# A STOCHASTIC BURGERS EQUATION FROM A CLASS OF MICROSCOPIC INTERACTIONS 

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#### Abstract

We consider a class of nearest-neighbor weakly asymmetric mass conservative particle systems evolving on $\mathbb{Z}$, which includes zero-range and types of exclusion processes, starting from a perturbation of a stationary state. When the weak asymmetry is of order $O\left(n^{-\gamma}\right)$ for $1 / 2<\gamma \leq 1$, we show that the scaling limit of the fluctuation field, as seen across process characteristics, is a generalized Ornstein-Uhlenbeck process. However, at the critical weak asymmetry when $\gamma=1 / 2$, we show that all limit points satisfy a martingale formulation which may be interpreted in terms of a stochastic Burgers equation derived from taking the gradient of the KPZ equation. The proofs make use of a sharp "Boltzmann-Gibbs" estimate which improves on earlier bounds.


1. Introduction. There has been much recent work on the classification of fluctuations of certain interfaces and currents, corresponding to mass conservative particle dynamics in one-dimensional nearest-neighbor interacting particle systems such as simple exclusion and its variants, with respect to so-called EdwardsWilkinson (EW) and Kardar-Parisi-Zhang (KPZ) classes (cf. [19] for a review and references). Following recent sensibilities, a $d=1$ particle system is in the EW class if the standard deviation of the associated "height" function $h_{t}(x)$ of the interface at time $t$ and space point $x$, or the integrated current at time $t \geq 0$ across the space point $x \in \mathbb{R}$, is of order $t^{1 / 4}$, and also spatial correlations are nontrivial at range $t^{1 / 2}$. Examples in this class are independent random walk systems, random averaging and reversible simple exclusion processes starting from a stationary state or even in nonstationary states [9, 22, 32, 37, 54].

On the other hand, a system is in the KPZ class if its "height" function and integrated current have standard deviation of order $t^{1 / 3}$, and nontrivial spatial correlations at range $t^{2 / 3}$. A well-studied particle system model in this class is the asymmetric simple exclusion process starting from deterministic initial configurations such as step profile and alternating conditions, or from a stationary state (cf. $[6,7,10,11,16,17,24,39,46,49,58]$ and references therein).

[^0]These two classes can be seen in the study of the famous KPZ stochastic partial differential equation first mentioned in [33]:

$$
\begin{equation*}
\partial_{t} h_{t}(x)=D \Delta h_{t}(x)+a\left(\nabla h_{t}(x)\right)^{2}+\sigma \dot{\mathcal{W}}_{t}(x) \tag{1.1}
\end{equation*}
$$

where $\dot{\mathcal{W}}_{t}(x)$ is a space-time white noise with unit variance. When $a=0$ and $D, \sigma>0$, then $h_{t}(x)$ is a generalized Ornstein-Uhlenbeck process in EW class. However, when $a \neq 0$ and $D, \sigma>0$, a physical argument indicates that $h_{t}(x)$ is in the KPZ class (cf. [33,59]). Also, in another sense, it has been shown that the "Cole-Hopf" solution of the KPZ equation, starting from certain initial conditions, interpolates between the two classes when the centered solution is examined in different asymptotic scaling regimes, that is, when normalized by $t^{1 / 3}$ as $t \uparrow \infty$ or when normalized by $t^{1 / 4}$ as $t \downarrow 0$, nontrivial limits are obtained (cf. [1, 8]).

Moreover, it is believed that in many "critical" weakly asymmetric, $d=1$ particle systems, that is, when the weak asymmetry is scaled at a critical level, the diffusively scaled "height" function or integrated current should converge to the solution of the KPZ equation with parameters depending on the structure of particle interactions and initial conditions. Recently, much progress has been made in making clear this convergence. Part of the difficulty is that, since "solutions" to the KPZ equation are expected to be distribution-valued, the nonlinear term in the equation does not make sense, and so the equation is ill-posed. Hence, what does it mean to solve the KPZ equation and also, when properly interpreted, how to derive the KPZ equation from microscopic particle interactions?

One way to approach these questions is to observe that the Cole-Hopf transformation $z_{t}(x):=\exp \left\{(a / D) h_{t}(x)\right\}$ linearizes the KPZ equation to a stochastic heat equation

$$
\begin{equation*}
\partial_{t} z_{t}(x)=D \Delta z_{t}(x)+(a \sigma / D) z_{t}(x) \dot{\mathcal{W}}_{t}(x) \tag{1.2}
\end{equation*}
$$

which can be solved uniquely starting from a class of initial conditions and is also strictly positive for times $t>0[44,61]$. Then the "Cole-Hopf" solution is defined as $h_{t}(x):=\log z_{t}(x)$. In [14], starting from near stationary measures in a certain weakly asymmetric simple exclusion process observed in diffusive scale, this sentiment was made rigorous. Namely, it was proved that the microscopic Cole-Hopf transform of the microscopic height function, using a clever device in [25] which linearizes the simple exclusion dynamics to a more manageable system, converges to the Cole-Hopf transform of the KPZ equation, the solution to the stochastic heat equation (1.2). More recently, in [1,52] this notion of solution further gained traction in that the result in [14] was nontrivially generalized to step profile deterministic initial configurations. At the same time, in [30], it has been shown that $\log z_{t}(x)$ is the unique solution of a well-posed equation on a torus, derived from a "rough paths" approximation of (1.1), so that it is clear what sort of KPZ equation the "Cole-Hopf" solution actually solves.

In this article, another approach is considered which allows to generalize the types of microscopic particle interactions considered, given that the device in [25]
seems limited to simple exclusion and a few variants such as $q$-TASEP dynamics [15]. At the microscopic level, the height function $H_{t}(x)$, evaluated for $t \geq 0$ and $x \in \mathbb{Z}$, takes form

$$
H_{t}(x)= \begin{cases}J_{0}(t)-\sum_{y=0}^{x-1} \eta_{t}(y), & \text { for } x \geq 1  \tag{1.3}\\ J_{0}(t), & \text { for } x=0 \\ J_{0}(t)+\sum_{y=x}^{-1} \eta_{t}(y), & \text { for } x \leq-1\end{cases}
$$

where $J_{y}(t)$ is the current across bond $(y-1, y)$ and $\eta_{t}(y)$ is the particle number at $y \in \mathbb{Z}$ at time $t \geq 0$. Then the discrete gradients of the microscopic height function are the particle numbers, $H_{t}(x+1)-H_{t}(x)=\eta_{t}(x)$, and the corresponding fluctuation field examined in diffusive scale, that is, when time is scaled in terms of $n^{2}$ and space is scaled by $n$, is the particle density fluctuation field $\mathcal{Y}_{t}^{n}$. The guiding idea is that $\mathcal{Y}_{t}^{n}$ should converge to $\mathcal{Y}_{t}=\nabla h_{t}$ in some sense.

Formally, by carrying through the " $\nabla$ " operation, $\mathcal{Y}_{t}$ satisfies a type of stochastic Burgers equation,

$$
\begin{equation*}
\partial_{t} \mathcal{Y}_{t}(x)=D \Delta \mathcal{Y}_{t}(x)+a \nabla\left(\mathcal{Y}_{t}(x)\right)^{2}+\sigma \nabla \dot{\mathcal{W}}_{t}(x) \tag{1.4}
\end{equation*}
$$

which again for the same reasons as for the original KPZ equation is ill-posed when $a \neq 0$. If $a=0$, however, it is a type of Ornstein-Uhlenbeck equation which possesses a unique solution when starting from a large class of initial distributions (cf. [13, 61]).

A main contribution of the article is to understand the derived stochastic Burgers equation (1.4) in the context of a general class of nearest-neighbor weakly asymmetric interacting particle systems on $\mathbb{Z}$, starting from perturbations of the invariant measure $v_{\rho}$. This class is composed of systems with "gradient" dynamics, not necessarily product invariant measures, sufficient spectral gap and "equivalence of ensembles" estimates among other technical conditions (cf. Section 2.1), which include in particular the already studied simple exclusion process, and also zerorange and exclusion models with kinetically constrained or speed-change interactions, which have varying and sometimes slow mixing behaviors. The initial distributions consist of "bounded entropy" perturbations of the invariant measure $v_{\rho}$ (cf. Section 2.1 for a precise statement).

Our results describe the limit points of the fluctuation field $\mathcal{Y}_{t}^{n, \gamma}$ in diffusive scale, in a reference frame moving with a process characteristic velocity $v_{n}(t) \sim n^{-1}\left\lfloor n^{2}\left(p_{n}-q_{n}\right) v t\right\rfloor$. Here, $p_{n}-q_{n}$ is the difference of the single particle jump rates which identifies the strength of the weak asymmetry considered, and $v$ is a homogenized velocity parameter depending on the particle dynamics. Given the size of $p_{n}-q_{n}$, a dichotomy emerges in the form of the limits derived. Namely, for $p_{n}-q_{n}=O\left(n^{-\gamma}\right)$, when $1 / 2<\gamma \leq 1$, we show a "crossover
result" (Theorem 2.2) that $\mathcal{Y}_{t}^{n, \gamma}$ converges to an Ornstein-Uhlenbeck field with certain homogenized parameters. When $\gamma=1$, convergence of $\mathcal{Y}_{t}^{n, \gamma}$ to the same Ornstein-Uhlenbeck field has been known for many particle systems since the work [18]. For discussions of "crossover" results with respect to simple exclusion, see [28, 53].

However, when $\gamma=1 / 2$, a critical value, we prove (Theorem 2.3) that limit points of $\mathcal{Y}_{t}^{n, \gamma}$ satisfy a martingale formulation, which we dub as an "energy" formulation (cf. Theorem 2.3), also with homogenized constants, which interprets the stochastic Burgers equation: Namely, the nonlinear term in (1.4) is understood in terms of a certain Cauchy limit of a function of the fluctuation field acting on an approximation of a point mass as the approximation becomes more refined. We remark, however, with respect to simple exclusion processes, convergence of $\mathcal{Y}_{t}^{n, \gamma}$ to a unique limit when $\gamma=1 / 2$ is already known, and this limit is understood in the "Cole-Hopf" sense as mentioned above [14]. Therefore, our results imply that the "Cole-Hopf" limit of the fluctuation field satisfies also our "energy" formulation. In this context, we note [3] further clarifies the "energy" formulation of the simple exclusion limit starting from the invariant state $v_{\rho}$ (cf. point 2 of Remark 2.4). Also, we note another martingale formulation was given with respect to the Burgers equation in [4].

In our general framework, convergence of $\mathcal{Y}_{t}^{n, \gamma}$ to a unique limit when $\gamma=1 / 2$ has not been shown, an important open question (cf. Remark 2.4). However, one may still try to characterize limit points of the height function across process characteristics, $H_{t}^{n, \gamma}(x):=n^{-1 / 2} H_{n^{2} t}\left(n x-n v_{n}(t)\right)$, via (1.3) given subsequential convergence of $\mathcal{Y}_{t}^{n, \gamma}$. Although this is not the purpose of this paper, we indicate how this might be accomplished to be more complete. Indeed, by (1.3) and $J_{0}(t)-J_{x}(t)=\sum_{y=0}^{x-1}\left(\eta_{t}(y)-\eta_{0}(y)\right)$, one has $H_{t}^{n, \gamma}(x)=$ $n^{-1 / 2} J_{n x-n v_{n}(t)}\left(n^{2} t\right)-n^{-1 / 2} \sum_{y=0}^{n x-n v_{n}(t)-1} \eta_{0}(y)$, say for $x>v_{n}(t)$. To write the current in terms of the fluctuation field, formally, $n^{-1 / 2} J_{n x-n v_{n}(t)}\left(n^{2} t\right)=$ $\mathcal{Y}_{t}^{n, \gamma}\left(1_{[x, \infty)}\right)-\mathcal{Y}_{0}^{n, \gamma}\left(1_{[x, \infty)}\right)+o(1)$, although as there are an infinite number of particles and $1_{[x, \infty)}$ is not a compactly supported function some sort of truncation is needed to make a rigorous argument. Using the method in [51] and [32], one can approximate $n^{-1 / 2} J_{n x-n v_{n}(t)}\left(n^{2} t\right)$ by $\mathcal{Y}_{t}^{n, \gamma}\left(G_{k, x}\right)-\mathcal{Y}_{0}^{n, \gamma}\left(G_{k, x}\right)$ for large $k$ where $G_{k, x}(z)=\left(1-(z-x)_{+} / k\right)_{+}$, and so it is possible to take subsequential limits of $H_{t}^{n, \gamma}$.

Finally, we remark if uniqueness of solution for the $\gamma=1 / 2$ "energy" formulation were known in our more general framework, one would be able to identify the solution, modulo parameters, as the limit already identified for simple exclusion through the Cole-Hopf apparatus. In this way, one should be able to determine that the height function limits, with respect to a general class of interactions starting from nearly the invariant measure, are in the KPZ class for instance.

We now remark on the argument for Theorems 2.2 and 2.3. We take a stochastic differential of $\mathcal{Y}_{t}^{n, \gamma}$, namely

$$
d \mathcal{Y}_{t}^{n, \gamma}=\left(\partial_{t} \mathcal{Y}_{t}^{n, \gamma}+L_{n} \mathcal{Y}_{t}^{n, \gamma}\right) d t+d \mathcal{M}_{t}^{n, \gamma}
$$

where $L_{n}$ is the system infinitesimal generator and $\mathcal{M}_{t}^{n, \gamma}$ is a martingale. We note, because the reference frame moves with velocity $v_{n}(t)$, the term $\partial_{t} \mathcal{Y}_{t}^{n, \gamma}$ does not vanish. Beginning in a perturbed invariant measure, the martingale term can be handled by an ergodic theorem. However, to write the drift term $\partial_{t} \mathcal{Y}_{t}^{n, \gamma}+L_{n} \mathcal{Y}_{t}^{n, \gamma}$, in terms of the fluctuation field itself and, therefore, "close" the equation, is a more difficult task, and requires what has been known as a "Boltzmann-Gibbs" principle. Such a principle was first proved in [18] when $\gamma=1$. In our context, we would like to recover a second-order term, and the principle would replace

$$
\int_{0}^{t} \frac{1}{n^{\gamma-1 / 2}} \sum_{x \in \mathbb{Z}} \nabla G(x / n) \tau_{x} V\left(\eta_{n^{2} s}\right) d s
$$

with

$$
\begin{aligned}
\frac{\varphi_{V}^{\prime \prime}(\rho)}{2} \int_{0}^{t} & \frac{1}{n^{\gamma+1 / 2}} \sum_{x \in \mathbb{Z}} \nabla G(x / n) \\
& \times\left\{\mathcal{Y}_{s}^{n, \gamma}\left(\frac{1}{2 \varepsilon} 1_{[x-\varepsilon, x+\varepsilon]}\right)^{2}-\mathbb{E}_{\nu_{\rho}}\left[\mathcal{Y}_{s}^{n, \gamma}\left(\frac{1}{2 \varepsilon} 1_{[x-\varepsilon, x+\varepsilon]}\right)^{2}\right]\right\} d s
\end{aligned}
$$

in $L^{2}\left(\mathbb{P}_{v_{\rho}}\right)$ as $n \uparrow \infty$ and $\varepsilon \downarrow 0$. Here, $G$ is a function in the Schwarz class, $\tau_{x}$ is the $x$-shift operator, $V$ is a mean-zero function with the property that the derivative of its "tilted mean" $\varphi_{V}(z)$ vanishes at $z=\rho$ [cf. definition near (2.4)]. Given such a replacement principle (cf. Section 3.2 for precise statements), one can prove the sequence $\mathcal{Y}_{t}^{n, \gamma}$ is tight and derive martingale formulations of limit points as desired.

The case $\gamma=1 / 2$ is the most difficult since there is no spatial averaging at all. However, there is much cancelation with respect to the time integral which helps to prove the estimate needed. We show the cases $1 / 2<\gamma \leq 1$ would follow from the $\gamma=1 / 2$ replacement. A similar replacement for symmetric simple exclusion, using specific duality methods, was performed in [5].

The method given here, in our general framework, is quite different. The main idea is to use an involved $H_{-1}$ renormalization scheme to bound errors in the replacement. Such a scheme makes good use of three assumed ingredients [cf. precise statements (R), (G), (EE) in Section 2.1]: First, the measure $v_{\rho}$ is invariant with respect to all asymmetric and symmetric versions of the process, the main reason for the "gradient dynamics" condition. Second, a spectral gap lower bound for the symmetric process localized on a interval $\Lambda_{\ell}$ with width $\ell$ and $\sum_{x \in \Lambda_{\ell}} \eta(x)$ particles which, after averaging with respect to $v_{\rho}$, is of order $O\left(\ell^{-2}\right)$. Also, third, an "equivalence of ensembles" estimate holds with respect to canonical $v_{\rho}\left(\cdot \mid \sum_{|x| \leq \ell} \eta(x)=k\right)$ and grand canonical $v_{\rho}$ measures.

We note the current article is an evolution of the arXiv paper [27], encompassing the work there on a type of exclusion model starting from a Bernoulli product invariant measure and a model specific Boltzmann-Gibbs principle. See also [3] for
a different type of resolvent method specific to simple exclusion. In this context, the current article is a nontrivial generalization to more diverse models, starting from perturbations of the stationary state, using a more general $H_{-1}$ renormalization scheme. We remark that part of this improvement, of its own interest, is that the Boltzmann-Gibbs principle (Theorem 3.2) shown here does not rely on the independence structure of a product measure, or on a sharp spectral gap estimate, or on a process "duality." Finally, we note elements of our $H_{-1}$ renormalization scheme go back to [26] and [56] in different contexts.

We now give the structure of the article. In Section 2, the general class of models studied, results and specific systems satisfying the class assumptions are discussed. Then, in Section 3, we outline the proof of the main results, Theorems 2.2 and 2.3, stating the form of Boltzmann-Gibbs principle used. In Section 4, this principle is proved. Finally, in Section 5, we prove for a class of systems, including the specific processes discussed in Section 2, the "equivalence of ensembles" estimate assumed for the proofs in Section 3.
2. Abstract framework, results and models. We now discuss the abstract framework we work with in Section 2.1, and state results in this framework in Section 2.2. This framework covers a wide class of models such as zero-range models and different types of exclusion processes which we detail in Sections 2.32.5. A reader focusing on one of these models, might skip to its subsection while referring to Section 2.1, and then proceed to results in Section 2.2.
2.1. Notation and assumptions. We consider a sequence of "weakly asymmetric" nearest-neighbor "mass conservative" particle systems $\left\{\eta_{t}^{n}: t \geq 0\right\}$ on the state space $\Omega=\mathbb{N}_{0}^{\mathbb{Z}}$ where $\mathbb{N}_{0}=\{0,1,2, \ldots\}$. The configuration of the system $\eta_{t}=\left\{\eta_{t}(x): x \in \mathbb{Z}\right\}$ is a collection of occupation numbers $\eta_{t}(x)$ which counts the numbers of particles at sites $x \in \mathbb{Z}$ at time $t \geq 0$. In some of the examples we will consider, the occupation number is bounded by 1 , in which case the effective state space reduces to $\{0,1\}^{\mathbb{Z}}$.
"Gradient" dynamics. The dynamics will be of "gradient" type. That is, we suppose there are functions $\left\{b_{x}^{R, n}\right\}_{n \geq 1}$ and $\left\{b_{x}^{L, n}\right\}_{n \geq 1}$ satisfying the following conditions (R1) and (R2). Let $\tau_{x}$ be the shift operator where $\left(\tau_{x} \eta\right)(z)=\eta(x+z)$ and $\tau_{x} f(\eta)=f\left(\tau_{x} \eta\right)$ for $x \in \mathbb{Z}$. Let also $\Lambda_{k}=\{j:|j| \leq k\} \subset \mathbb{Z}$ for $k \geq 1$.
(R1) For all $n \geq 1, b_{x}^{R, n}=\tau_{x} b_{0}^{R, n}$ and $b_{x}^{L, n}=\tau_{x} b_{0}^{L, n}$ are nonnegative finiterange functions on $\Omega$ such that $b_{0}^{R, n}$ and $b_{0}^{L, n}$ are supported on $\{\eta(y): y \in$ $\left.\Lambda_{R}\right\}$ for some $R \geq 1$. We suppose uniformly in $n$ that $\left|b_{0}^{R, n}(\eta)\right|+\left|b_{0}^{L, n}(\eta)\right| \leq$ $C \sum_{y \in \Lambda_{R}} \eta(y)$. Moreover, there are nonnegative functions $c_{x}^{n}=\tau_{x} c^{n}$ on $\Omega$, supported on $\left\{\eta(x): x \in \Lambda_{R}\right\}$ such that

$$
b_{x}^{R, n}(\eta)-b_{x}^{L, n}(\eta)=c_{x}^{n}(\eta)-c_{x+1}^{n}(\eta)
$$

In addition, suppose there are fixed functions $b_{0}^{R}, b_{0}^{L}$ and $c_{0}$ such that configurationwise

$$
\lim _{n \uparrow \infty} b_{0}^{R, n}(\eta)=b_{0}^{R}(\eta), \quad \lim _{n \uparrow \infty} b_{0}^{L, n}(\eta)=b_{0}^{L}(\eta) \quad \text { and } \quad \lim _{n \uparrow \infty} c_{0}^{n}(\eta)=c_{0}(\eta)
$$

In some of the models considered, such as zero-range processes in Section 2.3, the functions $b_{0}^{R, n}=b_{0}^{R}, b_{0}^{L, n}=b_{0}^{L}$ and $c_{0}^{n}=c_{0}$ are fixed and do not depend on the parameter $n$. However, for the kinetically constrained exclusion models in Section 2.4 , the rates do depend on $n$.
(R2) With respect to a fixed measure $\nu_{\rho}$ on $\Omega$, for all $n \geq 1$, we have

$$
b_{x}^{R, n}\left(\eta^{x+1, x}\right) \frac{d v_{\rho}^{x+1, x}}{d v_{\rho}}(\eta)=b_{x}^{L, n}(\eta)
$$

where $\nu_{\rho}^{x+1, x}$ is the measure of the variable $\zeta=\eta^{x+1, x}$ under $v_{\rho}$.
We also define $b_{x}^{n}(\eta)=b_{x}^{R, n}(\eta)+b_{x}^{L, n}(\eta), b^{n}(\eta)=b_{0}^{n}(\eta)$ and $c^{n}(\eta)=c_{0}^{n}(\eta)$ to simplify notation.

We now specify the process generator. For $a \in \mathbb{R}$ and $\gamma>0$, let

$$
p_{n}=\frac{1}{2}+\frac{a}{2 n^{\gamma}} \quad \text { and } \quad q_{n}=1-p_{n}=\frac{1}{2}-\frac{a}{2 n^{\gamma}}
$$

Let also $n_{0}$ be such that $0 \leq p_{n_{0}}, q_{n_{0}} \leq 1$, and $T>0$ be a fixed time.
(M) Suppose, for each $a \in \mathbb{R}, \gamma>0$ and $n \geq n_{0}$, that $\left\{\eta_{t}^{n}: t \in[0, T]\right\}$ is a $L^{2}\left(v_{\rho}\right)$ Markov process with strongly continuous Markov semigroup $P_{t}^{n}$ and Markov generator $L_{n}$ (cf. Chapter I; Section IV. 4 of [40]) with a core composed of local $L^{2}\left(v_{\rho}\right)$ functions on which

$$
\begin{equation*}
L_{n} f(\eta)=n^{2} \sum_{x \in \mathbb{Z}}\left\{b_{x}^{R, n}(\eta) p_{n} \nabla_{x, x+1} f(\eta)+b_{x}^{L, n}(\eta) q_{n} \nabla_{x+1, x} f(\eta)\right\} \tag{2.1}
\end{equation*}
$$

where $\nabla_{x, y} f(\eta)=f\left(\eta^{x, y}\right)-f(\eta)$, and $\eta^{x, y}$ is the configuration obtained from $\eta$ by moving a particle from $x$ to $y$ :

$$
\eta^{x, y}(z)= \begin{cases}\eta(y)+1, & \text { when } z=y \\ \eta(x)-1, & \text { when } z=x \\ \eta(z), & \text { otherwise }\end{cases}
$$

The role of $a \in \mathbb{R}$ and $\gamma>0$ is to control the strength of the "weak asymmetry" in the model.

