) \]
with thermal diffusivity \( D(T) \). In [1, 2], the authors conjectured that \( D(T) = D^s(T) \) where \( D^s(T) \) is the static part of the thermal diffusivity. However, with our main result, \( D(T) = D^s(T) \) implies the energy current is the gradient and it is not true, hence we conclude that the conjecture fails. Recently, they show in [3] how the variational characterization of the diffusion coefficient given in [7] can be put to use to obtain the correction to static (or instantaneous) part of the diffusion coefficient and carried out further molecular dynamics simulations, which on one side confirm our picture and on the other hand also show that the correction is very small.

5. Discussion on general Gibbs measures. To extend Theorem 1 for general Gibbs measures which are not necessarily product is an important future problem. In this section, we discuss some observations on this topic.

First, we emphasize that Lemma 2.2 holds for general Gibbs measures. Also, we may expect
\[
\lim_{k \to \infty} \frac{1}{(2k+1)^d} \langle (\sum_{|x| \leq k} \tau_x f)^2 \rangle = \sum_{x \in \mathbb{Z}^d} \langle f \tau_x f \rangle = \langle f | f \rangle
\]
for any \( f \in D_0 \), hence \( \langle f | f \rangle \geq 0 \) for general Gibbs measures. Then, the following holds.

**Lemma 5.1** For any \( \alpha = 1, 2, \ldots, d \), \( \langle j_{\alpha} | j_{\alpha} \rangle = 0 \) holds if and only if \( \langle j_{\alpha} | f \rangle = 0 \) for all \( f \in D_0 \).

**Proof.** If \( \langle j_{\alpha} | j_{\alpha} \rangle = 0 \) holds, for any \( f \in D_0 \) and \( c \in \mathbb{R} \),
\[
\langle j_{\alpha} + cf | j_{\alpha} + cf \rangle = 2c \langle j_{\alpha} | f \rangle + c^2 \langle f | f \rangle \geq 0.
\]
Hence, \( \langle j_{\alpha} | f \rangle = 0 \).

From above lemma, for any \( \psi \in D_0 \), \( \sum_x c_{x,\alpha,\psi} = 0 \) where \( c_{x,\alpha,\psi} = \langle j_{\alpha} \tau_x \psi \rangle \).

Now, we summarize two essential properties we used in the proof of Theorem 1, which are satisfied if the measure \( \mu \) is product.

(i) For any \( \psi \in D_0 \) and \( \alpha = 1, 2, \ldots, d \), \( \{c_{x,\alpha,\psi}\}_{x \in \mathbb{Z}^d} \) is in \( l_2 \), namely, the support of \( \{c_{x,\alpha,\psi}\}_{x \in \mathbb{Z}^d} \) is bounded. Therefore, we can apply Lemma 3.1 or 3.5.

(ii) There exists a set of countable functions \( \{\psi_n\}_{n \in \mathbb{N}_0} \) in \( L^2(\mu) \) satisfying \( \psi_0 \equiv 1 \), \( \psi_n \in D_0 \) for \( n \in \mathbb{N} \) and \( \{\tau_x \psi_n\}_{x \in \mathbb{Z}^d, n \in \mathbb{N}} \cup \{\psi_0\} \) forms an orthonormal basis of \( L^2(\mu) \). By definition, \( j_{\alpha} = \sum_{n \in \mathbb{N}} \sum_{x \in \mathbb{Z}^d} c_{x,\alpha,n} \tau_x \psi_n \) where \( c_{x,\alpha,n} = \langle j_{\alpha} \tau_x \psi_n \rangle \). Therefore, combining with (i), we can prove Theorem 1.
Even for general Gibbs measures, if we have these two properties, we can prove Theorem 1, but they do not hold under the typical exponential mixing condition. So far, it is difficult to conjecture whether Theorem 1 holds generally or not.

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