Invariant measure $v_{\rho}$. We now specify some technical properties which $v_{\rho}$ should satisfy. Define for a probability measure $\kappa$, the path measure $\mathbb{P}_{\kappa}$ governing the process $\left\{\eta_{t}^{n}: t \in[0, T]\right\}$ with initial configurations $\eta_{0}$ distributed according to $\kappa$. Let then $E_{\kappa}$ and $\mathbb{E}_{\kappa}$ denote expectations with respect to $\kappa$ and $\mathbb{P}_{\kappa}$, respectively.
(IM1) Suppose $v_{\rho}$ is a translation-invariant measure which is "spatially mixing." That is, for local $L^{2}\left(v_{\rho}\right)$ functions $f$ and $h$,

$$
\lim _{|x| \uparrow \infty} E_{v_{\rho}}\left[f(\eta) \tau_{x} h(\eta)\right]=E_{v_{\rho}}[f] E_{v_{\rho}}[h] .
$$

In addition, suppose the mean $E_{\nu_{\rho}}[\eta(0)]=\rho$, and moment-generating function $E_{\nu_{\rho}}\left[e^{\lambda \eta(0)}\right]<\infty$ for $|\lambda| \leq \lambda^{*}$ for a $\lambda^{*}>0$.
Although product measures $v_{\rho}$ are considered in most of the examples, we note, in Section 2.5 , a nonproduct measure $\nu_{\rho}$ corresponding to an exponentially mixing ergodic Markov chain is used.

Now, the measure $v_{\rho}$, by (IM1) and the "gradient dynamics" conditions (R1) and (R2), is an invariant measure with respect to $L_{n}$ for all $a \in \mathbb{R}$ and $\gamma>0$. Indeed, let $\phi$ be a local $L^{2}\left(v_{\rho}\right)$ function supported with respect to sites in $\Lambda_{k}$. Then, for $\ell>k$, we have

$$
\begin{aligned}
E_{v_{\rho}}\left[L_{n} \phi\right] & =-n^{2} E_{v_{\rho}}\left[\sum_{|x| \leq \ell}\left(p_{n}-q_{n}\right) \phi(\eta)\left[c_{x}^{n}(\eta)-c_{x+1}^{n}(\eta)\right]\right] \\
& =-n^{2}\left(p_{n}-q_{n}\right) E_{v_{\rho}}\left[\phi(\eta)\left(c_{-\ell}^{n}(\eta)-c_{\ell+1}^{n}(\eta)\right)\right] .
\end{aligned}
$$

The limit as $\ell \uparrow \infty$ vanishes, by translation-invariance and the spatial mixing assumption in (IM1).

One can also compute that the $L^{2}\left(v_{\rho}\right)$ adjoint $L_{n}^{*}$ is the generator with parameter $-a$, that is, when the jump probability is reversed. Define $S_{n}=\left(L_{n}+L_{n}^{*}\right) / 2$. Then the Dirichlet form $D_{v_{\rho}, n}(f):=E_{v_{\rho}}\left[f\left(-L_{n} f\right)\right]=E_{v_{\rho}}\left[f\left(-S_{n} f\right)\right]$ on local $L^{2}\left(v_{\rho}\right)$ functions is given by

$$
\begin{equation*}
D_{v_{\rho}}(f)=\frac{1}{2} \sum_{x \in \mathbb{Z}} E_{v_{\rho}}\left[b_{x}^{R, n}(\eta)\left(\nabla_{x, x+1} f(\eta)\right)^{2}\right] \tag{2.2}
\end{equation*}
$$

Moreover, when $a=0, S_{n}$ is the generator of the associated process and $\nu_{\rho}$ is a reversible measure.

Consider now the empirical measure

$$
\mathcal{Y}_{0}^{n}=\frac{1}{\sqrt{n}} \sum_{x \in \mathbb{Z}}(\eta(x)-\rho) \delta_{x / n}
$$

and its covariance under measure $\kappa$, on compactly supported functions,

$$
\mathcal{C}_{\kappa}^{n}(G, H)=E_{\kappa}\left[\left(\mathcal{Y}_{0}^{n}(G)-E_{\kappa}\left[\mathcal{Y}_{0}^{n}(G)\right]\right)\left(\mathcal{Y}_{0}^{n}(H)-E_{\kappa}\left[\mathcal{Y}_{0}^{n}(H)\right]\right)\right] .
$$

(IM2) We assume, starting from $\nu_{\rho}$, that $\mathcal{Y}_{0}^{n}$ converges weakly to a spatial Gaussian process with covariance $\mathcal{C}_{v_{\rho}}(G, H):=\lim _{n \uparrow \infty} \mathcal{C}_{v_{\rho}}^{n}(G, H)$ such that $\mathcal{C}_{v_{\rho}}(G, G) \leq C(\rho)\|G\|_{L^{2}(\mathbb{R})}^{2}$. Also, suppose the moment bound holds $\sup _{\ell \geq 1}\left\|\left(\frac{1}{\sqrt{\ell}} \sum_{x=1}^{\ell}(\eta(x)-\rho)\right)^{2}\right\|_{L^{4}\left(v_{\rho}\right)}<\infty$.

It will be convenient to define the variances

$$
\sigma_{n}^{2}(\rho):=\mathcal{C}_{v_{\rho}}^{n}(H, H)=E_{v_{\rho}}\left[\left(\frac{1}{\sqrt{2 n+1}} \sum_{x \in \Lambda_{n}}(\eta(x)-\rho)\right)^{2}\right]
$$

and $\sigma^{2}(\rho)=\mathcal{C}_{v_{\rho}}(H, H)=\lim _{n \uparrow \infty} \sigma_{n}^{2}(\rho)$ when $H(x)=1_{[-1,1]}(x)$.
When $v_{\rho}$ is sufficiently mixing, the case of our examples, (IM2) holds with $\mathcal{C}_{v_{\rho}}(G, H)=\sigma^{2}(\rho)\langle G, H\rangle_{L^{2}(\mathbb{R})}$.

Now, for $\lambda \in\left(-\lambda^{*}, \lambda^{*}\right)$, consider the tilted measure $v_{\rho}^{\lambda}$ with "tilt" or "chemical potential" $\lambda$ given by its finite-dimensional projections

$$
\begin{equation*}
\frac{d v_{\rho}^{\lambda}}{d v_{\rho}}\left(\eta(x)=e(x), x \in \Lambda_{\ell} \mid \eta(x)=\xi(x), x \notin \Lambda_{\ell}\right)=\frac{e^{\lambda \sum_{x \in \Lambda_{\ell}}(e(x)-\rho)}}{Z(\lambda, \ell, \xi)} \tag{2.3}
\end{equation*}
$$

where $e, \xi \in \Omega$ and $Z(\lambda, \ell, \xi)$ is the normalization.
(D1) We will assume the measures $\left\{\nu_{\rho}^{\lambda}:|\lambda|<\lambda^{*}\right\}$ are well defined on $\Omega$, that is a limit of (2.3) as $\Lambda_{\ell} \nearrow \mathbb{Z}$ can be taken, not depending on $\xi$. Also, we assume that the measures can be indexed by density, that is, $E_{\nu_{\rho}^{\lambda}}[\eta(0)]$ is strictly increasing in $\lambda$ for $|\lambda| \leq \lambda^{*}$.

These assumptions hold when $v_{\rho}$ is a nontrivial product measure satisfying (IM1): $\frac{d}{d \lambda} E_{v_{\rho}^{\lambda}}[\eta(0)]=E_{v_{\rho}^{\lambda}}\left[\left(\eta(0)-E_{v_{\rho}^{\lambda}}[\eta(0)]\right)^{2}\right]>0$. They also hold when $v_{\rho}$ corresponds to the ergodic Markov chain in the case for the exclusion with speed-change model (cf. details in Section 2.5).

The measures $\left\{v_{\rho}^{\lambda}:|\lambda|<\lambda^{*}\right\}$ are translation-invariant since $v_{\rho}$ is assumed translation-invariant (IM1). Also, given exponential moments of $v_{\rho}$ (IM1), $E_{\nu_{\rho}^{\lambda}}[\eta(0)]$ is continuous in $\lambda$ for $|\lambda|<\lambda^{*}$. Hence, by the strict increasing assumption in (D1), one can reparameterize $\left\{v_{\rho}^{\lambda}\right\}$ in terms of density: Let $z \in$ $\left(\rho_{*}, \rho^{*}\right)$ where $\rho_{*}=\lim _{\lambda \downarrow-\lambda^{*}} E_{\nu_{\rho}^{\lambda}}[\eta(0)]$ and $\rho^{*}=\lim _{\lambda \uparrow \lambda^{*}} E_{\nu_{\rho}^{\lambda}}[\eta(0)]$. Let $\lambda(z) \in$ $\left(-\lambda^{*}, \lambda^{*}\right)$ be the parameter such that $E_{\nu_{\rho}^{\lambda(z)}}[\eta(0)]=z$. Then we will define $v_{z}=v_{\rho}^{\lambda(z)}$.

Define also, for a local $L^{2}\left(v_{\rho}\right)$ function $f$, the "tilted mean" function $\varphi_{f}(z)$ : $\left(\rho_{*}, \rho^{*}\right) \rightarrow \mathbb{R}$ where

$$
\varphi_{f}(z)=E_{v_{z}}[f(\eta)]
$$

We consider the derivatives of $\varphi_{f}(z)$ as the formal limits of the derivatives of $E_{\nu_{z}}\left[f(\eta) \mid \eta(x)=\xi(x), x \in \Lambda_{\ell}\right]$ as $\ell \uparrow \infty$. Define

$$
\varphi_{f}^{\prime}(z):=\lambda^{\prime}(z) E_{\nu_{z}}\left[\left(f(\eta)-E_{\nu_{z}}[f]\right)\left(\sum_{x \in \mathbb{Z}}(\eta(x)-z)\right)\right]
$$

$$
\begin{align*}
\varphi_{f}^{\prime \prime}(z):= & \left(\lambda^{\prime}(z)\right)^{2} E_{v_{z}}\left[\left(f(\eta)-E_{v_{z}}[f]\right)\left(\sum_{x \in \mathbb{Z}}(\eta(x)-z)\right)^{2}\right] \\
& +\lambda^{\prime \prime}(z) E_{v_{z}}\left[\left(f(\eta)-E_{v_{z}}[f]\right)\left(\sum_{x \in \mathbb{Z}}(\eta(x)-z)\right)\right] \tag{2.4}
\end{align*}
$$

For the 0th derivative, we set $\varphi_{f}^{(0)}(z):=E_{\nu_{z}}[f]$.
(D2) For local $L^{2}\left(v_{\rho}\right)$ functions $f$, suppose the limits (2.4) are well defined and $\left|\varphi_{f}^{\prime}(\rho)\right|,\left|\varphi_{f}^{\prime \prime}(\rho)\right| \leq C(\rho)\|f\|_{L^{2}\left(\nu_{\rho}\right)}$; already, $\left|\varphi_{f}(\rho)\right| \leq\|f\|_{L^{2}\left(\nu_{\rho}\right)}$. Also, suppose

$$
\lim _{n \uparrow \infty} \varphi_{f_{n}}^{\prime}(\rho)=\varphi_{f}^{\prime}(\rho) \quad \text { and } \quad \lim _{n \uparrow \infty} \varphi_{f_{n}}^{\prime \prime}(\rho)=\varphi_{f}^{\prime \prime}(\rho)
$$

when $\left\{f_{n}\right\}$ and $f$ are local functions such that $\lim _{n \uparrow \infty} f_{n}(\eta)=f(\eta)$ and $f_{n}(\eta) \leq$ $\hat{f}(\eta)$ configuration-wise for each $n$ where $\hat{f} \in L^{2}\left(v_{\rho}\right)$.

When $\left\{v_{x}\right\}$ are product or rapidly mixing Markov measures, again the case for our examples, this condition also holds by calculation with (2.3).

Spectral gap. We now give a "spectral gap" condition. For $\ell \geq 1$, recall $\Lambda_{\ell}$ is the box of size $2 \ell+1$, namely $\Lambda_{\ell}:=\{x \in \mathbb{Z}:|x| \leq \ell\}$. Let also, for $k \geq 0$ and $\xi \in \Omega, \mathcal{G}_{k, \ell, \xi}:=\left\{\eta: \sum_{x \in \Lambda_{\ell}} \eta(x)=k, \eta(y)=\xi(y)\right.$ for $\left.y \notin \Lambda_{\ell}\right\}$ be the hyperplane of configurations on $\Lambda_{\ell}$ with $k$ particles which equal $\xi$ outside $\Lambda_{\ell}$. Denote by $\nu_{k, \ell, \xi}$ the canonical measure on $\mathcal{G}_{k, \ell, \xi}$, namely

$$
v_{k, \ell, \xi}(\cdot):=v_{\rho}\left(\cdot \mid \sum_{x \in \Lambda_{\ell}} \eta(x)=k, \eta(y)=\xi(y) \text { for } y \notin \Lambda_{\ell}\right)
$$

Consider now the process, restricted to the hyperplane $\mathcal{G}_{k, \ell, \xi}$ with generator

$$
\mathcal{S}_{n, \mathcal{G}_{k, \ell, \xi}} f(\eta)=\frac{1}{2} \sum_{\substack{|x-y|=1 \\ x, y \in \Lambda_{\ell}}} b_{x}^{n}(\eta) \nabla_{x, y} f(\eta)
$$

This is a finite-state Markov process with reversible invariant measure $v_{k, \ell, \xi}$. Denote by $\lambda_{k, \ell, \xi, n}$ the spectral gap, that is the second largest eigenvalue of $-\mathcal{S}_{n, \mathcal{G}_{k, \ell, \xi}}$ (with 0 being the largest). Let $W(k, \ell, \xi, n)$ denote the reciprocal of $\lambda_{k, \ell, \xi, n}$, which is set to $\infty$ if $\lambda_{k, \ell, \xi, n}=0$. Then the associated Poincaré-inequality reads as

$$
\begin{equation*}
\operatorname{Var}\left(f, v_{k, \ell, \xi}\right) \leq W(k, \ell, \xi, n) \mathcal{D}_{n}\left(f, v_{k, \ell, \xi}\right) \tag{2.5}
\end{equation*}
$$

where $\operatorname{Var}\left(f, v_{k, \ell, \xi}\right)$ is the variance of $f$ with respect to $v_{k, \ell, \xi}$ and the canonical Dirichlet form $\mathcal{D}_{n}\left(f, v_{k, \ell, \xi}\right)$ is given by

$$
\mathcal{D}_{n}\left(f, v_{k, \ell, \xi}\right):=\frac{1}{2} \sum_{x, x+1 \in \Lambda_{\ell}} E_{v_{k, \ell, \xi}}\left[b_{x}^{R, n}(\eta)\left(\nabla_{x, x+1} f(\eta)\right)^{2}\right]
$$

When $W(k, \ell, \xi, n)<\infty$, the process is ergodic and $\nu_{k, \ell, \xi}$ is the unique invariant measure.

Denote the "outside variables" by $\eta_{\ell}^{c}=\left\{\eta(x): x \notin \Lambda_{\ell}\right\}$. We will assume the following condition on $W(k, \ell, \xi, n)$.
(G) Suppose there is a constant $C=C(\rho)$ such that, for $n \geq 1$, we have

$$
E_{\nu_{\rho}}\left[W\left(\sum_{x \in \Lambda_{\ell}} \eta(x), \ell, \eta_{\ell}^{c}, n\right)^{2}\right] \leq C \ell^{4}
$$

We remark a sufficient condition to verify (G) would be the uniform bound $\sup _{k, \xi, n} \ell^{-2} W(k, \ell, \xi, n)<\infty$, which holds for some types but not all of the specific models discussed.

Equivalence of ensembles. We will also assume an "equivalence of ensembles" estimate between the canonical and grand-canonical measures. Define, for $\ell \geq 1$ and $\eta \in \Omega$, the empirical average

$$
\eta^{(\ell)}=\frac{1}{2 \ell+1} \sum_{y \in \Lambda_{\ell}} \eta(y)
$$

(EE) For local $L^{5}\left(v_{\rho}\right)$ functions $f$, supported on $\left\{\eta(x): x \in \Lambda_{\ell_{0}}\right\}$, such that $\varphi_{f}(\rho)=\varphi_{f}^{\prime}(\rho)=0$, and $\ell \geq \ell_{0}$, there exist constants $\alpha_{0}>0$ and $C=C\left(\rho, \ell_{0}, \alpha_{0}\right)$ where

$$
\left\|E_{v_{\rho}}\left[f \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]-\frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\left[\left(\eta^{(\ell)}-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1}\right]\right\|_{L^{4}\left(v_{\rho}\right)} \leq \frac{C\|f\|_{L^{5}\left(v_{\rho}\right)}}{\ell^{1+\alpha_{0} / 2}}
$$

On the other hand, when only $\varphi_{f}(\rho)=0$ is known,

$$
\left\|E_{v_{\rho}}\left[f \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]-\varphi_{f}^{\prime}(\rho)\left(\eta^{(\ell)}-\rho\right)\right\|_{L^{4}\left(v_{\rho}\right)} \leq \frac{C\|f\|_{L^{5}\left(v_{\rho}\right)}}{\ell^{1 / 2+\alpha_{0} / 2}}
$$

We remark, a weaker version, where the $L^{2}\left(v_{\rho}\right)$ norm, instead of the $L^{4}\left(v_{\rho}\right)$ norm of the difference, is say less than the same right-hand side expressions with $\|f\|_{L^{3}\left(v_{\rho}\right)}$ in place of $\|f\|_{L^{5}\left(v_{\rho}\right)}$ would be sufficient for our purposes if there is a uniform bound on the inverse gap: $\sup _{k, \xi, n} \ell^{-2} W(k, \ell, \xi, n)<\infty$.

Usually, such estimates follow from a local central limit theorem. In Proposition 5.1, we show, when $v_{\rho}$ is a nondegenerate product measure, that (EE) holds with $\alpha_{0}=1$. In Proposition 5.2, with respect to a Markovian measure, we prove (EE) holds with $\alpha_{0}=1-\varepsilon$ for any fixed $0<\varepsilon<1$. These two propositions cover the examples discussed in the article.

Initial conditions. We will start from initial measures $\left\{\mu^{n}\right\}$ which have bounded relative entropy $H\left(\mu^{n} ; v_{\rho}\right)$ with respect to $v_{\rho}$.
(BE) Suppose $\left\{\mu^{n}\right\}$ satisfies

$$
\sup _{n} H\left(\mu^{n} ; v_{\rho}\right)=\sup _{n} E_{v_{\rho}}\left[\frac{d \mu^{n}}{d v_{\rho}} \log \frac{d \mu^{n}}{d v_{\rho}}\right]<\infty .
$$

In addition, we presume a diffusive initial limit starting from $\left\{\mu^{n}\right\}$.
(CLT) Under initial measures $\left\{\mu^{n}\right\}$, we suppose $\mathcal{Y}_{0}^{n}$ converges weakly to a spatial Gaussian process $\overline{\mathcal{Y}}_{0}$ with covariance $\mathcal{C}(G, H)=\lim _{n \uparrow \infty} \mathcal{C}_{\mu^{n}}^{n}(G, H)$ for compactly supported functions $G, H$.

Of course, if $\mu^{n} \equiv v_{\rho}$, (BE) and (CLT) trivially hold with $\mathcal{C}(G, H)=$ $\mathcal{C}_{v_{\rho}}(G, H)$. When $v_{\rho}$ is a product measure, a possible way to get nontrivial examples of measures $\left\{\mu^{n}\right\}$ satisfying (BE) and (CLT) is the following. For simplicity, we consider the case on which $v_{\rho}$ is a Bernoulli product measure on $\{0,1\}^{\mathbb{Z}}$. Let $\left\{\kappa_{x}^{n}: x \in \mathbb{Z}\right\}$ be a given bounded sequence and define $\mu^{n}$ as the nonhomogeneous Bernoulli product measure satisfying

$$
\mu_{n}(\eta(x)=1)=\rho+\frac{\kappa_{x}^{n}}{\sqrt{n}}
$$

A simple computation shows that

$$
H\left(\mu^{n} ; v_{\rho}\right) \leq \frac{C\left(\|\kappa\|_{\ell \infty}\right)}{n} \sum_{x \in \mathbb{Z}}\left(\kappa_{x}^{n}\right)^{2} .
$$

Therefore, taking $\kappa_{x}^{n}=\kappa(x / n)$, where $\kappa: \mathbb{R} \rightarrow \mathbb{R}$ is bounded and in $L^{2}(\mathbb{R})$, we see that $\sup _{n} H\left(\mu^{n} ; v_{\rho}\right)<\infty$, and (BE) is satisfied. On the other hand, since the measure $\mu^{n}$ is product, a simple computation shows that, under $\left\{\mu^{n}\right\}$, the process $\mathcal{Y}_{0}^{n}$ converges in distribution to $\overline{\mathcal{Y}}_{0}+\kappa$, where $\overline{\mathcal{Y}}_{0}$ is a white noise with variance $\rho(1-\rho)$. In [48], the Cole-Hopf solution of KPZ is considered starting from such initial conditions.

One may relate probabilities of events $A$ under $\mu^{n}$ with those under $v_{\rho}$ by an application of the entropy inequality:

$$
\begin{equation*}
\mathbb{P}_{\mu^{n}}(A) \leq \frac{\log 2+H\left(\mu^{n} ; v_{\rho}\right)}{\log \left(1+\mathbb{P}_{v_{\rho}}(A)^{-1}\right)} \tag{2.6}
\end{equation*}
$$

For instance, let $r \in L^{2}\left(v_{\rho}\right)$ be a local function. By the spatial mixing assumption (IM2), under $v_{\rho}$, we have the convergence in probability,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \int_{0}^{T} \frac{1}{2 n+1} \sum_{x \in \Lambda_{n}} \tau_{x} r\left(\eta_{s}^{n}\right) d s=E_{v_{\rho}}[r(\eta)] \tag{2.7}
\end{equation*}
$$

Then, by the entropy relation, also under $\left\{\mu^{n}\right\}$, the same limit also holds in probability.

Of course, given that we begin from nearly the invariant measure $v_{\rho}$, (2.7) is a trivial case of "hydrodynamics." Formally, starting from more general measures, the hydrodynamic equation for the limiting empirical density $\rho=\rho(x, t)$ would read

$$
\begin{equation*}
\partial_{t} \rho(x, t)+\frac{a}{2} \nabla \varphi_{b}(\rho(x, t))=\frac{1}{2} \Delta \varphi_{c}(\rho(x, t)) . \tag{2.8}
\end{equation*}
$$

In a sense, the main results of the paper are on the different fluctuations from the law of large numbers (2.7) which arise for different regimes of the strength asymmetry parameters $a$ and $\gamma$.
2.2. Results. Denote by $\mathbb{S}(\mathbb{R})$ the standard Schwarz space of rapidly decreasing functions equipped with the usual metric, and let $\mathbb{S}^{\prime}(\mathbb{R})$ be its dual, namely the set of tempered distributions in $\mathbb{R}$, endowed with the strong topology. Denote the density fluctuation field acting on functions $H \in \mathbb{S}(\mathbb{R})$ as

$$
\mathcal{Y}_{t}^{n}(H)=\frac{1}{\sqrt{n}} \sum_{x \in \mathbb{Z}} H\left(\frac{x}{n}\right)\left(\eta_{t}^{n}(x)-\rho\right)
$$

Denote by $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$ and $C\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$ the spaces of right continuous functions with left limits and continuous functions respectively from $[0, T]$ to $\mathbb{S}^{\prime}(\mathbb{R})$.

We now state a result from the literature which has been proved for some processes (cf. [23], Chapter 11 in [34] for zero-range processes with bounded rate, [20, 50] for simple exclusion processes, and Section II.2.10 of [57] for exclusion systems with speed-change), sometimes from more general initial conditions, when the asymmetry is of order $O\left(n^{-1}\right)$.

Proposition 2.1. For $\gamma=1$, starting from $\left\{\mu^{n}\right\}$, the sequence $\left\{\mathcal{Y}_{t}^{n} ; n \geq 1\right\}$ converges in the uniform topology on $D\left([0, T], \mathcal{S}^{\prime}(\mathbb{R})\right)$ to the process $\mathcal{Y}_{t}$ which solves the Ornstein-Uhlenbeck equation

$$
\begin{equation*}
\partial_{t} \mathcal{Y}_{t}=\frac{1}{2} \varphi_{c}^{\prime}(\rho) \Delta \mathcal{Y}_{t}+\frac{a}{2} \varphi_{b}^{\prime}(\rho) \nabla \mathcal{Y}_{t}+\sqrt{\frac{1}{2} \varphi_{b}(\rho)} \nabla \dot{\mathcal{W}}_{t} \tag{2.9}
\end{equation*}
$$

where $\dot{\mathcal{W}}_{t}$ is a space-time white noise with unit variance, and $\mathcal{Y}_{0}=\overline{\mathcal{Y}}_{0}$, the field given in (CLT).

The Ornstein-Uhlenbeck equation (2.9) has a drift term coming from the weak asymmetry of the jump rates. The drift, as is well known, can be understood in terms of a characteristic velocity $v=(a / 2) \varphi_{b}^{\prime}(\rho)$ from considering the linearization of the hydrodynamic equation (2.8) (cf. Chapter II. 2 of [57]). However, it can be removed from the limit field by observing the density fluctuation field in the frame of an observer moving along the process characteristics. Define

$$
\mathcal{Y}_{t}^{n, \gamma}(H)=\frac{1}{\sqrt{n}} \sum_{x \in \mathbb{Z}} H\left(\frac{x}{n}-\frac{1}{n}\left\{\frac{a \varphi_{b^{n}}^{\prime}(\rho) t n^{2}}{2 n^{\gamma}}\right\}\right)\left(\eta_{t}^{n}(x)-\rho\right)
$$

If $\gamma=1$, Proposition 2.1 is equivalent to the statement that $\mathcal{Y}_{t}^{n, \gamma}$ converges in the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$ to $\mathcal{Y}_{t}$, the unique solution of the driftremoved Ornstein-Uhlenbeck equation

$$
\begin{equation*}
\partial_{t} \mathcal{Y}_{t}=\frac{1}{2} \varphi_{c}^{\prime}(\rho) \Delta \mathcal{Y}_{t}+\sqrt{\frac{1}{2} \varphi_{b}(\rho)} \nabla \dot{\mathcal{W}}_{t} \tag{2.10}
\end{equation*}
$$

This equation of course corresponds to (2.9) with $a=0$, is well posed and has a unique solution (cf. [61]).

Now we increase the strength of the asymmetry in the jump rates by decreasing the value of $\gamma$. We show for $1 / 2<\gamma<1$, starting from the measures $\left\{\mu^{n}\right\}$, that there is no effect in the convergence result of the fluctuation field.

THEOREM 2.2 (Crossover fluctuations). For $1 / 2<\gamma<1$, starting from initial measures $\left\{\mu^{n}\right\}$, the sequence $\left\{\mathcal{Y}_{t}^{n, \gamma} ; n \geq 1\right\}$ converges in the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$ to the process $\mathcal{Y}_{t}$ which is the solution of the OrnsteinUhlenbeck equation (2.10) with initial condition $\mathcal{Y}_{0}=\overline{\mathcal{Y}}_{0}$ given in (CLT).

However, for $\gamma=1 / 2$, which is a threshold, a much different qualitative limit behavior is obtained as the strength of the weak asymmetry in the jump rates is big enough to influence the limit field. As mentioned in the Introduction, the limit field $\mathcal{Y}_{t}$ should satisfy, in some sense, a stochastic Burgers equation, written in our framework as

$$
\begin{equation*}
\partial_{t} \mathcal{Y}_{t}=\frac{\varphi_{c}^{\prime}(\rho)}{2} \Delta \mathcal{Y}_{t}+\frac{a}{2} \varphi_{b}^{\prime \prime}(\rho) \nabla \mathcal{Y}_{t}^{2}+\sqrt{\frac{1}{2} \varphi_{b}(\rho)} \nabla \dot{\mathcal{W}}_{t} \tag{2.11}
\end{equation*}
$$

although it is ill-posed.
We now detail in what sense we mean to "solve" (2.11) in terms of a martingale formulation. Let $\iota: \mathbb{R} \rightarrow[0, \infty)$ be the function $\iota(z)=(1 / 2) 1_{[-1,1]}(z)$. Also, for $0<\varepsilon \leq 1$, define $\iota_{\varepsilon}(z)=\varepsilon^{-1} \iota\left(\varepsilon^{-1} z\right)$ and let $G_{\varepsilon}: \mathbb{R} \rightarrow[0, \infty)$ be a smooth compactly supported function in $\mathbb{S}(\mathbb{R})$ which approximates $l_{\varepsilon}$ : That is, $\left\|G_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2} \leq 2\left\|\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2}=\varepsilon^{-1}$ and

$$
\lim _{\varepsilon \downarrow 0} \varepsilon^{-1 / 2}\left\|G_{\varepsilon}-\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}=0
$$

Such choices can be readily found by convoluting $\iota_{\varepsilon}$ with smooth kernels. Also, for $x \in \mathbb{R}$, define the shift $\tau_{x}$ so that $\tau_{x} G_{\varepsilon}(z)=G_{\varepsilon}(x+z)$.

Consider now an $\mathbb{S}^{\prime}(\mathbb{R})$-valued process $\left\{\mathcal{Y}_{t} ; t \in[0, T]\right\}$ and for $0 \leq s \leq t \leq T$ let

$$
\mathcal{A}_{s, t}^{\varepsilon}(H)=\int_{s}^{t} \int_{\mathbb{R}} \nabla H(x)\left[\mathcal{Y}_{u}\left(\tau_{-x} G_{\varepsilon}\right)\right]^{2} d x d u
$$

We say the process $\mathcal{Y}$. satisfies the probability energy condition if for each $H \in$ $\mathbb{S}(\mathbb{R})$,

$$
\begin{equation*}
\left\{\mathcal{A}_{s, t}^{\varepsilon}(H)\right\} \text { is Cauchy in probability as } \varepsilon \downarrow 0 \tag{2.12}
\end{equation*}
$$

and the limit in probability does not depend on the particular smoothing family $\left\{G_{\varepsilon}\right\}$. This limit defines the process $\left\{\mathcal{A}_{s, t} ; 0 \leq s \leq t \leq T\right\}$ given by

$$
\mathcal{A}_{s, t}(H):=\lim _{\varepsilon \downarrow 0} \mathcal{A}_{s, t}^{\varepsilon}(H),
$$

which is $\mathbb{S}^{\prime}(\mathbb{R})$ valued (cf. pages 364-365; Theorem 6.15 of [61]).
We will say that $\left\{\mathcal{Y}_{t} ; t \in[0, T]\right\}$ is a probability energy solution of $(2.11)$ if the following conditions hold:
(i) Initially, $\mathcal{Y}_{0}$ is a spatial Gaussian process with covariance $\mathcal{C}(G, H)$ for $G, H \in \mathbb{S}(\mathbb{R})$.
(ii) The process $\left\{\mathcal{Y}_{t} ; t \in[0, T]\right\}$ satisfies the probability energy condition (2.12).
(iii) Then, the $\mathbb{S}^{\prime}(\mathbb{R})$ valued process $\left\{\mathcal{M}_{t}: t \in[0, T]\right\}$ where

$$
\begin{equation*}
\mathcal{M}_{t}(H):=\mathcal{Y}_{t}(H)-\mathcal{Y}_{0}(H)-\frac{\varphi_{c}^{\prime}(\rho)}{2} \int_{0}^{t} \mathcal{Y}_{s}(\Delta H) d s-\frac{a \varphi_{b}^{\prime \prime}(\rho)}{2} \mathcal{A}_{0, t}(H) \tag{2.13}
\end{equation*}
$$

is a continuous martingale with quadratic variation

$$
\left\langle\mathcal{M}_{t}(H)\right\rangle=\frac{\varphi_{b}(\rho) t}{2}\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}
$$

In particular, condition (iii) specifies by Lévy's theorem that $\mathcal{M}_{t}(H)$ is a Brownian motion with variance $\left(\varphi_{b}(\rho) / 2\right) t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}$.

We also define a stronger notion of solution to (2.11) which may be verified in some cases. We say that $\mathcal{Y}_{t}$ satisfies the $L^{2}$ energy condition if in (2.12), instead of in the probability sense, we assert $\left\{\mathcal{A}_{s, t}^{\varepsilon}(H)\right\}$ is Cauchy in $L^{2}$ with respect to the underlying probability measure, and $\mathcal{A}_{s, t}(H)$ is its $L^{2}$ limit. Then we say $\mathcal{Y}_{t}$ is an $L^{2}$ energy solution of (2.11) if (i) holds as before, (ii) the $L^{2}$ energy condition holds and (iii) holds with respect to the $L^{2}$ limit $\mathcal{A}_{s, t}(H)$.

THEOREM 2.3 (KPZ fluctuations). For $\gamma=1 / 2$, starting from initial measures $\left\{\mu^{n}\right\}$, the sequence of processes $\left\{\mathcal{Y}_{t}^{n, \gamma}: n \geq n_{0}\right\}$ is tight in the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$. Moreover, any limit point of $\mathcal{Y}_{t}^{n, \gamma}$ is a probability energy solution with respect to $(2.11)$ with initial field $\overline{\mathcal{Y}}_{0}$ given in (CLT).

If the initial measure is $\underline{\mu}^{n} \equiv v_{\rho}$, any limit point of $\mathcal{Y}_{t}^{n, \gamma}$ is an $L^{2}$ energy solution of (2.11) with initial field $\overline{\mathcal{Y}}_{0}$ given in (CLT).

REMARK 2.4. We now make the following comments:

1. Formally, equation (2.13) corresponds to the stochastic Burgers equation (2.11) where the nonlinear term is represented by $\mathcal{A}_{0, t}$. We remark, as in [3], by taking a fast subsequence in $\varepsilon$, one may write $\mathcal{A}_{0, t}$ as a function of $\left\{\mathcal{Y}_{u}: u \leq t\right\}$, and form an equation in which $\mathcal{Y}_{t}$ satisfies (2.11) a.s. on a type of negative order Hermite Hilbert space.
2. We also remark, as alluded to in the Introduction, if there were a unique probability or $L^{2}$ energy solution, that is uniqueness of process in the associated "martingale formulation," since with respect to simple exclusion the fluctuation field limit is known in terms of the "Cole-Hopf" solution of the KPZ equation [14], not only could one conclude a unique fluctuation field limit in Theorem 2.3 in the
framework of the particle systems considered, but also identify it in terms of the "Cole-Hopf" apparatus. What is required to show uniqueness of $\mathcal{Y}_{t}$ is to determine uniquely its finite dimensional distributions (cf. Section 4.4 of [21]), which the nonlinearity of $\mathcal{A}_{0, t}$ makes difficult.
3. We also note that the statement of Theorem 2.3 is nontrivial when $a \neq 0$ and $b$ is such that

$$
\begin{equation*}
\varphi_{b}^{\prime \prime}(\rho) \neq 0 \tag{2.14}
\end{equation*}
$$

Otherwise, when $\varphi_{b}^{\prime \prime}(\rho)=0$, the limit field $\mathcal{Y}_{t}$ satisfies the Ornstein-Uhlenbeck equation (2.10). Examples, fitting in our framework, where the second derivative vanishes include types of zero-range, that is independent particle systems where $\varphi_{b}(\rho)=2 \rho$ which are in the EW class.
2.3. Model 1: Zero-range processes. The one-dimensional weakly asymmetric zero-range process $\eta_{t}^{n}$, on the state space $\Omega:=\mathbb{N}_{0}^{\mathbb{Z}}$, consists of a collection of random walks which interact in that the jump rate of a particle at vertex $x$ only depends on the number of particles at $x$. More precisely, the generator is in form (2.1) where

$$
b_{x}^{R, n}(\eta)=g(\eta(x)) \quad \text { and } \quad b_{x}^{L, n}(\eta)=g(\eta(x+1))
$$

do not depend on $n$ and are fixed with respect to a function $g: \mathbb{N}_{0} \rightarrow \mathbb{R}_{+}$such that $g(0)=0, g(k)>0$ for $k \geq 1$ and $g$ is Lipschitz,
(LIP) $\sup _{k \geq 0}|g(k+1)-g(k)|<\infty$.
Under this specification, a Markov process $\eta_{t}^{n}$ can be constructed (on a subset of $\Omega$ ) [2]. Hence, (R1) holds and we identify the fixed function $c^{n} \equiv c$ as

$$
c(\eta)=g(\eta(0))
$$

The zero-range process possesses a family of invariant measures which are fairly explicit product measures. For $\alpha \geq 0$, define

$$
\mathcal{Z}(\alpha):=\sum_{k \geq 0} \frac{\alpha^{k}}{g(k)!},
$$

where $g(k)!=g(1) \cdots g(k)$ for $k \geq 1$ and $g(0)!=1$. Let $\alpha^{*}$ be the radius of convergence of this power series and notice that $\mathcal{Z}$ increases on [0, $\alpha^{*}$ ). Fix $0 \leq \alpha<\alpha^{*}$ and let $\bar{v}_{\alpha}$ be the product measure on $\mathbb{N}^{\mathbb{Z}}$ whose marginal at the site $x$ is given by

$$
\bar{v}_{\alpha}\{\eta: \eta(x)=k\}= \begin{cases}\frac{1}{\mathcal{Z}(\alpha)} \frac{\alpha^{k}}{g(k)!}, & \text { when } k \geq 1 \\ \frac{1}{\mathcal{Z}(\alpha)}, & \text { when } k=0\end{cases}
$$

We now reparameterize these measures in terms of the "density." Let $\rho(\alpha):=$ $E_{\bar{\nu}_{\alpha}}[\eta(0)]=\alpha \mathcal{Z}^{\prime}(\alpha) / \mathcal{Z}(\alpha)$. By computing the derivative, we obtain that $\rho(\alpha)$ is strictly increasing on $\left[0, \alpha^{*}\right)$. Then let $\alpha(\cdot)$ denote its inverse. Define

$$
v_{\rho}(\cdot):=\bar{v}_{\alpha(\rho)}(\cdot),
$$

so that $\left\{v_{\rho}: 0 \leq \rho<\rho^{*}\right\}$ is a family of invariant measures parameterized by the density. Here, $\rho^{*}=\lim _{\alpha \uparrow \alpha^{*}} \rho(\alpha)$, which may be finite or infinite depending on whether $\lim _{\alpha \rightarrow \alpha^{*}} \mathcal{Z}(\alpha)$ converges or diverges.

Note, since $\nu_{\rho}$ is a product measure, that $\nu_{\rho}^{\lambda(z)}=v_{z}$ for $0 \leq z<\rho^{*}$, and condition (D) holds. One can readily check that (R2) holds:

$$
\begin{aligned}
g\left(\eta^{x+1, x}(x)\right) \frac{d v_{\rho}^{x+1, x}}{d v_{\rho}} & =g(\eta(x)+1) \frac{g(\eta(x))!g(\eta(x+1))!}{g(\eta(x)+1)!g(\eta(x+1)-1)!} \\
& =g(\eta(x+1))
\end{aligned}
$$

Also, by the construction in [55], which extends the construction in [2] to an $L^{2}\left(v_{\rho}\right)$ process, we have that $L_{n}$ is a Markov $L^{2}\left(v_{\rho}\right)$ generator whose core can be taken as the space of all local $L^{2}\left(v_{\rho}\right)$ functions. Indeed, in [55], a core of bounded Lipschitz functions is identified; however, since any local $L^{2}\left(v_{\rho}\right)$ function is a limit of bounded Lipschitz functions, and the formula (2.1) is well defined and $L^{2}\left(v_{\rho}\right)$-bounded for a local $L^{2}\left(v_{\rho}\right)$ function, by dominated convergence the core can be extended. It follows that the measures $\left\{v_{\rho}: 0 \leq \rho<\rho^{*}\right\}$ are invariant for the zero-range process. Also, (IM) holds as $v_{\rho}$ is a product measure whose marginal has some exponential moments. In addition, one can check that (EE) holds by Proposition 5.1.

We now address the spectral gap properties of the system. Since the model interactions are range 0 , the gap does not depend on the outside variables $\xi$. However, the gap depends on $g$, as it should since $g$ controls the rate of jumps. We identify three types of rates for which a spectral gap bound has been proved. Let $\beta=k /(2 \ell+1)^{d}$.

- If $g$ is not too different from the independent case, for which the gap is of order $O\left(\ell^{-2}\right)$ uniform in $k$, one expects similar behavior as for a single particle. This has been proved for $d \geq 1$ in [38] under assumptions (LIP) and
(U) There exists $x_{0}$ and $\varepsilon_{0}>0$ such that $g\left(x+x_{0}\right)-g(x) \geq \varepsilon_{0}$ for all $x \geq 0$.
- If $g$ is sublinear, that is $C^{-1} x^{\gamma} \leq g(x+1)-g(x) \leq C x^{\gamma}$ for $0<\gamma<1$ and $C>0$, then it has been shown that the spectral gap depends on the number of particles $k$, namely the gap for $d \geq 1$ is $O\left((1+\beta)^{-\gamma} \ell^{-2}\right)$ [45].
- If $g(x)=1(x \geq 1)$, then it has been shown in $d \geq 1$ that the gap is $O((1+$ $\beta)^{-2} \ell^{-2}$ ) [43]. In $d=1$, this is true because of the connection between the zero-range and simple exclusion processes for which the gap estimate is well known [47]: The number of spaces between consecutive particles in simple exclusion correspond to the number of particles in the zero-range process.

In all these cases, (G) follows readily by straightforward moment calculations.
2.4. Model 2: Kinetically constrained exclusion systems. We consider a type of exclusion process, which may be thought of as a microscopic model for porous medium behavior, developed in [29] and references therein, in one dimension on $\Omega=\{0,1\}^{\mathbb{Z}}$ where particles more likely hop to unoccupied nearest-neighbor sites when at least $m-1 \geq 1$ other neighboring sites are full. When $m=2$, the rates are in the form

$$
\begin{aligned}
b_{x}^{R, n}(\eta ; \theta) & =\eta(x)(1-\eta(x+1))\left[\eta(x-1)+\eta(x+2)+\frac{\theta}{2 n}\right], \\
b_{x}^{L, n}(\eta ; \theta) & =\eta(x+1)(1-\eta(x))\left[\eta(x-1)+\eta(x+2)+\frac{\theta}{2 n}\right],
\end{aligned}
$$

with respect to a parameter $\theta>0$. If $\theta$ would vanish, particles can jump from site $x$ to $x+1$ exactly when there is at least 1 particle in the vicinity of the bond $(x, x+$ 1). However, with $\theta>0$, the jump from $x$ to $x+1$ may also occur irrespective of the neighboring particle structure with a small rate $\theta /(2 n)$.

When $m \geq 2$, the rates generalize to

$$
\begin{aligned}
& b_{x}^{R, n}(\eta ; \theta)=\eta(x)(1-\eta(x+1)) A_{n}(\eta ; \theta), \\
& b_{x}^{L, n}(\eta ; \theta)=\eta(x+1)(1-\eta(x)) A_{n}(\eta ; \theta),
\end{aligned}
$$

where $A_{n}(\eta ; \theta)$ equals
$\prod_{j=-(m-1)}^{-1} \eta(x+j)+\prod_{\substack{j=-(m-2) \\ j \neq 0,1}}^{2} \eta(x+j)+\cdots+\prod_{\substack{j=-1 \\ j \neq 0,1}}^{m-1} \eta(x+j)+\prod_{j=2}^{m} \eta(x+j)+\frac{\theta}{2 n}$.
The role of $\theta>0$ is to make the system "ergodic." If $\theta=0$, there would be an infinite number of invariant measures, such as Dirac measures supported on configurations which cannot evolve under the dynamics. The hydrodynamic limit for this model corresponds to the porous medium equation, $\partial_{t} \rho_{t}(t, u)=\Delta \rho^{m}(t, u)$.

Now, one may calculate that $b_{x}^{R, n}(\eta ; \theta)-b_{x}^{L, n}(\eta ; \theta)=c_{x}^{n}(\eta)-c_{x+1}^{n}(\eta)$ where, for $m \geq 2$,

$$
\begin{aligned}
c^{n}(\eta ; \theta)= & \prod_{j=-(m-1)}^{0} \eta(j)+\cdots+\prod_{j=0}^{m-1} \eta(j) \\
& -\prod_{\substack{j=-(m-1) \\
j \neq 0}}^{1} \eta(j)-\cdots-\prod_{\substack{j=-1 \\
j \neq 0}}^{m-1} \eta(j)+\frac{\theta}{2 n} \eta(0) .
\end{aligned}
$$

In the case $m=2$, the last formula reduces to $c^{n}(\eta ; \theta)=\eta(-1) \eta(0)+\eta(0) \eta(1)-$ $\eta(-1) \eta(1)+\frac{\theta}{2 n} \eta(0)$.

Of course, uniformly in $\eta$, as $n \uparrow \infty$, the terms involving $\theta$ vanish,

$$
\begin{aligned}
b_{x}^{R, n}(\eta ; \theta) & \rightarrow b_{x}^{R}(\eta):=b_{x}^{R, 1}(\eta ; 0), \quad b_{x}^{L, n}(\eta ; \theta) \rightarrow b_{x}^{L}(\eta):=b_{x}^{L, 1}(\eta ; 0) \quad \text { and } \\
c^{n}(\eta ; \theta) & \rightarrow c:=c^{1}(\eta ; 0)
\end{aligned}
$$

Consider now the Bernoulli product measure on $\Omega$ :

$$
v_{\rho}=\prod_{x \in \mathbb{Z}} \mu_{\rho} \quad \text { where } \mu_{\rho}(1)=1-\mu_{\rho}(0)=\rho
$$

for $\rho \in[0,1]$. By the construction in [40], it is now standard that $L_{n}$ is a Markov $L^{2}\left(v_{\rho}\right)$ generator. One may also inspect that condition (R2) holds with respect to $v_{\rho}$. Hence, $v_{\rho}$ is invariant for $\rho \in[0,1]$. Condition (IM) also holds as $v_{\rho}$ supports two-state configurations. In addition, as $v_{\rho}$ is a product measure, $\nu_{\rho}^{\lambda(z)}=v_{z}$ and (D) holds. Also, by Proposition 5.1, (EE) is satisfied.

We now discuss the spectral gap behavior of the process.
Proposition 2.5. For kinetically constrained exclusion processes evolving on $\Lambda_{\ell}$, when $m \geq 2$, there exists a constant $C$, uniform over $\xi$ and $n$, such that

$$
W(k, \ell, \xi, n) \leq C \ell^{2}\left(\frac{\ell}{k}\right)^{m} 1(k \geq 1)
$$

When $m=2$ and $k \leq \ell / 3$, the above spectral gap estimate is already given in Proposition 6.2 of [29]. However, a straightforward modification of the proof of Proposition 6.2 in [29] yields the more general estimate in Proposition 2.5. Indeed, the difference when $m \geq 2$ is that to bound equation (6.10) in [29] in the general case, one uses that there are at most $C j^{m-1}$ ways to arrange $m-1$ particles in an interval of width $j$. Now, a similar optimization on $j$ as given in the proof of Proposition 6.2 of [29] leads to the desired generalized spectral gap estimate.

LEMMA 2.6. For the kinetically constrained exclusion model, the spectral gap condition (G) is satisfied.

Proof. With respect to a constant $C$, which may change line to line,

$$
\begin{aligned}
& E_{v_{\rho}}\left[\left(W\left(\sum_{x \in \Lambda_{\ell}} \eta(x), \ell, \xi, n\right)\right)^{2}\right] \\
& \quad \leq C \ell^{4} E_{v_{\rho}}\left[1\left(\frac{1}{2 \ell+1} \leq \eta^{(\ell)}\right)\left(\eta^{\ell}\right)^{-2 m}\right] \\
& \quad \leq C \ell^{4}\left\{\varepsilon^{-2 m}+E_{v_{\rho}}\left[1\left(\frac{1}{2 \ell+1} \leq \eta^{(\ell)}<\varepsilon\right)\left(\eta^{(\ell)}\right)^{-2 m}\right]\right\} \\
& \quad \leq C \ell^{4}\left\{\varepsilon^{-2 m}+\ell^{2 m} P_{v_{\rho}}\left(\eta^{(\ell)}<\varepsilon\right)\right\}
\end{aligned}
$$

for a fixed $\varepsilon<\rho$. Then, as $v_{\rho}$ is a Bernoulli product measure with density $\rho$, by a large deviations estimate say, $E_{v_{\rho}}\left[W\left(\sum_{x \in \Lambda_{\ell}} \eta(x), \ell, \xi, n\right)^{2}\right] \leq C \ell^{4}$ for all $\ell \geq 1$.
2.5. Model 3: Gradient exclusion with speed change. In this version of exclusion on $\Omega=\{0,1\}^{\mathbb{Z}}$, rates are chosen which correspond to a Hamiltonian with nearest-neighbor interactions,

$$
Q_{\beta}(\eta)=-\beta \sum_{x \in \mathbb{Z}}(\eta(x)-1 / 2)(\eta(x+1)-1 / 2)
$$

for $\beta \in \mathbb{R}$, which will be reversible with respect to a stationary Markovian measure $v_{1 / 2}$. That is, specify $\nu_{1 / 2}$ by its finite-dimensional distributions

$$
\nu_{1 / 2}\left(\eta(x)=e(x): x \in \Lambda_{\ell} \mid \eta(y)=\xi(y) \text { for } y \notin \Lambda_{\ell}\right)=\frac{e^{-Q_{\beta, \ell}(e, \xi)}}{\mathcal{Z}}
$$

where

$$
\begin{aligned}
Q_{\beta, \ell}(e, \xi)= & -\beta \sum_{x, x+1 \in \Lambda_{\ell}}(e(x)-1 / 2)(e(x+1)-1 / 2) \\
& -\beta(\xi(-\ell-1)-1 / 2)(e(-\ell)-1 / 2) \\
& -\beta(e(\ell)-1 / 2)(\xi(\ell+1)-1 / 2)
\end{aligned}
$$

$e, \xi \in \Omega$ and $\mathcal{Z}=\mathcal{Z}(\ell, \xi)$ is the normalization. It is not difficult to see that $\nu_{1 / 2}$ is Markovian with transition matrix

$$
P=\frac{1}{e^{\beta / 4}+e^{-\beta / 4}}\left[\begin{array}{cc}
e^{\beta / 4} & e^{-\beta / 4} \\
e^{-\beta / 4} & e^{\beta / 4}
\end{array}\right]
$$

and marginal distribution $\langle 1 / 2,1 / 2\rangle$ so that $E_{v_{1 / 2}}[\eta(0)]=1 / 2$.
Then, as discussed in [57], Section II.2.4, (R2) is ensured if we take the rates $b_{x}^{R, n}=b_{x}^{R}$ and $b_{x}^{L, n}=b_{x}^{L}$ which do not depend on $n$ as

$$
\begin{aligned}
b_{x}^{R}(\eta)= & \eta(x)(1-\eta(x+1)) \\
& \times\left[\alpha_{1} \eta(x-1) \eta(x+2)+\alpha_{2}(1-\eta(x-1)) \eta(x+2)\right. \\
& \left.\quad+\alpha_{3} \eta(x-1)(1-\eta(x+2))+\alpha_{4}(1-\eta(x-1))(1-\eta(x+2))\right] \\
b_{x}^{L}(\eta)= & \eta(x+1)(1-\eta(x)) \\
\times & {\left[\alpha_{1} \eta(x-1) \eta(x+2)+\alpha_{3}(1-\eta(x-1)) \eta(x+2)\right.} \\
& \left.+\alpha_{2} \eta(x-1)(1-\eta(x+2))+\alpha_{4}(1-\eta(x-1))(1-\eta(x+2))\right]
\end{aligned}
$$

where $\alpha_{1}, \alpha_{2}=e^{\beta} \alpha_{3}, \alpha_{4}>0$. The condition (R1) also follows if we also assume that $\alpha_{1}-\alpha_{2}-\alpha_{3}+\alpha_{4}=0$ so that, as can be checked, $c(\eta)$ takes the form

$$
\begin{aligned}
c(\eta)= & \alpha_{4} \eta(0)+\left(\alpha_{3}-\alpha_{4}\right) \eta(-1) \eta(0)+\left(\alpha_{3}-\alpha_{4}\right) \eta(0) \eta(1) \\
& +\left(\alpha_{4}-\alpha_{2}\right) \eta(-1) \eta(1)+\left(\alpha_{2}-\alpha_{3}\right) \eta(-1) \eta(0) \eta(1) .
\end{aligned}
$$

Again, by [40], $L_{n}$ is a Markov $L^{2}\left(v_{\rho}\right)$ generator for the process. We note when $\beta=0$ and $\alpha_{i}=1$ for $i=1,2,3,4$, the model is the simple exclusion process and $\nu_{1 / 2}$ is the Bernoulli product measure with density $1 / 2$.

We now introduce a family of stationary, reversible measures by use of a "tilt" or "chemical potential" $\lambda$. Define $\nu_{1 / 2}^{\lambda}$, again specified by its finite-dimensional distributions, through the relation

$$
\frac{d \nu_{1 / 2}^{\lambda}}{d \nu_{1 / 2}}\left(\eta(x)=e(x): x \in \Lambda_{\ell} \mid \eta(y)=\xi(y) \text { for } y \notin \Lambda_{\ell}\right)=\frac{e^{\lambda \sum_{x \in \Lambda_{\ell}}(e(x)-1 / 2)}}{\mathcal{Z}^{\prime}}
$$

where $e, \xi \in \Omega$ and $\mathcal{Z}^{\prime}=\mathcal{Z}^{\prime}(\ell, \xi)$ is another normalization. These measures are also Markovian with transition matrix

$$
P_{\lambda}=\left[\begin{array}{cc}
1-u_{1} & u_{1}  \tag{2.15}\\
v_{1} & 1-v_{1}
\end{array}\right],
$$

where

$$
u_{1}=\frac{r_{1}(\lambda, \beta)+\sinh (\lambda / 2)}{\cosh (\lambda / 2)+r_{1}(\lambda, \beta)}, \quad v_{1}=\frac{r_{1}(\lambda, \beta)-\sinh (\lambda / 2)}{\cosh (\lambda / 2)+r_{1}(\lambda, \beta)}
$$

and $r_{1}(\lambda, \beta)=\sqrt{\sinh ^{2}(\lambda / 2)+e^{-\beta}}$. The stationary distribution equals $\pi_{\lambda}=\left(v_{1}+\right.$ $\left.u_{1}\right)^{-1}\left\langle v_{1}, u_{1}\right\rangle$, which is the marginal distribution of $v_{1 / 2}^{\lambda}$. These derivations are performed in [12] and [62].

The measures $\left\{\nu_{1 / 2}^{\lambda}: \lambda \in \mathbb{R}\right\}$ are uniformly mixing: Indeed, the eigenvalues of $P_{\lambda}$ are 1 and $1-u_{1}-v_{1}$, and the spectral gap $u_{1}+v_{1}$ is uniformly bounded away from 0 for $\lambda \in \mathbb{R}$.

One can calculate $E_{\nu_{1 / 2}^{\lambda}}[\eta(0)]=u_{1} /\left(u_{1}+v_{1}\right)$ strictly increases in $\lambda$. To parameterize in terms of "density", recall $\nu_{1 / 2}^{\lambda(z)}=v_{z}$ where $\lambda=\lambda(z)$ is chosen so that $E_{\nu_{z}}[\eta(0)]=z$. Here, as $z \downarrow 0=\rho_{*}, \lambda(z) \downarrow-\infty$ and, as $z \uparrow 1=\rho^{*}, \lambda(z) \uparrow \infty$; also $\lambda(1 / 2)=0$. Hence, since also $v_{z}$ is exponentially mixing, both (IM) and (D) hold.

From the defining relation for $\lambda(z), u_{1} /\left(u_{1}+v_{1}\right)=z$, one can differentiate at $z=1 / 2$ to find $\lambda^{\prime}(1 / 2)\left[e^{-\beta / 2} / 4\right]=1$.

Also, we note the additive functional variance $\sigma^{2}(z)$ [cf. (IM2)] satisfies the formula $\sigma^{2}(z)=E_{\pi_{\lambda(z)}}\left[u^{2}\right]-E_{\pi_{\lambda(z)}}\left[\left(P_{\lambda(z)} u\right)^{2}\right]$ where $\left(I-P_{\lambda(z)}\right) u=f$ and $f=$ $\langle-z, 1-z\rangle$ represents the values of the function $f(\eta)=\eta(0)-z$; see Section 6.5 of [60]. In fact, we find $\sigma^{2}(1 / 2)=e^{-\beta / 2} / 4$ and so $\lambda^{\prime}(1 / 2) \sigma^{2}(1 / 2)=1$.

The spectral gap for a more general process, including this one, has been bounded as follows [41]: Uniformly over $k$ and $\xi$ (it does not depend on $n$ ), we have

$$
W(k, \ell, \xi, n) \leq C \ell^{2}
$$

Hence, (G) holds.
Also, in Proposition 5.2, we show that (EE) holds.
3. Proofs-outline. The strategies of the proofs for Theorems 2.2 and 2.3 are similar. We consider the stochastic differential of $\mathcal{Y}_{t}^{n, \gamma}$ and represent it in terms of corrector and martingale terms. Tightness is shown for each term in the decomposition of $\mathcal{Y}_{t}^{n, \gamma}$. Under the assumption that the initial measure is the invariant state $v_{\rho}$, limit points are identified using a Boltzmann-Gibbs principle, and shown to satisfy (2.10) when $1 / 2<\gamma \leq 1$ and to be energy solutions of (2.11) when $\gamma=1 / 2$. When the initial measures $\left\{\mu^{n}\right\}$ satisfy (BE), the entropy inequality then allows to characterize the limit points as desired.

In the following Sections 3.1-3.3, associated martingales, Boltzmann-Gibbs principles and tightness are discussed. In Section 3.4, limit points are identified and Theorems 2.2 and 2.3 are proved.

To reduce some of the notation, we will drop the superscript " $n$ " in the rate functions and write $b_{x}^{R, n}=b_{x}^{R}, b_{x}^{L, n}=b_{x}^{L}, b_{x}^{n}=b_{x}, b^{n}=b, c_{x}^{n}=c_{x}$ and $c^{n}=c$ until Section 3.4.
3.1. Associated martingales. For $H \in \mathbb{S}(\mathbb{R}), x \in \mathbb{Z}$ and $n \geq 1$, define

$$
\begin{aligned}
\Delta_{x}^{n} H & =n^{2}\left\{H\left(\frac{x+1}{n}\right)+H\left(\frac{x-1}{n}\right)-2 H\left(\frac{x}{n}\right)\right\} \\
\nabla_{x}^{n} H & =n\left\{H\left(\frac{x+1}{n}\right)-H\left(\frac{x}{n}\right)\right\}
\end{aligned}
$$

Define also, for $\gamma, s \geq 0$, the functions

$$
\begin{align*}
& H_{\gamma, s}(\cdot)=H\left(\cdot-\frac{1}{n}\left\lfloor\frac{a \varphi_{b}^{\prime}(\rho) s n^{2}}{2 n^{\gamma}}\right\rfloor\right) \text { and }  \tag{3.1}\\
& \widetilde{H}_{\gamma, s}(\cdot)=H\left(\cdot-\frac{1}{n}\left\{\frac{a \varphi_{b}^{\prime}(\rho) s n^{2}}{2 n^{\gamma}}\right\}\right) .
\end{align*}
$$

We note, in $H_{\gamma, s}$, the process characteristic shift is along $n^{-1} \mathbb{Z}$, which helps make tidy some proofs [in applying a Boltzmann-Gibbs principle (Theorem 3.2) in proofs of Propositions 3.3 and 3.5], instead of along $\mathbb{R}$ as in $\widetilde{H}_{\gamma, s}$.

Let $F\left(s, \eta_{s}^{n} ; H, n\right)=\mathcal{Y}_{s}^{n, \gamma}(H)$, and $F(\eta ; H, n)=n^{-1 / 2} \sum_{x \in \mathbb{Z}} H\left(\frac{x}{n}\right)(\eta(x)-\rho)$. Although $F(\eta ; H, n)$ is an $L^{2}\left(v_{\rho}\right)$ function, in general, it is not a local function. However, by approximation with local functions and noting by condition (R1) that $|b(\eta)| \leq C \sum_{|x| \leq R} \eta(x)$, one may conclude $F(\eta ; H, n)$ and also $F^{2}(\eta ; H, n)$ belong to the domain of $L_{n}$. In particular,

$$
L_{n} F\left(s, \eta_{s}^{n} ; H, n\right)=\frac{1}{2 \sqrt{n}} \sum_{x \in \mathbb{Z}} c_{x}\left(\eta_{s}^{n}\right) \Delta_{x}^{n} \widetilde{H}_{\gamma, s}+\frac{a}{2 n^{\gamma-1 / 2}} \sum_{x \in \mathbb{Z}} b_{x}\left(\eta_{s}^{n}\right) \nabla_{x}^{n} \widetilde{H}_{\gamma, s}
$$

Also,

$$
\frac{\partial}{\partial_{s}} F\left(s, \eta_{s}^{n} ; H, n\right)=\left\{\frac{-a \varphi_{b}^{\prime}(\rho) n^{2}}{2 n^{\gamma}}\right\} \frac{1}{n^{3 / 2}} \sum_{x \in \mathbb{Z}} \nabla \tilde{H}_{\gamma, s}\left(\frac{x}{n}\right)\left(\eta_{s}^{n}(x)-\rho\right)
$$

Then

$$
\begin{aligned}
\mathcal{M}_{t}^{n, \gamma}(H):= & F\left(t, \eta_{t}^{n} ; H, n\right)-F\left(0, \eta_{0}^{n} ; H, n\right) \\
& -\int_{0}^{t} \frac{\partial}{\partial_{s}} F\left(s, \eta_{s}^{n} ; H, n\right)+L_{n} F\left(s, \eta_{s}^{n} ; H, n\right) d s
\end{aligned}
$$

is a martingale. We may decompose
(3.2) $\mathcal{M}_{t}^{n, \gamma}(H)=\mathcal{Y}_{t}^{n, \gamma}(H)-\mathcal{Y}_{0}^{n, \gamma}(H)-\mathcal{I}_{t}^{n, \gamma}(H)-\mathcal{B}_{t}^{n, \gamma}(H)-\mathcal{K}_{t}^{n, \gamma}(H)$,
where

$$
\begin{aligned}
& \mathcal{I}_{t}^{n, \gamma}(H)= \frac{1}{2} \int_{0}^{t} \frac{1}{\sqrt{n}} \sum_{x \in \mathbb{Z}}\left(c_{x}\left(\eta_{s}^{n}\right)-\varphi_{c}(\rho)\right) \Delta_{x}^{n} H_{\gamma, s} d s \\
& \mathcal{B}_{t}^{n, \gamma}(H)= \frac{a}{2 n^{\gamma-1 / 2}} \int_{0}^{t} \sum_{x \in \mathbb{Z}}\left(b_{x}\left(\eta_{s}^{n}\right)-\varphi_{b}(\rho)-\varphi_{b}^{\prime}(\rho)\left(\eta_{s}^{n}(x)-\rho\right)\right) \nabla_{x}^{n} H_{\gamma, s} d s \\
& \mathcal{K}_{t}^{n, \gamma}(H) \\
&= \int_{0}^{t}\left[\frac{1}{\sqrt{n}} \sum_{x \in \mathbb{Z}} \kappa_{x}^{n, 1}(H, s)\left(c_{x}\left(\eta_{s}^{n}\right)-\varphi_{c}(\rho)\right)\right. \\
&\left.+\frac{a}{2 n^{\gamma-1 / 2}} \sum_{x \in \mathbb{Z}} \kappa_{x}^{n, 2}(H, s)\left(b_{x}\left(\eta_{s}^{n}\right)-\varphi_{b}(\rho)-\varphi_{b}^{\prime}(\rho)\left(\eta_{s}^{n}(x)-\rho\right)\right)\right] d s
\end{aligned}
$$

Here, we introduced the centering constants $\varphi_{c}(\rho)$ and $\varphi_{b}(\rho)$ in $\mathcal{I}_{t}^{n, \gamma}$ and $\mathcal{B}_{t}^{n, \gamma}$ as $\Delta_{x}^{n} H_{\gamma, s}$ and $\nabla_{x}^{n} H_{\gamma, s}$ both sum to zero. Also,

$$
\begin{aligned}
\kappa_{x}^{n, 1}(H, s) & =\Delta_{x}^{n}\left(\tilde{H}_{\gamma, s}-H_{\gamma, s}\right)=O\left(n^{-1}\right) \cdot \Delta_{x}^{n} H_{\gamma, s}^{\prime}+O\left(n^{-2}\right) \cdot H_{\gamma, s}^{(4)}\left(x^{\prime} / n\right) \\
\kappa_{x}^{n, 2}(H, s) & =\nabla_{x}^{n}\left(\widetilde{H}_{\gamma, s}-H_{\gamma, s}\right) \\
& =O\left(n^{-1}\right) \cdot \Delta H_{\gamma, s}(x / n)+O\left(n^{-2}\right) \cdot H_{\gamma, s}^{\prime \prime \prime}\left(x^{\prime \prime} / n\right)
\end{aligned}
$$

where $\left|x^{\prime}-x\right|,\left|x^{\prime \prime}-x\right| \leq 2$.
To capture the quadratic variation $\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle$, we compute

$$
\begin{aligned}
& L_{n} F^{2}\left(s, \eta_{s}^{n} ; H, n\right)-2 F\left(s, \eta_{s}^{n} ; H, n\right) L_{n} F\left(s, \eta_{s}^{n} ; H, n\right) \\
& \quad=\frac{1}{2 n} \sum_{x \in \mathbb{Z}} b_{x}\left(\eta_{s}^{n}\right)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, s}\right)^{2}+\frac{a}{2 n^{1+\gamma}} \sum_{x \in \mathbb{Z}}\left(c_{x}\left(\eta_{s}^{n}\right)-c_{x+1}\left(\eta_{s}^{n}\right)\right)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, s}\right)^{2}
\end{aligned}
$$

so that $\left(\mathcal{M}_{t}^{n, \gamma}(H)\right)^{2}-\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle$ is a martingale with

$$
\begin{aligned}
\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle= & \int_{0}^{t} \frac{1}{2 n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, s}\right)^{2} b_{x}\left(\eta_{s}^{n}\right) d s \\
& +\int_{0}^{t} \frac{a}{2 n^{1+\gamma}} \sum_{x \in \mathbb{Z}}\left(c_{x}\left(\eta_{s}^{n}\right)-c_{x+1}\left(\eta_{s}^{n}\right)\right)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, s}\right)^{2} d s
\end{aligned}
$$

When starting from the invariant measure $v_{\rho}$, noting the bounds in (R1), we have

$$
\begin{align*}
\mathbb{E}_{v_{\rho}} & {\left[\left(\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)\right)^{2}\right] } \\
\leq & \left\{\int_{s}^{t}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, s}\right)^{2}\right) d s\right\} \\
& \times\left[\frac{1}{2} E_{v_{\rho}}[b(\eta)]+\frac{a}{2 n^{\gamma}} E_{\nu_{\rho}}\left[\left|c_{0}(\eta)-c_{1}(\eta)\right|\right]\right]  \tag{3.3}\\
\leq & C(a)\|b\|_{L^{1}\left(v_{\rho}\right)} \int_{s}^{t}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, s}\right)^{2}\right) d s .
\end{align*}
$$

To express an exponential martingale, we now observe for $0 \leq \lambda \leq \lambda(H, n)$ small that $\exp \{\lambda F(\eta ; H, n)\}$ is in the domain of $L_{n}$. Indeed, if $H$ is a local function, as $v_{\rho}$ is assumed in (IM) to have small parameter exponential moments, then $\exp \{\lambda F(\eta ; H, n)\} \in L^{2}\left(v_{\rho}\right)$ for all small $\lambda$. Again, an approximation argument when $H \in \mathbb{S}(\mathbb{R})$ is not local shows also $\exp \{\lambda F(\eta ; H, n)\}$ belongs to the domain of $L_{n}$. We calculate

$$
\begin{aligned}
& \exp \left\{-\lambda F\left(u, \eta_{u}^{n} ; H, n\right)\right\}\left(\frac{\partial}{\partial_{u}}+L_{n}\right) \exp \left\{\lambda F\left(u, \eta_{u}^{n} ; H, n\right)\right\} \\
&=n^{2} \sum_{x \in \mathbb{Z}} {\left[b_{x}^{R}(\eta) p_{n}\left(\exp \left\{\lambda n^{-3 / 2}\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right)\right\}-1\right)\right.} \\
&\left.+b_{x}^{L}(\eta) q_{n}\left(\exp \left\{-\lambda n^{-3 / 2}\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right)\right\}-1\right)\right] \\
&- \frac{1}{n^{3 / 2}}\left\{\frac{a \lambda \varphi_{b}^{\prime}(\rho) n^{2}}{2 n^{\gamma}}\right\} \sum_{x \in \mathbb{Z}} \nabla \widetilde{H}_{\gamma, u}(x / n)\left(\eta_{u}^{n}(x)-\rho\right),
\end{aligned}
$$

which, given the assumptions on $b$ in (R1) and on moments of $v_{\rho}$ in (IM), belongs to $L^{2}\left(v_{\rho}\right)$.

Hence, by the proof of Lemma IV.3.2 of [21],

$$
\mathcal{Z}_{s, t}=\exp \left\{\lambda F\left(t, \eta_{t}^{n}\right)-\lambda F\left(s, \eta_{s}^{n}\right)-\int_{s}^{t} e^{-\lambda F\left(u, \eta_{u}^{n}\right)}\left(\frac{\partial}{\partial_{u}}+L_{n}\right) e^{\lambda F\left(u, \eta_{u}^{n}\right)} d u\right\}
$$

is a martingale. We may expand $\mathcal{Z}_{s, t}$ in terms of $\lambda$ as

$$
\begin{aligned}
\mathcal{Z}_{s, t}=\exp \{ & \lambda\left(\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)\right) \\
& -\frac{\lambda^{2}}{2}\left\langle\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)\right\rangle \\
& \left.+\frac{\lambda^{3}}{3!} \int_{s}^{t} \mathcal{R}_{1} d u+\frac{\lambda^{4}}{4!} \int_{s}^{t} \mathcal{R}_{2} d u+\lambda^{5} \int_{s}^{t} \mathcal{R}_{3} d u\right\}
\end{aligned}
$$

where

$$
\begin{aligned}
\mathcal{R}_{1}(u)= & \frac{n^{2}}{2 n^{9 / 2}} \sum_{x \in \mathbb{Z}}\left(b_{x}^{R}(\eta)-b_{x}^{L}(\eta)\right)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right)^{3} \\
& +\frac{a n^{2}}{2 n^{9 / 2+(1 / 2+\gamma)}} \sum_{x \in \mathbb{Z}} b_{x}(\eta)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right)^{3} \\
\mathcal{R}_{2}(u)= & \frac{n^{2}}{2 n^{6}} \sum_{x \in \mathbb{Z}} b_{x}(\eta)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right)^{4} \\
& +\frac{a n^{2}}{2 n^{6+(1 / 2+\gamma)}} \sum_{x \in \mathbb{Z}}\left(b_{x}^{R}(\eta)-b_{x}^{L}(\eta)\right)\left(\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right)^{4} .
\end{aligned}
$$

By the gradient condition and the bound on $b$ in assumption (R), one may compute for $i=1,2$ that

$$
\begin{equation*}
\left\|\mathcal{R}_{i}(u)\right\|_{L^{4}\left(v_{\rho}\right)} \leq \frac{C(a)}{n^{3 / 2}}\|b(\eta)\|_{L^{4}\left(v_{\rho}\right)}\left(\frac{1}{n} \sum_{x}\left|\nabla_{x}^{n} \widetilde{H}_{\gamma, u}\right|^{2+i}\right) \tag{3.4}
\end{equation*}
$$

Since $\mathbb{E}_{\nu_{\rho}}\left[\mathcal{Z}_{s, t}\right]=1$, by expanding in powers of $\lambda$, using Schwarz inequality, the bound on the quadratic variation (3.3), bounds on $\mathcal{R}_{i}$ (3.4) and invariance of $v_{\rho}$, we obtain a bound for the fourth moment of $\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)$ :

$$
\begin{align*}
& \mathbb{E}_{\nu_{\rho}}\left[\left(\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)\right)^{4}\right] \\
& \quad \leq C(a, H)\|b\|_{L}^{4}\left(v_{\rho}\right)\left(|t-s|^{2}+n^{-3 / 2}|t-s|\right) \tag{3.5}
\end{align*}
$$

3.2. Generalized Boltzmann-Gibbs principles. To treat the stochastic differential of $\mathcal{Y}_{t}^{n, \gamma}$, we replace the spatial terms of form $\sum_{x \in \mathbb{Z}} h(x) \tau_{x} f(\eta)$, where $h$ is a function on $\mathbb{Z}$ and $f$ is a local function, in terms of the fluctuation field itself to close the evolution equations. Such replacements fall under the term "BoltzmannGibbs principles" coined by Brox-Rost in [18] which have general validity. For instance, the following result forms the backbone of the argument for Proposition 2.1, when starting from the invariant measure $v_{\rho}$, with respect to the papers cited just before the proposition statement.

## Proposition 3.1.

Let $f$ be a local $L^{2}\left(v_{\rho}\right)$ function. For $t \geq 0$ and $h \in \ell^{2}(\mathbb{Z})$, we have

$$
\lim _{n \rightarrow \infty} \mathbb{E}_{v_{\rho}}\left[\left(\int_{0}^{t} \frac{1}{\sqrt{n}} \sum_{x \in \mathbb{Z}}\left(\tau_{x} f\left(\eta_{s}^{n}\right)-\varphi_{f}(\rho)-\varphi_{f}^{\prime}(\rho)\left(\eta_{s}^{n}(x)-\rho\right)\right) h(x) d s\right)^{2}\right]=0
$$

We now state a main result of this paper which provides a sharper estimate, perhaps of independent interest, when starting from $v_{\rho}$. To simplify expressions,
we will use the notation

$$
\left(\eta_{s}^{n}\right)^{(\ell)}(x):=\frac{1}{2 \ell+1} \sum_{y \in \Lambda_{\ell}} \eta_{s}^{n}(x+y)
$$

THEOREM 3.2 ( $L^{2}$ generalized Boltzmann-Gibbs principle). Let $f$ be a local $L^{5}\left(v_{\rho}\right)$ function supported on sites $\Lambda_{\ell_{0}}$ such that $\varphi_{f}(\rho)=\varphi_{f}^{\prime}(\rho)=0$. There exists a constant $C=C\left(\rho, \ell_{0}\right)$ such that, for $t \geq 0, \ell \geq \ell_{0}^{3}$ and $h \in \ell^{1}(\mathbb{Z}) \cap \ell^{2}(\mathbb{Z})$,

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}} & {\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}}\left(\tau_{x} f\left(\eta_{s}^{n}\right)-\frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\left\{\left(\left(\eta_{s}^{n}\right)^{(\ell)}(x)-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1}\right\}\right) h(x) d s\right)^{2}\right] } \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2}\left(\frac{t \ell}{n}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}} h^{2}(x)\right)+\frac{t^{2} n^{2}}{\ell^{2+\alpha_{0}}}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}|h(x)|\right)^{2}\right)
\end{aligned}
$$

On the other hand, when only $\varphi_{f}(\rho)=0$ is known,

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}} & {\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}}\left(\tau_{x} f\left(\eta_{s}^{n}\right)-\varphi_{f}^{\prime}(\rho)\left\{\left(\eta_{s}^{n}\right)^{(\ell)}(x)-\rho\right\}\right) h(x) d s\right)^{2}\right] } \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2}\left(\frac{t \ell^{2}}{n}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}} h^{2}(x)\right)+\frac{t^{2} n^{2}}{\ell^{1+\alpha_{0}}}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}|h(x)|\right)^{2}\right) .
\end{aligned}
$$

Here, $\alpha_{0}>0$ is the power in assumption (EE).
The proof of Theorem 3.2 is given in Section 4. We note, if the uniform spectral gap holds, $\sup _{k, \xi, n} \ell^{-2} W(k, \ell, \xi, n)<\infty$, then the argument shows one can replace in the right-hand sides above $\|f\|_{L^{5}\left(v_{\rho}\right)}$ with $\|f\|_{L^{3}\left(v_{\rho}\right)}$.
3.3. Tightness. We prove tightness of the fluctuation fields, first starting from the invariant measure $v_{\rho}$, using the $L^{2}$ generalized Boltzmann-Gibbs principle. Then by the relative entropy bound (2.6), we deduce tightness when beginning from initial measures $\left\{\mu^{n}\right\}$.

Proposition 3.3. The sequences $\left\{\mathcal{Y}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{M}_{t}^{n, \gamma}: t \in\right.$ $[0, T]\}_{n \geq 1},\left\{\mathcal{I}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{B}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{K}_{t}^{n, \gamma}: t \in[0, T]\right\}$ and $\left\{\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle: t \in[0, T]\right\}_{n \geq 1}$, when starting from the invariant measure $v_{\rho}$, are tight in the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$.

Proof. By Mitoma's criterion [42], to prove tightness of the sequences with respect to the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$, it is enough to show tightness of $\left\{\mathcal{Y}_{t}^{n, \gamma}(H) ; t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{M}_{t}^{n, \gamma}(H): t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{I}_{t}^{n, \gamma}(H): t \in\right.$ $[0, T]\}_{n \geq 1},\left\{\mathcal{B}_{t}^{n, \gamma}(H): t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{K}_{t}^{n, \gamma}(H): t \in[0, T]\right\}$ and $\left\{\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle: t \in\right.$
$[0, T]\}_{n \geq 1}$, with respect to the uniform topology for all $H \in \mathbb{S}(\mathbb{R})$. Note that all initial values vanish, except $\mathcal{Y}_{0}^{n, \gamma}(H)$.

Tightness of $\mathcal{Y}_{t}^{n, \gamma}(H)$, in view of the decomposition $\mathcal{Y}_{t}^{n, \gamma}(H)=\mathcal{Y}_{0}^{n, \gamma}(H)+$ $\mathcal{I}_{t}^{n, \gamma}(H)+\mathcal{B}_{t}^{n, \gamma}(H)+\mathcal{K}_{t}^{n, \gamma}(H)+\mathcal{M}_{t}^{n, \gamma}(H)$, will follow from tightness of each term. The tightness of $\mathcal{Y}_{0}^{n, \gamma}(H)$, given that we begin under $v_{\rho}$, follows from assumption (IM).

For the martingale term, we use Doob's inequality and stationarity to obtain

$$
\begin{aligned}
& \mathbb{P}_{v_{\rho}}\left(\sup _{\substack{|t-s| \leq \delta \\
0 \leq s, t \leq T}}\left|\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)\right|>\varepsilon\right) \\
& \quad \leq \varepsilon^{-4} \mathbb{E}_{v_{\rho}}\left[\sup _{|t-s| \leq \delta}\left|\mathcal{M}_{t}^{n, \gamma}(H)-\mathcal{M}_{s}^{n, \gamma}(H)\right|^{4}\right] \\
& \quad \leq C \varepsilon^{-4} \delta^{-1} \mathbb{E}_{\nu_{\rho}}\left[\left(\mathcal{M}_{\delta}^{n, \gamma}(H)\right)^{4}\right] .
\end{aligned}
$$

Now, by the fourth moment estimate (3.5), we have

$$
\delta^{-1} \mathbb{E}_{\nu_{\rho}}\left[\left(\mathcal{M}_{\delta}^{n}(H)\right)^{4}\right] \leq C\|b\|_{L^{4}\left(\nu_{\rho}\right)}\left(\delta+n^{-3 / 2}\right)
$$

which vanishes as $n \uparrow \infty$ and then $\delta \downarrow 0$. This is enough to conclude that $\left\{\mathcal{M}_{t}^{n, \gamma}(H): t \in[0, T]\right\}_{n \geq 1}$ is tight in the uniform topology.

We now prove tightness for $\mathcal{B}_{t}^{n, \gamma}(H)$ through the Kolmogorov-Centsov criterion. The argument for $\mathcal{I}_{t}^{n, \gamma}(H)$ is similar. Also, the proofs for $\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle$ and $\mathcal{K}_{t}^{n, \gamma}(H)$, given their forms, are simpler and can be done using invariance of $v_{\rho}$ by squaring all terms. We focus on the case $\gamma=1 / 2$, given that the estimates are analogous and simpler when $1 / 2<\gamma \leq 1$. Let

$$
V_{b}(\eta)=b(\eta)-\varphi_{b}(\rho)-\varphi_{b}^{\prime}(\rho)(\eta(0)-\rho)
$$

By assumption (R1), $V_{b}$ has range $R$. Also, by its form, $\varphi_{V_{b}}(\rho)=\varphi_{V_{b}}^{\prime}(\rho)=0$ and also $\varphi_{V_{b}}^{\prime \prime}(\rho)=\varphi_{b}^{\prime \prime}(\rho)$.

Then

$$
\mathcal{B}_{t}^{n, \gamma}(H)=\frac{a}{2} \int_{0}^{t} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H_{\gamma, s}\right) \tau_{x} V_{b}\left(\eta_{s}\right) d s
$$

By invoking Theorem 3.2 and translation-invariance of $v_{\rho}$ which allows to replace $\nabla_{x}^{n} H_{\gamma, s}$ with $\nabla_{x}^{n} H$ (which does not depend on time $s$ ), for $\ell \geq \ell_{0}^{3}=R^{3}$, with respect to a constant $C=C(a, \rho, R)$, we have

$$
\begin{align*}
\mathbb{E}_{v_{\rho}}[ & \left(\mathcal{B}_{t}^{n, \gamma}(H)\right. \\
& \left.\left.-\frac{a}{4} \int_{0}^{t} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H_{\gamma, s}\right) \varphi_{b}^{\prime \prime}(\rho)\left\{\left(\left(\eta_{s}^{n}\right)^{(\ell)}(x)-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1}\right\} d s\right)^{2}\right]  \tag{3.6}\\
\leq & C\|b\|_{L^{4}\left(v_{\rho}\right)}^{2}\left\{\frac{t \ell}{n}+\frac{t^{2} n^{2}}{\ell^{2+\alpha_{0}}}\right\}\left[\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right)^{2}\right)+\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left|\nabla_{x}^{n} H\right|\right)^{2}\right] .
\end{align*}
$$

On the other hand, given $\sup _{\ell \geq R} E_{\nu_{\rho}}\left[\left(\sqrt{\ell}\left(\eta^{\ell}-\rho\right)\right)^{4}\right]<\infty$ by assumption (IM) and $\left|\varphi_{b}^{\prime \prime}(\rho)\right| \leq C\|b\|_{L^{2}\left(v_{\rho}\right)}$ by assumption (D), and the Schwarz inequality $\left(\sum_{x} h(x) r(x)\right)^{2} \leq\left(\sum_{x}|h(x)|\right) \sum_{x}|h(x)| r^{2}(x)$, we have for $\ell>R^{3}$ that

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}} & {\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H_{\gamma, s}\right) \frac{\varphi_{b}^{\prime \prime}(\rho)}{2}\left\{\left(\left(\eta_{s}^{n}\right)^{(\ell)}(x)-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1}\right\} d s\right)^{2}\right] } \\
& \leq C(\rho)\|b\|_{L^{2}\left(v_{\rho}\right)}^{2} \frac{t^{2} n^{2}}{\ell^{2}}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left|\nabla_{x}^{n} H\right|\right)^{2}
\end{aligned}
$$

Hence, for $\ell>R^{3}$, we have $\mathbb{E}_{\nu_{\rho}}\left[\left(\mathcal{B}_{t}^{n, \gamma}(H)\right)^{2}\right] \leq C(a, \rho, R, H)\|b\|_{L^{4}\left(v_{\rho}\right)}^{2}[t \ell / n+$ $t^{2} n^{2} / \ell^{2}$ ], noting the domination $n^{2} / \ell^{2+\alpha_{0}} \leq n^{2} / \ell^{2}$. Then, if $\ell$ is taken as $\ell=$ $t^{1 / 3} n>R^{3}$, we conclude $\mathbb{E}_{\nu_{\rho}}\left[\left(\mathcal{B}_{t}^{n, \gamma}(H)\right)^{2}\right] \leq C(a, \rho, R, H)\|b\|_{L^{4}\left(v_{\rho}\right)}^{2} t^{4 / 3}$.

However, when $t^{1 / 3} n \leq R^{3}$, we have by the same Schwarz bound that

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}}\left[\left(\mathcal{B}_{t}^{n, \gamma}(H)\right)^{2}\right] & \leq C(\rho, a)\|b\|_{L^{2}\left(v_{\rho}\right)}^{2} t^{2} n^{2}\left(\frac{1}{n} \sum_{x}\left|\nabla_{x}^{n} H\right|\right)^{2} \\
& \leq C(\rho, a, H, R)\|b\|_{L^{2}\left(v_{\rho}\right)}^{2} t^{4 / 3}
\end{aligned}
$$

This shows tightness of $\mathcal{B}_{t}^{n, \gamma}(H)$.
Combining these estimates, we conclude the proof of the proposition.

We now update to when the process begins from the measures $\left\{\mu^{n}\right\}$.

Proposition 3.4. The fluctuation field sequences $\left\{\mathcal{Y}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1}$, $\left\{\mathcal{M}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{I}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{B}_{t}^{n, \gamma}: t \in[0, T]\right\}_{n \geq 1},\left\{\mathcal{K}_{t}^{n, \gamma}: t \in\right.$ $[0, T]\}_{n \geq 1}$ and $\left\{\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle: t \in[0, T]\right\}_{n \geq 1}$ are tight in the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$ when starting from $\left\{\mu^{n}\right\}$ satisfying assumption (BE).

Proof. As before, all initial values vanish except $\mathcal{Y}_{0}^{n, \gamma}$ which, however, is tight by (CLT). Next, by Proposition 3.3, we have $\lim _{\delta \downarrow 0} \lim _{n \uparrow \infty} \mathbb{P}_{\nu_{\rho}}\left(O_{\delta, \varepsilon}^{n}\right)=0$ where

$$
O_{\delta, \varepsilon}^{n}=\left\{\sup _{\substack{|t-s| \leq \delta \\ s, t \in[0, T]}}\left\|X_{t}^{n}-X_{s}^{n}\right\|>\varepsilon\right\},
$$

and $X_{t}^{n}$ may be equal to $\mathcal{Y}_{t}^{n, \gamma}, \mathcal{M}_{t}^{n, \gamma}, \mathcal{I}_{t}^{n, \gamma}, \mathcal{B}_{t}^{n, \gamma}, \mathcal{K}_{t}^{n, \gamma}$ or $\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle$. Then we have by the entropy inequality (2.6) that also $\lim _{\delta \downarrow 0} \lim _{n \uparrow \infty} \mathbb{P}_{\mu^{n}}\left(O_{\delta, \varepsilon}\right)=0$ which allows to conclude.
3.4. Identification of limit points: Proofs of Theorems 2.2 and 2.3. With tightness (Proposition 3.4) in hand, we now identify the limit points of $\left\{\mathcal{Y}_{t}^{n, \gamma}: t \in\right.$ $[0, T]\}_{n \geq 1}$ and its parts in decomposition (3.2). Let $Q^{n}$ be the distribution of

$$
\left(\mathcal{Y}_{t}^{n, \gamma}, \mathcal{M}_{t}^{n, \gamma}, \mathcal{I}_{t}^{n, \gamma}, \mathcal{B}_{t}^{n, \gamma}, \mathcal{K}_{t}^{n, \gamma},\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle: t \in[0, T]\right)
$$

and let $n^{\prime}$ be a subsequence where $Q^{n^{\prime}}$ converges to a limit point $Q$. Let also $\mathcal{Y}_{t}$, $\mathcal{M}_{t}, \mathcal{I}_{t}, \mathcal{B}_{t}, \mathcal{K}_{t}$ and $\mathcal{D}_{t}$ be the respective limits in distribution of the components. Since tightness is shown in the uniform topology on $D\left([0, T], \mathbb{S}^{\prime}(\mathbb{R})\right)$, we have that $\mathcal{Y}_{t}, \mathcal{M}_{t}, \mathcal{I}_{t}, \mathcal{B}_{t}, \mathcal{K}_{t}$ and $\mathcal{D}_{t}$ have a.s. continuous paths.

Let now $G_{\varepsilon}: \mathbb{R} \rightarrow[0, \infty)$ be a smooth compactly supported function for $0<$ $\varepsilon \leq 1$ which approximates $\iota_{\varepsilon}(z)=\varepsilon^{-1} 1_{[-1,1]}\left(z \varepsilon^{-1}\right)$ as in the definition of energy solution before Theorem 2.3. That is, $\left\|G_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2} \leq 2\left\|\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2}=\varepsilon^{-1}$ and $\lim _{\varepsilon \downarrow 0} \varepsilon^{-1 / 2}\left\|G_{\varepsilon}-\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}=0$. Define

$$
\mathcal{A}_{s, t}^{n, \gamma, \varepsilon}(H):=\int_{s}^{t} \frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right)\left[\tau_{x} \mathcal{Y}_{u}^{n, \gamma}\left(G_{\varepsilon}\right)\right]^{2} d u
$$

Since for fixed $0<\varepsilon \leq 1$ the map $\pi . \mapsto \int_{S}^{t} d u \int d x(\nabla H(x))\left\{\pi_{u}\left(\tau_{-x} G_{\varepsilon}\right)\right\}^{2}$ is continuous in the uniform topology on $D\left([0, T] ; \mathbb{S}^{\prime}(\mathbb{R})\right)$, we have subsequentially in distribution that

$$
\lim _{n^{\prime} \uparrow \infty} \mathcal{A}_{s, t}^{n^{\prime}, \gamma, \varepsilon}(H)=\int_{s}^{t} d u \int d x(\nabla H(x))\left\{\mathcal{Y}_{u}\left(\tau_{-x} G_{\varepsilon}\right)\right\}^{2}=: \mathcal{A}_{s, t}^{\varepsilon}(H)
$$

Proposition 3.5. Suppose the initial distribution is the invariant measure $v_{\rho}$ and $t \in[0, T]$.

When $\gamma=1 / 2$, there is a constant $C=C(a, \rho, R)$ such that

$$
\begin{aligned}
& \lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}}\left[\left|\mathcal{B}_{t}^{n, \gamma}(H)-\frac{a \varphi_{b}^{\prime \prime}(\rho)}{4} \mathcal{A}_{0, t}^{n, \gamma, \varepsilon}(H)\right|^{2}\right] \\
& \quad \leq C t\left(\varepsilon+\varepsilon^{-1}\left\|G_{\varepsilon}-\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2}\right)\|b\|_{L^{4}\left(v_{\rho}\right)}^{2}\left[\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}+\|\nabla H\|_{L^{1}(\mathbb{R})}^{2}\right]
\end{aligned}
$$

Then, in $L^{2}\left(\mathbb{P}_{\nu_{\rho}}\right), A_{0, t}^{\varepsilon}(H)$ is a Cauchy $\varepsilon$-sequence. Hence,

$$
\frac{a \varphi_{b}^{\prime \prime}(\rho)}{4} \mathcal{A}_{0, t}(H):=\lim _{\varepsilon \downarrow 0} \frac{a \varphi_{b}^{\prime \prime}(\rho)}{4} \mathcal{A}_{0, t}^{\varepsilon}(H)=\mathcal{B}_{t}(H)
$$

In particular, we conclude $\mathcal{A}_{s, t}(H) \stackrel{d}{=} \mathcal{A}_{0, t-s}(H)$ does not depend on the specific smoothing family $\left\{G_{\varepsilon}\right\}$. Moreover, when $1 / 2<\gamma \leq 1$, we have $\mathcal{B}_{t}(H)=0$.

In addition, when $1 / 2 \leq \gamma \leq 1$,

$$
\lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}}\left[\left|\mathcal{I}_{t}^{n, \gamma}(H)-\frac{\varphi_{c}^{\prime}(\rho)}{2} \int_{0}^{t} \mathcal{Y}_{s}^{n, \gamma}(\Delta H) d s\right|^{2}\right]=0
$$

$$
\begin{aligned}
\lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}}\left[\left|\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle-\frac{\varphi_{b}(\rho)}{2} t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}\right|^{2}\right] & =0 \\
\lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}}\left[\left|\mathcal{K}_{t}^{n, \gamma}(H)\right|^{2}\right] & =0
\end{aligned}
$$

Then, in $L^{2}\left(\mathbb{P}_{v_{\rho}}\right), \mathcal{K}_{t}(H)=0$ and

$$
\mathcal{I}_{t}(H)=\frac{\varphi_{c}^{\prime}(\rho)}{2} \int_{0}^{t} \mathcal{Y}_{s}(\Delta H) d s \quad \text { and } \quad D_{t}(H)=\frac{\varphi_{b}(\rho)}{2} t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}
$$

Proof. We prove the limit display for $\mathcal{B}_{t}(H)$ when $\gamma=1 / 2$ which shows, by a Fatou's lemma that $\mathbb{E}_{\nu_{\rho}}\left[\left|\mathcal{B}_{t}(H)-\left(a \varphi_{b}^{\prime \prime}(\rho) / 4\right) \mathcal{A}_{0, t}^{\varepsilon}(H)\right|^{2}\right] \leq C(a, \rho, R, H) t[\varepsilon+$ $\left.\varepsilon^{-1}\left\|G_{\varepsilon}-\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2}\right]$. Therefore, $\mathcal{A}_{0, t}^{\varepsilon}(H)$, as a sequence in $\varepsilon$, is Cauchy in $L^{2}\left(\mathbb{P}_{\nu_{\rho}}\right)$. The arguments for $\mathcal{I}_{t}(H), \mathcal{D}_{t}(H)$, and $\mathcal{K}_{t}(H)$, noting their forms, are similar; for $\mathcal{D}_{t}(H)$ and $\mathcal{K}_{t}(H)$, one might also use spatial mixing assumed in (IM). To simplify notation, we will call $n=n^{\prime}$.

Note, for $\ell=\varepsilon n$, that

$$
\begin{aligned}
\sum_{x \in \mathbb{Z}} & \left(\nabla_{x}^{n} H_{\gamma, s}\right)\left(\left(\eta_{s}^{n}\right)^{(\ell)}(x)-\rho\right)^{2} \\
& =\sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H_{\gamma, s}\right)\left(\frac{1}{2 n \varepsilon+1} \sum_{|z| \leq n \varepsilon}\left(\eta_{s}^{n}(z+x)-\rho\right)\right)^{2} \\
& =\frac{1+O\left(n^{-1}\right)}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right)\left[\tau_{x} \mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)\right]^{2}
\end{aligned}
$$

Here, the shift by $n^{-1}\left\lfloor a \varphi_{b}^{\prime}(\rho) s n^{2} /\left(2 n^{\gamma}\right)\right\rfloor$ in $\nabla_{x}^{n} H_{\gamma, s}$ [cf. (3.1)] was transferred to $\tau_{x} \mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)$.

Then, with $\ell=\varepsilon n$, by Theorem 3.2, as in the bound (3.6), we have

$$
\begin{aligned}
\lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}} & {\left[\left(\mathcal{B}_{t}^{n, \gamma}(H)-\frac{a \varphi_{b^{n}}^{\prime \prime}(\rho)}{4} \int_{0}^{t} \frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right) \tau_{x} \mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)^{2} d s\right)^{2}\right] } \\
= & \lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}}\left[\left(\mathcal{B}_{t}^{n, \gamma}(H)\right.\right. \\
& \left.\left.\quad-\frac{a \varphi_{b^{n}}^{\prime \prime}(\rho)}{4} \int_{0}^{t} \frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right) \tau_{x}\left\{\mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \varepsilon}\right\} d s\right)^{2}\right] \\
\leq & \lim _{n \uparrow \infty} C(a, \rho, R)\left\|b^{n}\right\|_{L^{4}\left(v_{\rho}\right)}^{2} t\left(\varepsilon+\frac{1}{\varepsilon^{2+\alpha_{0}} n^{\alpha}}\right) \\
& \times\left[\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right)^{2}\right)+\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left|\nabla_{x}^{n} H\right|\right)^{2}\right] .
\end{aligned}
$$

Here, as the sum of $\nabla_{x}^{n} H_{\gamma, s}$ on $x$ vanishes, we introduced the centering constant $(2 \varepsilon)^{-1} \sigma_{\ell}^{2}(\rho)$ in the second line.

Now,

$$
\mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)^{2}-\mathcal{Y}_{x}^{n, \gamma}\left(G_{\varepsilon}\right)^{2}=\left[\mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)-\mathcal{Y}_{s}^{n, \gamma}\left(G_{\varepsilon}\right)\right]\left[\mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)+\mathcal{Y}_{s}^{n, \gamma}\left(G_{\varepsilon}\right)\right]
$$

and by (IM2)

$$
\mathcal{C}_{v_{\rho}}\left(\iota_{\varepsilon}-G_{\varepsilon}, \iota_{\varepsilon}-G_{\varepsilon}\right)^{1 / 2} \mathcal{C}_{v_{\rho}}\left(\iota_{\varepsilon}+G_{\varepsilon}, \iota_{\varepsilon}+G_{\varepsilon}\right)^{1 / 2} \leq C(\rho) \varepsilon^{-1 / 2}\left\|G_{\varepsilon}-\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})} .
$$

Hence, by Schwarz inequality,

$$
\begin{gathered}
\lim _{n \uparrow \infty} \mathbb{E}_{v_{\rho}}\left[\left(\int_{0}^{t} \frac{1}{n} \sum_{x \in \mathbb{Z}}\left(\nabla_{x}^{n} H\right) \tau_{x} \mathcal{Y}_{s}^{n, \gamma}\left(\iota_{\varepsilon}\right)^{2} d s-\mathcal{A}_{0, t}^{n, \gamma, \varepsilon}(H)\right)^{2}\right] \\
\leq C(\rho) \varepsilon^{-1}\left\|G_{\varepsilon}-\iota_{\varepsilon}\right\|_{L^{2}(\mathbb{R})}^{2} t^{2}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}\left|\nabla_{x}^{n} H\right|\right)^{2}
\end{gathered}
$$

Finally, combining these estimates with the inequality $(a+b)^{2} \leq 2 a^{2}+2 b^{2}$, and by assumption (D) that $\lim _{n \uparrow \infty} \varphi_{b^{n}}^{\prime \prime}(\rho)=\varphi_{b}^{\prime \prime}(\rho)$, we complete the proof.

Proposition 3.6. Suppose the initial measures $\left\{\mu^{n}\right\}$ satisfy assumption (BE), and $t \in[0, T]$.

When $\gamma=1 / 2$, we have $\mathcal{A}_{0, t}^{\varepsilon}(H)$ is a Cauchy $\varepsilon$-sequence in probability with respect to a limit measure $Q$, and hence

$$
\frac{a \varphi_{b}^{\prime \prime}(\rho)}{4} \mathcal{A}_{0, t}(H):=\lim _{\varepsilon \downarrow 0} \frac{a \varphi_{b}^{\prime \prime}(\rho)}{4} \mathcal{A}_{0, t}^{\varepsilon}(H)=\mathcal{B}_{t}(H)
$$

On the other hand, when $1 / 2<\gamma \leq 1$, we have $\mathcal{B}_{t}(H) \equiv 0$.
When $1 / 2 \leq \gamma \leq 1$, we have $\mathcal{K}_{t}^{n}(H) \equiv 0$,

$$
\mathcal{I}_{t}(H)=\frac{\varphi_{c}^{\prime}(\rho)}{2} \int_{0}^{t} \mathcal{Y}_{s}(\Delta H) d s \quad \text { and } \quad \mathcal{D}_{t}(H)=\frac{\varphi_{b}(\rho)}{2} t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}
$$

Proof. By assumption (BE), and lower semicontinuity of entropy, the limit measure $Q$ also has bounded entropy with respect to $\mathbb{P}_{\nu_{\rho}}, H\left(Q ; \mathbb{P}_{v_{\rho}}\right)<\infty$. When $\gamma=1 / 2$, by the $L^{2}\left(\mathbb{P}_{v_{\rho}}\right)$ statements in Proposition 3.5 and the entropy inequality (2.6), we have for $\delta>0$ that $\lim _{\varepsilon \downarrow 0} Q\left(\left|\mathcal{B}_{t}(H)-\left(a \varphi_{b}^{\prime \prime}(\rho) / 4\right) \mathcal{A}_{t}^{\varepsilon}(H)\right|>\right.$ $\delta)=0$, and so $\mathcal{A}_{t}^{\varepsilon}(H)$ is Cauchy in probability with respect to $Q$. Therefore, $\lim _{\varepsilon \downarrow 0}\left(a \varphi_{b}^{\prime \prime}(\rho) / 4\right) \mathcal{A}_{t}^{\varepsilon}(H)=\mathcal{B}_{t}(H)$.

The other claims follow similarly.
Proof of Theorems 2.2 and 2.3. Let $H \in \mathbb{S}(\mathbb{R}), t \in[0, T]$, and suppose the initial measures are $\left\{\mu^{n}\right\}$. When $\gamma=1 / 2$, by the decomposition (3.2), Proposition 3.6, and tightness of the constituent processes $\mathcal{Y}_{t}^{n, \gamma}, \mathcal{M}_{t}^{n, \gamma}, \mathcal{I}_{t}^{n, \gamma}, \mathcal{B}_{t}^{n, \gamma}, \mathcal{K}_{t}^{n, \gamma}$ and $\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle$ in the uniform topology, any limit point of

$$
\left(\mathcal{Y}_{t}^{n, \gamma}, \mathcal{M}_{t}^{n, \gamma}, \mathcal{I}_{t}^{n, \gamma}, \mathcal{B}_{t}^{n, \gamma}, \mathcal{K}_{t}^{n, \gamma},\left\langle\mathcal{M}_{t}^{n, \gamma}\right\rangle: t \in[0, T]\right)
$$

satisfies

$$
\mathcal{M}_{t}(H)=\mathcal{Y}_{t}(H)-\mathcal{Y}_{0}(H)-\frac{\varphi_{c}^{\prime}(\rho)}{2} \int_{0}^{t} \mathcal{Y}_{s}(\Delta H) d s-\left(a \varphi_{b}^{\prime \prime}(\rho) / 2\right) \mathcal{A}_{t}(H)
$$

However, when $1 / 2<\gamma \leq 1$,

$$
\begin{equation*}
\mathcal{M}_{t}(H)=\mathcal{Y}_{t}(H)-\mathcal{Y}_{0}(H)-\frac{\varphi_{c}^{\prime}(\rho)}{2} \int_{0}^{t} \mathcal{Y}_{s}(\Delta H) d s \tag{3.7}
\end{equation*}
$$

Also, in both cases, $\mathcal{Y}_{0}(H)=\overline{\mathcal{Y}}_{0}(H)$ by assumption (CLT).
We also claim in both cases that $\mathcal{M}_{t}(H)$ is a continuous martingale with a quadratic variation

$$
\left\langle\mathcal{M}_{t}(H)\right\rangle=\frac{\varphi_{b}(\rho)}{2} t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}
$$

Indeed, by Proposition 3.6, any limit point of the quadratic variation sequence equals $\mathcal{D}_{t}(H)=\left(\varphi_{b}(\rho) / 2\right) t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}$. Next, $\mathcal{M}_{t}(H)$ as the limit of martingales with respect to the uniform topology is a continuous martingale. Also, by the triangle inequality, Doob's inequality and the quadratic variation bound (3.3),

$$
\begin{aligned}
\sup _{n} & \mathbb{E}_{v_{\rho}}\left[\sup _{0 \leq s \leq t}\left|\mathcal{M}_{s}^{n, \gamma}(H)-\mathcal{M}_{s^{-}}^{n, \gamma}(H)\right|\right] \\
& \leq 2 \sup _{n} \mathbb{E}_{v_{\rho}}\left[\sup _{u \in[0, t]}\left|\mathcal{M}_{u}^{n, \gamma}(H)\right|^{2}\right]^{1 / 2} \\
& \leq 2 \sup _{n} \mathbb{E}_{v_{\rho}}\left[\left\langle M_{t}^{n, \gamma}(H)\right\rangle\right]^{1 / 2} \leq C(a, T)\|b\|_{L^{1}\left(v_{\rho}\right)}\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}
\end{aligned}
$$

Then, by Corollary VI.6.30 of [31], $\left(\mathcal{M}_{t}^{n, \gamma}(H),\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle\right)$ converges on a subsequence in distribution to $\left(\mathcal{M}_{t}(H),\left\langle\mathcal{M}_{t}(H)\right\rangle\right)$. Since, also $\left\langle\mathcal{M}_{t}^{n, \gamma}(H)\right\rangle$ converges on a subsequence in distribution to $\mathcal{D}_{t}(H)=\left(\varphi_{b}(\rho) / 2\right) t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}$, we have $\left\langle\mathcal{M}_{t}(H)\right\rangle=\left(\varphi_{b}(\rho) / 2\right) t\|\nabla H\|_{L^{2}(\mathbb{R})}^{2}$.

By Proposition 3.6, when $\gamma=1 / 2, \mathcal{Y}_{t}$ is a "probability energy solution" corresponding to the stochastic Burgers equation (2.11). But, if initially $\mu^{n} \equiv v_{\rho}$, by Proposition 3.5, $\mathcal{Y}_{t}$ is an " $L^{2}$ energy solution." This completes the proof of Theorem 2.3.

However, when $1 / 2<\gamma \leq 1$, by the form of $\mathcal{M}_{t}(H)$ in (3.7), we conclude $\mathcal{Y}_{t}(H)$ solves the Ornstein-Uhlenbeck equation (2.10). By uniqueness, all subsequences converge to the same limit, and we obtain Theorem 2.2.
4. Proof of the generalized Boltzmann-Gibbs principle. We start by recalling the notion of $H_{1, n}$ and $H_{-1, n}$ spaces [35]. For $n \geq n_{0}$, recall $S_{n}=\left(L_{n}+L_{n}^{*}\right) / 2$ [cf. near (2.2)], and define the $H_{1, n}$ seminorm $\|\cdot\|_{1, n}$ on $L^{2}\left(v_{\rho}\right)$ functions by

$$
\|f\|_{1, n}^{2}:=E_{v_{\rho}}\left[f\left(-S_{n}\right) f\right]=n^{2} D_{v_{\rho}}(f)
$$

The Hilbert space $H_{1, n}$ is then the completion of functions with finite $H_{1, n}$ norm modulo norm-zero functions. In particular, local bounded functions are dense in $H_{1, n}$.

Correspondingly, one can define the dual seminorm $\|\cdot\|_{-1, n}$ with respect to the $L^{2}\left(v_{\rho}\right)$ inner-product by

$$
\|f\|_{-1, n}:=\sup \left\{\frac{E_{\mathcal{V}_{\rho}}[f \phi]}{\|\phi\|_{1, n}}: \phi \neq 0 \text { local, bounded }\right\}
$$

and the Hilbert space $H_{-1, n}$ which is the completion over those functions with finite $\|\cdot\|_{-1, n}$ norm modulo norm-zero functions.

We now state a helping lemma for the results in this section. Define the restricted Dirichlet form on local, bounded functions with respect to the grand canonical measure $v_{\rho}$ as

$$
D_{v_{\rho}, \ell}(\phi)=\sum_{x, x+1 \in \Lambda_{\ell}} E_{v_{\rho}}\left[b_{x}^{R, n}(\eta)\left(\nabla_{x, x+1} \phi(\eta)\right)^{2}\right]
$$

Recall the collection $\eta_{r}^{c}:=\left\{\eta(x): x \notin \Lambda_{r}\right\}$.
Proposition 4.1. Let $r: \Omega \rightarrow \mathbb{R}$ be an $L^{4}\left(v_{\rho}\right)$ function and $\ell_{0} \geq 2$. Suppose that $E_{\nu_{\rho}}\left[r \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]=0$ a.s. Then, for local, bounded functions $\phi$, we have

$$
\left|E_{v_{\rho}}[r(\eta) \phi(\eta)]\right| \leq E_{v_{\rho}}\left[W\left(\sum_{x \in \Lambda_{\ell_{0}}} \eta(x), \ell_{0}, \eta_{\ell_{0}}^{c}, n\right)^{2}\right]^{1 / 4}\|r\|_{L^{4}\left(v_{\rho}\right)} D_{v_{\rho}, \ell_{0}}^{1 / 2}(\phi)
$$

Proof. Recall, from Section 2.1, for $k \geq 0, \ell_{0} \geq 2$ and $\xi \in \Omega$, the space

$$
\mathcal{G}_{k, \ell_{0}, \xi}=\left\{\eta: \sum_{x \in \Lambda_{\ell_{0}}} \eta(x)=k, \eta(y)=\xi(y) \text { for } y \notin \Lambda_{\ell_{0}}\right\}
$$

and generator $S_{n, \mathcal{G}}:=S_{n, \mathcal{G}_{k, \ell_{0}, \xi}}$ which governs the evolution of the symmetrized process on $\mathcal{G}_{k, \ell_{0}, \xi}$. Suppose $W\left(k, \ell_{0}, \xi, n\right)<\infty$ so that the measure $v_{k, \ell_{0}, \xi}$ is the unique invariant measure for the process.

Given $E_{v_{\rho}}\left[r \mid \sum_{|x|<\ell_{0}} \eta(x)=k, \eta(y)=\xi(y)\right.$ for $\left.y \notin \Lambda_{\ell_{0}}\right]=E_{v_{k, \ell_{0}, \xi}}[r]=0$, we have $r$ restricted to $\overline{\mathcal{G}}_{k, \ell_{0}, \xi}$ is orthogonal to constant functions and therefore belongs to the range of $-S_{n, \mathcal{G}}$, that is the equation $r=-S_{n, \mathcal{G}} u$ can be solved for some function $u: \mathcal{G}_{k, \ell_{0}, \xi} \rightarrow \mathbb{R}$.

Now, with $k=\sum_{x \in \Lambda_{\ell_{0}}} \eta(x)$ and $\xi=\eta_{\ell_{0}}^{c}, W\left(k, \ell_{0}, \eta_{\ell_{0}}^{c}, n\right)<\infty$ a.s. by assumption (G). Hence,

$$
\begin{aligned}
\left|E_{v_{\rho}}[r \phi]\right| & =\left|E_{v_{\rho}}\left[E_{v_{\rho}}\left[r \phi \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]\right]\right| \\
& =\left|E_{v_{\rho}}\left[E_{v_{\rho}}\left[\left(-S_{n, \mathcal{G}} u\right) \phi \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]\right]\right| \\
& \leq E_{v_{\rho}}\left[E_{v_{\rho}}\left[u\left(-S_{n, \mathcal{G}} u\right) \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]^{1 / 2} E_{v_{\rho}}\left[\phi\left(-S_{n, \mathcal{G}} \phi\right) \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]^{1 / 2}\right]
\end{aligned}
$$

The last line follows as $-S_{n, \mathcal{G}}$ is a nonnegative symmetric operator and, therefore, has a square root.

Further, since $W\left(k, \ell_{0}, \xi, n\right)$ is the reciprocal of the spectral gap for $-S_{n, \mathcal{G}}$, we have

$$
E_{v_{\rho}}\left[r u \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right] \leq W\left(k, \ell_{0}, \eta_{\ell_{0}}^{c}, n\right) E_{v_{\rho}}\left[r^{2} \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right] .
$$

Therefore, we conclude

$$
\left|E_{v_{\rho}}[r \phi]\right| \leq E_{v_{\rho}}\left[W\left(\sum_{x \in \Lambda_{\ell_{0}}} \eta(x), \ell_{0}, \eta_{\ell_{0}}^{c}, n\right) E_{\nu_{\rho}}\left[r^{2} \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]\right]^{1 / 2} D_{\nu_{\rho}, \ell_{0}}^{1 / 2}(\phi)
$$

The desired bound now follows from Schwarz inequality.
The following bound on the variance of additive functionals is the main way we control the fluctuations of several quantities in the sequel. A proof of Proposition 4.2 can be found in Appendix 1.6 of [34].

To simplify notation, for the rest of the section, we will drop the superscript " $n$ " and write $\eta^{n}=\eta$.

Proposition 4.2. Letr $: \Omega \rightarrow \mathbb{R}$ be a mean-zero $L^{2}\left(v_{\rho}\right)$ function, $\varphi_{r}(\rho)=0$. Then

$$
\mathbb{E}_{v_{\rho}}\left[\left(\int_{0}^{t} r\left(\eta_{s}\right) d s\right)^{2}\right] \leq 20 t\|r\|_{-1, n}^{2}
$$

The proof of Theorem 3.2, given at the end of the section, is made through a succession of steps, labeled "one-block," "renormalization step," "two-blocks" and "equivalence of ensembles" estimates.

LEMMA 4.3 (One-block estimate). Let $f: \Omega \rightarrow \mathbb{R}$ be a local $L^{4}\left(v_{\rho}\right)$ function supported on sites in $\Lambda_{\ell_{0}}$ such that $\varphi_{f}(\rho)=0$. Then there exists a constant $C=$ $C(\rho)$ such that for $\ell \geq \ell_{0}, t \geq 0$ and $h \in \ell^{1}(\mathbb{Z}) \cap \ell^{2}(\mathbb{Z})$ :

$$
\begin{aligned}
& \mathbb{E}_{v_{\rho}}\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}} h(x) \tau_{x}\left\{f\left(\eta_{s}\right)-E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(\ell)},\left(\eta_{s}\right)_{\ell}^{c}\right]\right\} d s\right)^{2}\right] \\
& \quad \leq C t \frac{\ell^{3}}{n^{2}}\|f\|_{L^{4}\left(v_{\rho}\right)}^{2} \sum_{x \in \mathbb{Z}} h^{2}(x) .
\end{aligned}
$$

Proof. By Proposition 4.2, we need only to estimate the $H_{-1, n}$ norm of the integrand [which is in $L^{2}\left(v_{\rho}\right)$ since $h \in \ell^{1}(\mathbb{Z})$ ]. Bound the $H_{-1, n}$ norm multiplied by $n$, using Proposition 4.1, as follows:

$$
\begin{align*}
& \sup _{\phi} D_{v_{\rho}}^{-1 / 2}(\phi) E_{v_{\rho}}\left[\sum_{x \in \mathbb{Z}} h(x) \tau_{x}\left\{f-E_{v_{\rho}}\left[f \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]\right\} \phi\right] \\
& \text { 1) }=\sup _{\phi} \sum_{x \in \mathbb{Z}} D_{v_{\rho}}^{-1 / 2}(\phi) E_{v_{\rho}}\left[h(x) \tau_{x}\left(f-E_{v_{\rho}}\left[f \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]\right) \phi\right] \tag{4.1}
\end{align*}
$$

$$
\begin{aligned}
& \leq \sup _{\phi} D_{v_{\rho}}^{-1 / 2}(\phi) \\
& \quad \times \sum_{x \in \mathbb{Z}}|h(x)| E_{v_{\rho}}\left[W\left(\sum_{x \in \Lambda_{\ell}} \eta(x), \ell, \eta_{\ell}^{c}, n\right)^{2}\right]^{1 / 4}\|f\|_{L^{4}\left(v_{\rho}\right)} D_{v_{\rho}, \ell}^{1 / 2}\left(\tau_{-x} \phi\right)
\end{aligned}
$$

Observe now, by translation-invariance of $v_{\rho}$, that

$$
\sum_{x \in \mathbb{Z}} D_{v_{\rho}, \ell}\left(\tau_{-x} \phi\right) \leq(2 \ell+1) D_{v_{\rho}}(\phi)
$$

Then, noting the spectral gap assumption (G), and using the relation $2 a b=$ $\inf _{\kappa>0}\left[a^{2} \kappa+\kappa^{-1} b^{2}\right]$, we bound (4.1) by

$$
\begin{aligned}
& \sup _{\phi} D_{v_{\rho}}^{-1 / 2}(\phi) \inf _{\kappa>0}\left\{\kappa C \ell^{2}\|f\|_{L^{4}\left(v_{\rho}\right)}^{2} \sum_{x \in \mathbb{Z}} h^{2}(x)+\kappa^{-1} C \ell D_{v_{\rho}}(\phi)\right\} \\
& \quad \leq\left(C \ell^{3}\|f\|_{L^{4}\left(v_{\rho}\right)}^{2} \sum_{x \in \mathbb{Z}} h^{2}(x)\right)^{1 / 2}
\end{aligned}
$$

where $C=C(\rho)$ is a constant. This completes the proof.
Now we double the size of the box in the conditional expectation.
Lemma 4.4 (Renormalization step). Let $f: \Omega \rightarrow \mathbb{R}$ be a local $L^{5}\left(v_{\rho}\right)$ function supported on sites in $\Lambda_{\ell_{0}}$ such that $\varphi_{f}(\rho)=\varphi_{f}^{\prime}(\rho)=0$. There exists a constant $C=C\left(\rho, \ell_{0}\right)$ such that for $\ell \geq \ell_{0}, t \geq 0$ and $h \in \ell^{1}(\mathbb{Z}) \cap \ell^{2}(\mathbb{Z})$ :

$$
\begin{aligned}
& \mathbb{E}_{v_{\rho}}\left[\left(\int _ { 0 } ^ { t } \sum _ { x \in \mathbb { Z } } \tau _ { x } \left\{E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(\ell)},\left(\eta_{s}\right)_{\ell}^{c}\right]\right.\right.\right. \\
& \left.\left.\left.\quad-E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(2 \ell)},\left(\eta_{s}\right)_{2 \ell}^{c}\right]\right\} h(x) d s\right)^{2}\right] \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} t \frac{\ell}{n^{2}} \sum_{x \in \mathbb{Z}} h^{2}(x)
\end{aligned}
$$

On the other hand, when only $\varphi_{f}(\rho)=0$ is known,

$$
\begin{aligned}
& \mathbb{E}_{v_{\rho}}\left[\left(\int _ { 0 } ^ { t } \sum _ { x \in \mathbb { Z } } \tau _ { x } \left\{E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(\ell)},\left(\eta_{s}\right)_{\ell}^{c}\right]\right.\right.\right. \\
& \left.\left.\left.\quad-E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(2 \ell)},\left(\eta_{s}\right)_{2 \ell}^{c}\right]\right\} h(x) d s\right)^{2}\right] \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} t \frac{\ell^{2}}{n^{2}} \sum_{x \in \mathbb{Z}} h^{2}(x)
\end{aligned}
$$

Proof. We prove the first statement as the second is similar. Since

$$
E_{\nu_{\rho}}\left[E_{\nu_{\rho}}\left[f(\eta) \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right] \mid \eta^{(2 \ell)}, \eta_{2 \ell}^{c}\right]=E_{\nu_{\rho}}\left[f(\eta) \mid \eta^{(2 \ell)}, \eta_{2 \ell}^{c}\right],
$$

we follow now the same steps as in the proof of Lemma 4.3 to the last line. To finish the proof, we now give an order $O\left(\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} \ell^{-2}\right)$ bound on the variance

$$
\left\|E_{v_{\rho}}\left[f(\eta) \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]-E_{v_{\rho}}\left[f(\eta) \mid \eta^{(2 \ell)}, \eta_{2 \ell}^{c}\right]\right\|_{L^{4}\left(v_{\rho}\right)}^{2}
$$

Adding and subtracting terms, and the inequality $(a+b+c)^{2} \leq 3 a^{2}+3 b^{2}+3 c^{2}$, the variance is bounded by

$$
\begin{aligned}
& \leq 3\left\|E_{v_{\rho}}\left[\left.f(\eta)-\frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\left\{\left(\eta^{(\ell)}-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1}\right\} \right\rvert\, \eta^{(\ell)}, \eta_{\ell}^{c}\right]\right\|_{L^{4}\left(v_{\rho}\right)}^{2} \\
&+3\left\|E_{v_{\rho}}\left[\left.f(\eta)-\frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\left\{\left(\eta^{(2 \ell)}-\rho\right)^{2}-\frac{\sigma_{2 \ell}^{2}(\rho)}{2(2 \ell+1)}\right\} \right\rvert\, \eta^{(2 \ell)}, \eta_{2 \ell}^{c}\right]\right\|_{L^{4}\left(v_{\rho}\right)}^{2} \\
&+3 \| \frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\left\{E_{v_{\rho}}\left[\left.\left(\eta^{(\ell)}-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1} \right\rvert\, \eta^{(\ell)}, \eta_{\ell}^{c}\right]\right. \\
&\left.+E_{v_{\rho}}\left[\left.\left(\eta^{(2 \ell)}-\rho\right)^{2}-\frac{\sigma_{2 \ell}^{2}(\rho)}{2(2 \ell+1)} \right\rvert\, \eta^{(2 \ell)}, \eta_{2 \ell}^{c}\right]\right\} \|_{L^{4}\left(v_{\rho}\right)}^{2}
\end{aligned}
$$

The last term, by the fourth moment bound of $\left(\eta^{(k)}-\rho\right)^{2}$ in (IM2) with $k=\ell$ and $k=2 \ell$ and that $\left|\varphi_{f}^{\prime \prime}(\rho)\right| \leq C(\rho)\|f\|_{L^{2}\left(v_{\rho}\right)}$ in (D), is of order $O\left(\|f\|_{L^{2}\left(\nu_{\rho}\right)}^{2} \ell^{-2}\right)$. But the first two terms are of order $O\left(\|f\|_{L^{5}\left(\nu_{\rho}\right)}^{2} \ell^{-2+\alpha_{0}}\right)$ by applying the equivalence of ensembles assumption (EE).

Lemma 4.5 (Two-blocks estimate). Let $f: \Omega \rightarrow \mathbb{R}$ be a local $L^{5}\left(v_{\rho}\right)$ function supported on sites in $\Lambda_{\ell_{0}}$ such that $\varphi_{f}(\rho)=\varphi_{f}^{\prime}(\rho)=0$. Then, there exists a constant $C=C\left(\rho, \ell_{0}\right)$ such that for $\ell \geq \ell_{0}, t \geq 0$ and $h \in \ell^{1}(\mathbb{Z}) \cap \ell^{2}(\mathbb{Z})$ :

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}} & {\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}} \tau_{x}\left\{E_{v_{\rho}}\left[f(\eta) \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]-E_{v_{\rho}}\left[f(\eta) \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]\right\} h(x) d s\right)^{2}\right] } \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} t \frac{\ell}{n^{2}} \sum_{x \in \mathbb{Z}} h^{2}(x) .
\end{aligned}
$$

On the other hand, when only $\varphi_{f}(\rho)=0$ is known,

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}} & {\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}} \tau_{x}\left\{E_{v_{\rho}}\left[f(\eta) \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]-E_{v_{\rho}}\left[f(\eta) \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right]\right\} h(x) d s\right)^{2}\right] } \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} t \frac{\ell^{2}}{n^{2}} \sum_{x \in \mathbb{Z}} h^{2}(x)
\end{aligned}
$$

Proof. We prove the first display as the second is analogous. Again, we invoke Proposition 4.2 and bound the square of the $H_{-1, n}$ norm. To this end, write $\ell=2^{m+1} \ell_{0}+r$ where $0 \leq r \leq 2^{m+1} \ell_{0}-1$. Then

$$
\begin{aligned}
E_{v_{\rho}} & {\left[f(\eta) \mid \eta^{\left(\ell_{0}\right)}, \eta_{\ell_{0}}^{c}\right]-E_{v_{\rho}}\left[f(\eta) \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right] } \\
= & E_{v_{\rho}}\left[f(\eta) \mid \eta^{\left(2^{m+1} \ell_{0}\right)}, \eta_{2^{m+1} \ell_{0}}^{c}\right]-E_{v_{\rho}}\left[f(\eta) \mid \eta^{(\ell)}, \eta_{\ell}^{c}\right] \\
& \quad+\sum_{i=0}^{m}\left\{E_{v_{\rho}}\left[f(\eta) \mid \eta^{\left(2^{i} \ell_{0}\right)}, \eta_{2^{i} \ell_{0}}^{c}\right]-E_{v_{\rho}}\left[f(\eta) \mid \eta^{\left(2^{i+1} \ell_{0}\right)}, \eta_{2^{i+1} \ell_{0}}^{c}\right]\right\} .
\end{aligned}
$$

Now, by Minkowski's inequality, with respect to the $H_{-1, n}$ norm, over the $m+2$ terms, and Lemma 4.4, we obtain that the left-hand side of the display in the lemma statement is bounded by

$$
\begin{aligned}
& \left\{\left(\frac{C t 2^{m+1} \ell_{0}}{n^{2}}\right)^{1 / 2}+\sum_{i=0}^{m}\left(\frac{C t 2^{i} \ell_{0}}{n^{2}}\right)^{1 / 2}\right\}^{2}\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} \sum_{x \in \mathbb{Z}} h^{2}(x) \\
& \quad \leq \frac{C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} t \ell}{n^{2}} \sum_{x \in \mathbb{Z}} h^{2}(x)
\end{aligned}
$$

to finish the proof.
LEmmA 4.6 (Equivalence of ensembles estimate). Let $f: \Omega \rightarrow \mathbb{R}$ be a local $L^{5}\left(v_{\rho}\right)$ function supported on sites in $\Lambda_{\ell_{0}}$ such that $\varphi_{f}(\rho)=\varphi_{f}^{\prime}(\rho)=0$. Then, there exists a constant $C=C\left(\rho, \ell_{0}\right)$ such that for $\ell \geq \ell_{0}, t \geq 0$ and $h \in \ell^{1}(\mathbb{Z})$ :

$$
\begin{aligned}
& \mathbb{E}_{v_{\rho}}\left[\left(\int _ { 0 } ^ { t } \sum _ { x \in \mathbb { Z } } \tau _ { x } \left\{E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(\ell)},\left(\eta_{s}\right)_{\ell}^{c}\right]\right.\right.\right. \\
& \left.\left.\left.\quad-\frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\left(\left(\eta_{s}^{(\ell)}-\rho\right)^{2}-\frac{\sigma_{\ell}^{2}(\rho)}{2 \ell+1}\right)\right\} h(x) d s\right)^{2}\right] \\
& \quad \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{2} t^{2} \frac{n^{2}}{\ell^{2+\alpha_{0}}}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}|h(x)|\right)^{2} .
\end{aligned}
$$

On the other hand, when only $\varphi_{f}(\rho)=0$ is known,

$$
\begin{aligned}
\mathbb{E}_{v_{\rho}} & {\left[\left(\int_{0}^{t} \sum_{x \in \mathbb{Z}} \tau_{x}\left\{E_{v_{\rho}}\left[f\left(\eta_{s}\right) \mid \eta_{s}^{(\ell)},\left(\eta_{s}\right)_{\ell}^{c}\right]-\varphi_{f}^{\prime}(\rho)\left(\eta_{s}^{(\ell)}-\rho\right)\right\} h(x) d s\right)^{2}\right] } \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)^{2}}^{2} \frac{n^{2}}{\ell^{1+\alpha_{0}}}\left(\frac{1}{n} \sum_{x \in \mathbb{Z}}|h(x)|\right)^{2} .
\end{aligned}
$$

Here, $\alpha_{0}>0$ is the power mentioned in assumption (EE).

Proof. By squaring and using invariance of $v_{\rho}$, the left-hand side of the display is bounded by

$$
2 t^{2} \mathbb{E}_{v_{\rho}}\left[\left(\sum_{x \in \mathbb{Z}}|h(x)||r(x)|\right)^{2}\right]
$$

where $r(x)$ is the $\tau_{x}$-shifted expression in curly braces in the display of Lemma 4.6. Now, by Schwarz inequality,

$$
\left(\sum_{x \in \mathbb{Z}}|h(x)| r(x)\right)^{2} \leq\left(\sum_{x \in \mathbb{Z}}|h(x)|\right) \sum_{x \in \mathbb{Z}}|h(x)| r^{2}(x)
$$

Since $v_{\rho}$ is translation-invariant, the desired bound is now obtained by noting the form of $r(x)$ and the equivalence of ensembles assumption (EE).

Proof of Theorem 3.2. By combining Lemma 4.3 with $\ell=\ell_{0}$, and Lemmas 4.5 and 4.6, we straightforwardly obtain the result.
5. Equivalence of ensembles. We prove, as a consequence of Proposition 5.1, that condition (EE) holds for a large class of systems with product invariant measures. In this case, $v_{k, \ell, \xi}$ does not depend on $\xi$, which simplifies the conditional expectation in the statement of (EE).

Next, we show in Proposition 5.2 that (EE) also holds for the Markov chain measure $\nu_{1 / 2}$ defined in Section 2.5. Some parts of the proofs of these statements are similar to those in [56].

Define $\Lambda_{m}^{+}=\{x: 1 \leq x \leq m\}$.
PROPOSITION 5.1. Let $v_{\rho}$ be a product measure on $\Omega$ such that (IM) holds, and $0<v_{\rho}(\eta(0)=j)<1$ for $j=0$, Let also $f$ be a local $L^{5}\left(v_{\rho}\right)$ function, supported on sites $\Lambda_{\ell_{0}}^{+}$, such that $\varphi_{f}(\rho)=\varphi_{f}^{\prime}(\rho)=0$. Then there exists a constant $C=C\left(\rho, \ell_{0}\right)$, such that for $n \geq \ell_{0}$ we have

$$
\left\|E_{v_{\rho}}[f(\eta) \mid y]-\left\{y^{2}-\frac{\sigma^{2}(\rho)}{n}\right\} \frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\right\|_{L^{4}\left(v_{\rho}\right)} \leq \frac{C\|f\|_{L^{5}\left(v_{\rho}\right)}}{n^{3 / 2}}
$$

On the other hand, when only $\varphi_{f}(\rho)=0$ is known,

$$
\left\|E_{v_{\rho}}[f(\eta) \mid y]-y \varphi_{f}^{\prime}(\rho)\right\|_{L^{4}\left(v_{\rho}\right)} \leq \frac{C\|f\|_{L^{5}\left(v_{\rho}\right)}}{n}
$$

Here, $y:=\frac{1}{n} \sum_{x \in \Lambda_{n}^{+}} \eta(x)-\rho$.
Proof. We prove the first display as the second statement, following the same scheme, has a simpler argument. At the expense of the constant, we need only to consider all large $n>\ell_{0}$. To simplify notation, we will call $\ell=\ell_{0}$. The proof follows in several steps.

Step 1. Recall the tilted measures $\left\{v_{z}: \rho_{*}<z<\rho^{*}\right\}$ given after assumption (D1) which are well defined as $\nu_{\rho}$ is a product measure. Let $\sigma^{2}(z)=E_{\nu_{z}}\left[(\eta(0)-z)^{2}\right]$. Note also the canonical expectation $E_{\nu_{z}}[f \mid y]$ does not depend on $z$, and that we are free to choose it as desired.

Develop

$$
\begin{aligned}
E_{v_{\rho}}[f(\eta) \mid y] & =E_{v_{y+\rho}}\left[f(\eta) \left\lvert\, \frac{1}{n} \sum_{x \in \Lambda_{n}^{+}} \eta(x)-\rho=y\right.\right] \\
& =\frac{E_{v_{y+\rho}}\left[f(\eta) 1\left((1 / n) \sum_{x \in \Lambda_{n}^{+}} \eta(x)-\rho=y\right)\right]}{v_{y+\rho}\left((1 / n) \sum_{x \in \Lambda_{n}^{+}} \eta(x)-\rho=y\right)} .
\end{aligned}
$$

Define $\theta_{m}(z)=\sqrt{m} \nu_{y+\rho}\left(\sum_{x \in \Lambda_{m}^{+}} \eta(x)-\rho-y=z\right)$, and write the last expression as

$$
E_{v_{y+\rho}}\left[f(\eta) \frac{\sqrt{n} \theta_{n-\ell}\left(-\sum_{x \in \Lambda_{\ell}^{+}}(\eta(x)-y-\rho)\right)}{\sqrt{n-\ell} \theta_{n}(0)}\right]
$$

The goal will be now to expand $\theta_{n-\ell}(z)$ to recover the main terms approximating $E_{\nu_{\rho}}[f \mid y]$ when $|y|$ is small. We will treat the case when $|y|$ is bounded away from 0 afterward.

Step 2. To expand $\theta_{m}(z)$, let $\psi_{y}(t)=E_{v_{y+\rho}}\left[e^{i t(\eta(x)-\rho-y)}\right]$ be the characteristic function. Then one can write

$$
\begin{aligned}
\theta_{m}(x) & =\frac{\sqrt{m}}{2 \pi} \int_{-\pi}^{\pi} e^{-i t x} \psi_{y}^{m}(t) d t \\
& =\frac{1}{2 \pi} \int_{-\pi \sqrt{m}}^{\pi \sqrt{m}} e^{-i t x / \sqrt{m}} \psi_{y}^{m}(t / \sqrt{m}) d t
\end{aligned}
$$

By Taylor expansion,

$$
\begin{align*}
2 \pi \theta_{m}(x)= & \int_{-\pi \sqrt{m}}^{\pi \sqrt{m}} \psi_{y}^{m}(t / \sqrt{m}) d t-\int_{-\pi \sqrt{m}}^{\pi \sqrt{m}} \frac{i x t}{\sqrt{m}} \psi_{y}^{m}(t / \sqrt{m}) d t \\
& -\frac{1}{2} \int_{-\pi \sqrt{m}}^{\pi \sqrt{m}} \frac{x^{2} t^{2}}{m} \psi_{y}^{m}(t / \sqrt{m}) d t  \tag{5.1}\\
& +O\left(\frac{|x|^{3}}{m^{3 / 2}}\right) \int_{-\pi \sqrt{m}}^{\pi \sqrt{m}}|t|^{3}\left|\psi_{y}^{m}(t / \sqrt{m})\right| d t
\end{align*}
$$

Step 3. Let $\delta>0$ be such that $(\rho-\delta, \rho+\delta) \subset\left(\rho_{*}, \rho^{*}\right)$ and sufficiently small in the following estimates. Let also $0<\varepsilon \leq \pi$.

First, $\sup _{|y| \leq \delta, \varepsilon \leq|t| \leq \pi}\left|\psi_{y}^{m}(t)\right|<C_{0}^{m}$ where $C_{0}<1$ : Write

$$
\begin{aligned}
\left|\psi_{y}(t)\right| & \leq\left|v_{y+\rho}(\eta(0)=0)+e^{i t} v_{y+\rho}(\eta(0)=1)\right|+\sum_{k \geq 2} v_{y+\rho}(\eta(0)=k) \\
& \leq\left(A^{2}-2 v_{y+\rho}(\eta(0)=0) v_{y+\rho}(\eta(0)=1)[1-\cos (t)]\right)^{1 / 2}+1-A
\end{aligned}
$$

where $A=v_{y+\rho}(\eta(0)=0)+v_{y+\rho}(\eta(0)=1)$. By the proposition assumptions and continuity of $v_{y+\rho}(\eta(0)=k)$ in $y, 0<v_{y+\rho}(\eta(0)=j)<1$ for $j=0$, 1 uniformly for $|y| \leq \delta$. Hence, uniformly over $\varepsilon \leq|t| \leq \pi,|y| \leq \delta$, the right-hand side of the display above is strictly bounded by a constant $C_{0}<1$.

Second, for $0 \leq|t / \sqrt{m}|<\varepsilon$ and $|y| \leq \delta$,

$$
\psi_{y}^{m}(t / \sqrt{m})=\left[1-\left(t^{2} \sigma^{2}(y+\rho) /(2 m)\right)+O\left(C(\delta)|t|^{3} m^{-3 / 2}\right)\right]^{m}
$$

so that $\left|\psi_{y}^{m}(t / \sqrt{m})\right| \leq e^{-C_{1}(y, \varepsilon) t^{2}}$. Here, by continuity in $y$ and $\sigma^{2}(\rho)>0$, when $\varepsilon$ is small, $\inf _{|y| \leq \delta} C_{1}(y, \varepsilon)>0$. Similarly, we note $\sup _{|y| \leq \delta} \sigma^{2}(y+\rho)<\infty$, and $\inf _{|y| \leq \delta} \sigma^{2}(y+\rho)>0$.

Last, by the classical local limit theorem, $\lim _{m \uparrow \infty} \theta_{m}(0)=\left(2 \pi \sigma^{2}(y+\rho)\right)^{-1 / 2}$.
Step 4. We now observe, for $|y| \leq \delta$ and $m \geq 1$, as a consequence of the estimates in step 3, the integral in the last term in (5.1) is uniformly bounded: Split the integral over the ranges $|t / \sqrt{m}|<\varepsilon$ and $|t / \sqrt{m}| \geq \varepsilon$ and bound each part separately.

Also, similarly, we split the second integral in (5.1), when $|y| \leq \delta$, over ranges $|t / \sqrt{m}| \geq \varepsilon$ and $|t / \sqrt{m}|<\varepsilon$. On the first range, the restricted integral exponentially decays, and on the range $|t / \sqrt{m}|<\varepsilon$, the restricted integral is almost the integral of an odd function since here

$$
\psi_{y}^{m}(t / \sqrt{m})=\left(1-\frac{t^{2} \sigma^{2}(y+\rho)}{2 m}\right)^{m}\left[1+O\left(C(\delta)|t|^{3} m^{-1 / 2}\right)\right]
$$

Therefore, we conclude that the second integral in (5.1) is of order $O\left(m^{-1 / 2}\right)$.
Step 5. Then, for $|y| \leq \delta$, we have

$$
\begin{aligned}
E_{v_{\rho}}[f(\eta) \mid y]= & \kappa_{0} E_{v_{y+\rho}}[f(\eta)]+\frac{\kappa_{1}}{\sqrt{n-\ell}} E_{v_{y+\rho}}\left[f(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-\rho-y\right)\right] \\
& +\frac{\kappa_{2}}{n-\ell} E_{v_{y+\rho}}\left[f(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-\rho-y\right)^{2}\right]+\varepsilon_{f}(n)
\end{aligned}
$$

where $\left|\varepsilon_{f}(n)\right| \leq C(\rho, \ell, \delta)\|f\|_{L^{2}\left(v_{\rho}\right)} n^{-3 / 2}$ and $\kappa_{i}=\kappa_{i}(n)$ for $i=0,1,2$ are explicit expressions. Indeed, one observes

$$
\begin{aligned}
\kappa_{0}(n) & =\frac{\sqrt{n}}{\sqrt{n-\ell}} \frac{\theta_{n-\ell}(0)}{\theta_{n}(0)}=1+O\left(n^{-1 / 2}\right) \\
\kappa_{1}(n) & =\frac{\sqrt{n}}{\theta_{n}(0) \sqrt{n-\ell}} \frac{1}{2 \pi} \int_{-\pi \sqrt{n-\ell}}^{\pi \sqrt{n-\ell}} \text { it } \psi_{y}^{n-\ell}\left(\frac{t}{\sqrt{n-\ell}}\right) d t=O\left(n^{-1 / 2}\right) \\
\kappa_{2}(n) & =\frac{-\sqrt{n}}{2 \theta_{n}(0) \sqrt{n-\ell}} \frac{1}{2 \pi} \int_{-\pi \sqrt{n-\ell}}^{\pi \sqrt{n-\ell}} t^{2} \psi_{y}^{n-\ell}\left(\frac{t}{\sqrt{n-\ell}}\right) d t \\
& =\frac{-1}{2 \sigma^{2}(y+\rho)}+O\left(n^{-1 / 2}\right)
\end{aligned}
$$

Step 6. We now develop expansions of $E_{v_{y+\rho}}[h]$ for a local $L^{2}\left(v_{\rho}\right)$ function $h$ supported on coordinates in $\Lambda_{\ell}^{+}$. The "tilting" given in the Introduction, (2.3) reduces to

$$
E_{v_{y+\rho}}[h]=E_{v_{\rho}}\left[h(\eta) \frac{e^{\lambda(y+\rho) \sum_{x \in \Lambda_{\ell}^{+}}(\eta(x)-\rho)}}{M^{\ell}(\lambda(y+\rho))}\right]
$$

where $\lambda(y+\rho)$ is the "tilt" chosen to change the density to $y+\rho$ and $M(\lambda)=$ $E_{v_{\rho}}\left[e^{\lambda(\eta(x)-\rho)}\right]$. Note that $z-\rho=M^{\prime}(\lambda(z)) / M(\lambda(z))$ and

$$
\lambda^{\prime}(z)=\left[\frac{M^{\prime \prime}(\lambda(z))}{M(\lambda(z))}-\left(\frac{M^{\prime}(\lambda(z))}{M(\lambda(z))}\right)^{2}\right]^{-1}=\frac{1}{\sigma^{2}(z)}
$$

Consider the first and second derivatives of $E_{v_{y+\rho}}[h]$ given exactly in (2.4) as $v_{y+\rho}$ is a product measure. The third derivative takes the form

$$
\begin{aligned}
\frac{d^{3}}{d y^{3}} E_{v_{y+\rho}}[h(\eta)]= & \lambda^{\prime \prime \prime}(y+\rho) E_{v_{y+\rho}}\left[\bar{h}(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-y-\rho\right)\right] \\
& +3 \lambda^{\prime}(y+\rho) \lambda^{\prime \prime}(y+\rho) E_{v_{y+\rho}}\left[\bar{h}(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-y-\rho\right)^{2}\right] \\
+ & \left(\lambda^{\prime}(y+\rho)\right)^{3} E_{v_{y+\rho}}\left[\bar{h}(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-y-\rho\right)^{3}\right] \\
& -3\left(\lambda^{\prime}(y+\rho)\right)^{3} E_{v_{y+\rho}}\left[\bar{h}(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-y-\rho\right)\right] \\
& \quad \times E_{v_{y+\rho}}\left[\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-y-\rho\right)^{2}\right]
\end{aligned}
$$

where $\bar{h}(\eta)=h(\eta)-E_{v_{y+\rho}}[h]$.
Then, for $|y| \leq \delta$, when $\varphi_{h}(\rho)=\varphi_{h}^{\prime}(\rho)=0$, we may expand around $y=0$ :

$$
E_{v_{y+\rho}}[h(\eta)]=\left(\lambda^{\prime}(\rho)\right)^{2} \frac{y^{2}}{2} E_{v_{\rho}}\left[h(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-\rho\right)^{2}\right]+|y|^{3} r(\rho, \delta, h)
$$

When only $\varphi_{h}(\rho)=0$ is known,

$$
E_{\nu_{y+\rho}}[h(\eta)]=\lambda^{\prime}(\rho) y E_{\nu_{\rho}}\left[h(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}} \eta(x)-\rho\right)\right]+|y|^{2} r(\rho, \delta, h)
$$

When possibly $\varphi_{h}(\rho) \neq 0$,

$$
E_{v_{y+\rho}}[h(\eta)]=E_{v_{\rho}}[h(\eta)]+|y| r(\rho, \delta, h)
$$

Here, as the first and second derivatives in (2.4) and the third derivative above are bounded for $|y| \leq \delta$, we may conclude that the remainders $|r(\rho, \delta, h)| \leq$ $C(\rho, \delta)\|h\|_{L^{2}\left(v_{\rho}\right)}$.

We now relate the terms $E_{v_{\rho}}\left[h(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)^{k}\right]$ to derivatives $\varphi_{h}^{(k)}(\rho)$ : From (2.4), for $k=1,2$, when $\varphi_{h}^{(k-1)}(\rho)=\varphi_{h}(\rho)=0$, we have

$$
\begin{equation*}
\varphi_{h}^{(k)}(\rho)=\left(\lambda^{\prime}(\rho)\right)^{k} E_{v_{\rho}}\left[h(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)^{k}\right] \tag{5.2}
\end{equation*}
$$

Step 7. Consider the expansion of $E_{v_{\rho}}[f \mid y]$ in step 5 when $|y| \leq \delta$. With $h$ equal to variously $f, f(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)$, and $f(\eta)\left(\sum_{x \in \Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)^{2}$, we may write

$$
\begin{aligned}
E_{v_{\rho}}[f \mid y]= & \frac{\kappa_{0}}{2}\left(\lambda^{\prime}(\rho)\right)^{2} y^{2} E_{v_{\rho}}\left[f(\eta)\left(\sum_{\Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)^{2}\right]+\kappa_{0}|y|^{3} r(f) \\
& +\frac{\kappa_{1} \lambda^{\prime}(\rho) y}{\sqrt{n-\ell}} E_{v_{\rho}}\left[f(\eta)\left(\sum_{\Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)^{2}\right]+\frac{\kappa_{1}}{\sqrt{n-\ell}}|y|^{2} r(f) \\
& +\frac{\kappa_{2}}{n-\ell} E_{v_{\rho}}\left[f(\eta)\left(\sum_{\Lambda_{\ell}^{+}}(\eta(x)-\rho)\right)^{2}\right]+\frac{\kappa_{2}}{n-\ell}|y| r(f)+\varepsilon_{f}(n)
\end{aligned}
$$

where $|r(f)| \leq C(\rho, \ell, \delta)\|f\|_{L^{2}\left(v_{\rho}\right)}^{2}$.
Hence, noting the assumptions on $\varphi_{f}(\rho),(5.2)$, and $E_{v_{\rho}}\left[y^{2 p}\right]=O\left(n^{-p}\right)$ so that each $y$ factor is $O\left(n^{-1 / 2}\right)$, we can group the dominant terms so that

$$
\begin{aligned}
E_{v_{\rho}} & {[1(|y| \leq \delta)} \\
& \left.\times\left(E_{v_{\rho}}[f(\eta) \mid y]-\left\{\frac{\kappa_{0} y^{2}}{2}+\frac{1}{\lambda^{\prime}(\rho)} \frac{\kappa_{1} y}{\sqrt{n}}+\frac{1}{\left(\lambda^{\prime}(\rho)\right)^{2}} \frac{\kappa_{2}}{n}\right\} \varphi_{f}^{\prime \prime}(\rho)\right)^{4}\right] \\
& \leq C(\rho, \delta)\|f\|_{L^{2}\left(v_{\rho}\right)}^{4} n^{-6}
\end{aligned}
$$

Noting $\kappa_{0}(n)=1+O\left(n^{-1 / 2}\right), \kappa_{1}(n)=O\left(n^{-1 / 2}\right)$, formula $\lambda^{\prime}(\rho)=\sigma^{-2}(\rho)$ in step $6,\left|\varphi^{\prime \prime}(\rho)\right| \leq C\|f\|_{L^{2}\left(\nu_{\rho}\right)}$ and, by Taylor expansion of $\sigma^{2}(y+\rho)$ around $y=0$, $\kappa_{2}(n)=-2^{-1} \sigma^{-2}(\rho)+O\left(n^{-1 / 2}\right)$, we have further

$$
E_{v_{\rho}}\left[1(|y| \leq \delta)\left(E_{v_{\rho}}[f(\eta) \mid y]-\left\{y^{2}-\frac{\sigma^{2}(\rho)}{n}\right\} \frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\right)^{4}\right] \leq C\|f\|_{L^{2}\left(v_{\rho}\right)}^{4} n^{-6}
$$

Step 8. On the other hand, by say large deviations estimates, we bound

$$
\begin{aligned}
E_{v_{\rho}} & {\left[1(|y|>\delta)\left(E_{v_{\rho}}[f(\eta) \mid y]-\left\{y^{2}-\frac{\sigma^{2}(\rho)}{n}\right\} \frac{\varphi_{f}^{\prime \prime}(\rho)}{2}\right)^{4}\right] } \\
& \leq C\|f\|_{L^{5}\left(v_{\rho}\right)}^{4} O\left(n^{-6}\right)
\end{aligned}
$$

to complete the proof.

We now prove the equivalence ensembles estimate (EE) with respect to a Markovian measure. Recall the Gibbs measures $v_{1 / 2}$ and $v_{z}=v_{1 / 2}^{\lambda(z)}$, and transition matrix $P$ defined in Section 2.5. To see how the next proposition can be used to satisfy assumption (EE), we note (1) the estimate is uniform in the "outside variables" $\eta_{\ell}^{c}$, and (2) since the transition matrix $P$ is positive, the $L^{\infty}$ norm of any local function supported on sites $\Lambda_{\ell_{0}}$ can be bounded $\|f\|_{L^{\infty}} \leq$ $C\left(\ell_{0}, \beta\right)\|f\|_{L^{p}\left(\nu_{1 / 2}\right)}$ for $p>0$. Recall also the definitions of $\varphi_{f}(\rho)$ and its derivatives in (2.4).

Proposition 5.2. Let $f$ be a local function, supported on sites indexed by $\Lambda_{\ell_{0}}$, such that $\varphi_{f}(1 / 2)=\varphi_{f}^{\prime}(1 / 2)=0$. Then, for each $0<\varepsilon<1$, there is a constant $C=C\left(\ell_{0}, \varepsilon\right)$ such that for $a, b \in\{0,1\}$ and $n \geq \ell_{0}$,

$$
\begin{aligned}
& \left\|E_{\nu_{1 / 2}}[f \mid y, \eta(-n-1)=a, \eta(n+1)=b]-\frac{\varphi_{f}^{\prime \prime}(1 / 2)}{2}\left[y^{2}-\frac{\sigma_{n}^{2}(1 / 2)}{2 n+1}\right]\right\|_{L^{4}\left(\nu_{1 / 2}\right)} \\
& \quad \leq \frac{C\|f\|_{L^{\infty}}}{n^{3 / 2-\varepsilon}}
\end{aligned}
$$

On the other hand, when only $\varphi_{f}(1 / 2)=0$ is known,

$$
\left\|E_{\nu_{1 / 2}}[f \mid y, \eta(-n-1)=a, \eta(n+1)=b]-y \varphi_{f}^{\prime}\left(\frac{1}{2}\right)\right\|_{L^{4}\left(\nu_{1 / 2}\right)} \leq \frac{C\|f\|_{L^{\infty}}}{n^{1-\varepsilon}} .
$$

Here, $y=(2 n+1)^{-1} \sum_{x \in \Lambda_{n}}\left(\eta(x)-\frac{1}{2}\right)$.
Proof. The argument has the same structure as for Proposition 5.1. We will concentrate on the first display for all large $n$; the second statement has a similar argument. Since $\nu_{1 / 2}$ corresponds to an ergodic finite-state Markov chain with uniform invariant measure, it is exponentially mixing and allows standard approximations, which are used in many steps.

Step 1 . Let $\chi>0$ be small and $n^{\prime}=n-n^{\chi}$. Write

$$
\begin{aligned}
& E_{v_{1 / 2}}[f \mid y, \eta(-n-1)=a, \eta(n+1)=b] \\
& \quad=E_{v_{y+1 / 2}}[f(\eta) \mid y, \eta(-n-1)=a, \eta(n+1)=b]
\end{aligned}
$$

$$
\begin{array}{r}
=E_{v_{y+1 / 2}}\left[\left.f(\eta) \frac{\sqrt{2 n+1} \theta_{n, y, a, b}^{\chi}\left(-\sum_{x \in \Lambda_{n} x}(\eta(x)-y-1 / 2)\right)}{\sqrt{2 n^{\prime}} \theta_{n, y, a, b}(0)} \right\rvert\,\right. \\
\eta(-n-1)=a, \eta(n+1)=b]
\end{array}
$$

where

$$
\begin{array}{r}
\theta_{n, y, a, b}^{\chi}(z)=\sqrt{2 n^{\prime}} v_{y+1 / 2}\left(\sum_{n x<|x| \leq n} \eta(x)-y-1 / 2=z \mid\right. \\
\eta\left(n^{\chi}\right), \eta\left(-n^{\chi}\right), \eta(-n-1)=a \\
\eta(n+1)=b), \\
\theta_{n, y, a, b}(z)=\sqrt{2 n+1} v_{y+1 / 2}\left(\sum_{x \in \Lambda_{n}} \eta(x)-y-1 / 2=z\right. \\
\eta(-n-1)=a, \eta(n+1)=b) .
\end{array}
$$

Step 2. Let the characteristic function $\psi_{n, y, \chi, a, b}(t)$ for $|t| \leq \pi$ be defined by

$$
E_{v_{y+1 / 2}}\left[e^{i t \sum_{n} \chi<|x| \leq n}(\eta(x)-y-1 / 2) \mid \eta\left(n^{\chi}\right), \eta\left(-n^{\chi}\right), \eta(-n-1)=a, \eta(n+1)=b\right] .
$$

Let $\delta>0$ be such that $(\rho-\delta, \rho+\delta) \subset(0,1)$ and sufficiently small in the following estimates. Suppose $|y| \leq \delta$. Let also $r>0$ be a small number. We now state a few relations, and then argue them. First, for $|t|<r$, we claim

$$
\begin{equation*}
\psi_{n, y, \chi, a, b}\left(\frac{t}{\sqrt{2 n^{\prime}}}\right)=\left(1-\frac{t^{2} \sigma^{2}(y+1 / 2)}{2\left(2 n^{\prime}\right)}\right)^{2 n^{\prime}}\left[1+O\left(n^{\prime-1 / 2}\right)\right] \tag{5.3}
\end{equation*}
$$

where $\sigma^{2}(z)=\lim _{n \uparrow \infty}(2 n+1)^{-1} E_{v_{z}}\left[\left(\sum_{x \in \Lambda_{n}} \eta(x)-z\right)^{2}\right]$ is the limiting variance of the additive functional $n^{-1 / 2} \sum_{x=1}^{n} \eta(x)-z$ with respect to measure $\nu_{z}$ (cf. formula in Section 2.5).

Therefore, for $|t|<r$ and $C=C(\delta, r)>0$, we have $\left|\psi_{n, y, \chi, a, b}\left(t / \sqrt{2 n^{\prime}}\right)\right|<$ $\exp \left\{-C t^{2}\right\}$.

Next, we claim, for $r \leq|t| \leq \pi$ that

$$
\begin{equation*}
\left|\psi_{n, y, \chi, a, b}(t)\right|<A^{2 n^{\prime}} \tag{5.4}
\end{equation*}
$$

where $A=A(\delta, r)<1$.
We also state a case of a local central limit theorem for ergodic Markov chains [36],

$$
\lim _{n \uparrow \infty} \theta_{n, y, a, b}(0)=\frac{1}{\sqrt{2 \pi \sigma^{2}(y+1 / 2)}}
$$

We now give an argument for the above claims which may be skipped on first reading. Recall $u_{1}$ and $v_{1}$ near (2.15) with $\lambda=\lambda(y+1 / 2)$, and consider the transfer matrix:

$$
\widetilde{P}(s)=\left[\begin{array}{cc}
\left(1-u_{1}\right) e^{s(-1 / 2-y)} & u_{1} e^{s(1 / 2-y)} \\
v_{1} e^{s(-1 / 2-y)} & \left(1-v_{1}\right) e^{s(1 / 2-y)}
\end{array}\right]
$$

By the Markov property, one writes

$$
\begin{equation*}
\psi_{n, y, \chi, a, b}(t)=\frac{\widetilde{P}(i t)^{n^{\prime}}\left(\eta\left(-n^{\chi}\right), a\right) \widetilde{P}(i t)^{n^{\prime}}\left(\eta\left(n^{\chi}\right), b\right)}{\widetilde{P}(0)^{n^{\prime}}\left(\eta\left(-n^{\chi}\right), a\right) \widetilde{P}(0)^{n^{\prime}}\left(\eta\left(n^{\chi}\right), b\right)} \tag{5.5}
\end{equation*}
$$

We may diagonalize $\widetilde{P}(i t / \sqrt{m})^{m}=Q(i t / \sqrt{m}) D^{m}(i t / \sqrt{m}) Q^{-1}(i t / \sqrt{m})$, for large $m$, where $D(t)$ is a diagonal matrix with eigenvalues $w_{1}(t)$ and $w_{2}(t)$ and $Q(t)$ is the matrix of the corresponding eigenvectors. Of course, when $t=0$, $1=w_{1}(0)>w_{2}(0)=1-u_{1}-v_{1}$ with corresponding eigenvectors $\langle 1,1\rangle$ and $\left\langle u_{1},-v_{1}\right\rangle$. For large $m, w_{1}(i t / \sqrt{m})$ is the eigenvalue with maximum absolute value and is expressed as

$$
\begin{aligned}
& \frac{e^{-i t y / \sqrt{m}}}{2}\left[\left(1-u_{1}\right) e^{-i t /(2 \sqrt{m})}+\left(1-v_{1}\right) e^{i t /(2 \sqrt{m})}\right. \\
& \left.\quad+\left\{\left(\left(1-u_{1}\right) e^{-i t /(2 \sqrt{m})}-\left(1-v_{1}\right) e^{i t /(2 \sqrt{m})}\right)^{2}+4 u_{1} v_{1}\right\}^{1 / 2}\right]
\end{aligned}
$$

It is not difficult to check that

$$
\begin{aligned}
w_{1}^{\prime}(0) & =-i t y / \sqrt{m}+i t /(2 \sqrt{m})\left[\left(u_{1}-v_{1}\right) /\left(u_{1}+v_{1}\right)\right] \\
& =(i t / \sqrt{m}) E_{\pi_{1}(y+1 / 2)}[\eta(0)-1 / 2-y]=0
\end{aligned}
$$

where $\pi_{1}(y+1 / 2)=\left(u_{1}+v_{1}\right)^{-1}\left\langle v_{1}, u_{1}\right\rangle$ is the marginal of $v_{y+1 / 2}$ (cf. Section 2.5). One now expands, as all quantities are smooth, $w_{1}(i t / m)=1-$ $\left(t^{2} / 2 m\right) w_{1}^{\prime \prime}(0)+O\left(m^{-3 / 2}\right)$ where the error is uniform for $|y| \leq \delta$ and $|t| \leq \pi$. Similarly, $Q(i t / \sqrt{m})=Q(0)+O\left(m^{-1 / 2}\right)$. One can identify $w_{1}^{\prime \prime}(0)$ as the variance $\sigma^{2}(y+1 / 2)$ since we know

$$
\begin{aligned}
& E_{v_{y+1 / 2}}\left[e^{(i t / \sqrt{m}) \sum_{x=1}^{m} \eta(x)-y-1 / 2}\right] \\
& \quad=\pi_{1}(y+1 / 2) P(i t / \sqrt{m})^{m} \mathbf{1} \\
& \quad=\left(1-t^{2} w_{1}^{\prime \prime}(0) /(2 m)\right)^{m}\left(1+O\left(m^{-1 / 2}\right)\right)
\end{aligned}
$$

must converge to $e^{-t^{2} \sigma^{2}(y+1 / 2) / 2}$. Here, $\pi_{1}(y+1 / 2)$ is thought of as a row vector, and $\mathbf{1}$ is the column vector with entries equal to 1 .

Putting these estimates together, we may conclude (5.3). To verify (5.4), from equation (5.5), we need only show the moduli $\left|w_{1}(i t)\right|,\left|w_{2}(i t)\right|<1$ uniformly for $r \leq|t| \leq \pi$ and $|y| \leq \delta$. One way is the following. Suppose $y=0$ and note that the moduli are less than 1 at $|t|=\pi$. For $r \leq|t| \leq \pi$, from the determinant of
$\widetilde{P}(i t), w_{1}(i t) w_{2}(i t)=1-u_{1}-v_{1}$. In particular, if $w_{1}(i t)$ say is of the form $e^{i \theta}$ with $|\theta| \leq \pi$, then $w_{2}(i t)=e^{-i \theta}\left(1-u_{1}-v_{1}\right)$. From the trace, we obtain equation $e^{i \theta}+e^{-i \theta}\left(1-u_{1}-v_{1}\right)=\left(1-u_{1}\right) e^{-i t / 2}+\left(1-v_{1}\right) e^{i t / 2}$ which is absurd: The real part is $\cos (\theta)\left(2-u_{1}-v_{2}\right)=\cos (t / 2)\left(2-u_{1}-v_{2}\right)$ which yields $\theta=t / 2$. But the imaginary part is $\sin (\theta)\left(u_{1}+v_{1}\right)=\sin (t / 2)\left(u_{1}-v_{1}\right)$ which is a contradiction as $r / 2<|\theta|=|t| / 2 \leq \pi / 2$ and $v_{1} \neq 0$ for $y=0$. Hence, by continuity, for $|y| \leq \delta$ small, we conclude the claim.

Step 3. Now, write

$$
\begin{aligned}
\theta_{n, y, a, b}^{\chi}(x) & =\frac{\sqrt{2 n^{\prime}}}{2 \pi} \int_{-\pi}^{\pi} e^{-i t x} \psi_{n, y, \chi, a, b}(t) d t \\
& =\frac{1}{2 \pi} \int_{-\pi \sqrt{2 n^{\prime}}}^{\pi \sqrt{2 n^{\prime}}} e^{-i t x / \sqrt{2 n^{\prime}}} \psi_{n, y, \chi, a, b}\left(t / \sqrt{2 n^{\prime}}\right) d t
\end{aligned}
$$

The last expression is rewritten as

$$
\begin{aligned}
& \frac{1}{2 \pi} \int_{-\pi \sqrt{2 n^{\prime}}}^{\pi \sqrt{2 n^{\prime}}} \psi_{n, y, \chi, a, b}\left(t / \sqrt{2 n^{\prime}}\right) d t-\frac{i x}{2 \pi \sqrt{2 n^{\prime}}} \int_{-\pi \sqrt{2 n^{\prime}}}^{\pi \sqrt{2 n^{\prime}}} t \psi_{n, y, \chi, a, b}\left(t / \sqrt{2 n^{\prime}}\right) d t \\
& \quad-\frac{x^{2}}{4 \pi n^{\prime}} \int_{-\pi \sqrt{2 n^{\prime}}}^{\pi \sqrt{2 n^{\prime}}} t^{2} \psi_{n, y, \chi, a, b}\left(t / \sqrt{2 n^{\prime}}\right) d t+r_{0}(x) n^{-3 / 2}
\end{aligned}
$$

in terms of error $r_{0}(x)$ which, by the estimates in step 2 , is of order $O\left(|x|^{3}\right)$.
The second integral in the last display is also estimated of order $O\left(n^{\prime-1 / 2}\right)$ by the same argument as given in step 3 of the proof of Proposition 5.1.

Step 4. Then, for $|y| \leq \delta$, we have

$$
\begin{aligned}
E_{\nu_{1 / 2}}[f & f y, \eta(-n-1)=a, \eta(n+1)=b] \\
= & \kappa_{0} E_{v_{y+1 / 2}}[f(\eta) \mid \eta(-n-1)=a, \eta(n+1)=b] \\
& +\frac{\kappa_{1}}{\sqrt{2 n^{\prime}}} E_{v_{y+1 / 2}}\left[\left.f(\eta)\left(\sum_{|x| \leq n^{x}} \eta(x)-\frac{1}{2}-y\right) \right\rvert\, \eta(-n-1)=a\right. \\
& \eta(n+1)=b] \\
& +\frac{\kappa_{2}}{2 n^{\prime}} E_{v_{y+1 / 2}}\left[\left.f(\eta)\left(\sum_{|x| \leq n^{x}} \eta(x)-\frac{1}{2}-y\right)^{2} \right\rvert\, \eta(-n-1)=a,\right. \\
& \\
& \quad \eta(n+1)=b] \\
& \\
& \\
&
\end{aligned}
$$

where $\left|\varepsilon_{f}(n)\right| \leq C\|f\|_{L^{\infty}\left(\nu_{\rho}\right)} n^{-3 / 2+3 \chi}$ and $\kappa_{i}=\kappa_{i}(n)$ for $i=0,1,2$ have the same asymptotics as in step 5 of the proof of Proposition 5.1.

Step 5. Recall the tilted measures and the formula for the tilt $\lambda(z)$ in Section 2.5. Recall also the definitions of $\varphi_{f}^{(i)}(\rho)$ (2.4). For $|y| \leq \delta$ and $i=0,1,2$, using the uniform exponentially mixing property of the measures $\left\{v_{y+1 / 2}:|y| \leq \delta\right\}$ and $\varphi_{f}(1 / 2)=\varphi_{f}^{\prime}(1 / 2)=0$, we claim

$$
\begin{align*}
E_{v_{y+1 / 2}} & {\left[f(\eta)\left(\sum_{|x| \leq n^{x}} \eta(x)-1 / 2-y\right)^{i} \mid \eta(-n-1)=a, \eta(n+1)=b\right] } \\
= & \frac{\lambda^{\prime}(1 / 2)^{2-i} y^{2-i}}{(2-i)!} E_{v_{1 / 2}}\left[f(\eta)\left(\sum_{|x| \leq n^{2} x} \eta(x)-1 / 2\right)^{2}\right]  \tag{5.6}\\
& +|y|^{3-i} r_{1}(f, n)+r_{2}(f, n)
\end{align*}
$$

Here, the error $r_{1}(f, n)$ stands for the error made first in Taylor approximation around $y=0$ with respect to the conditioned measure: Using that $\nu_{y+1 / 2}$ is exponentially mixing, one can bound the first, second and third derivatives below (5.9), (5.10) and (5.11), uniformly in $a, b$ and $|y| \leq \delta$ after a calculation so that $\left|r_{1}(f, n)\right| \leq C(\delta) n^{4 \chi}\|f\|_{L^{\infty}}$. The error $r_{2}(f, n)$ represents other errors made by exponential approximations and $\left|r_{2}(h, n)\right| \leq C\|f\|_{L^{\infty} n^{-3 / 2}}$. The reader, on first reading, may like to skip now to step 6.

Indeed, in more detail, when $i=1$,

$$
\begin{align*}
& E_{v_{y+1 / 2}}\left[f(\eta)\left(\sum_{|x| \leq n^{x}} \eta(x)-1 / 2-y\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right] \\
& =E_{v_{1 / 2}}\left[f(\eta)\left(\sum_{|x| \leq n^{x}} \eta(x)-1 / 2\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right]  \tag{5.7}\\
& \quad+B y+y^{2} r_{1}(f, n),
\end{align*}
$$

where, referring to the first derivative expression (5.9), $B$ equals

$$
\begin{array}{r}
\lambda^{\prime}\left(\frac{1}{2}\right) E_{\nu_{1 / 2}}\left[\left.f(\eta)\left(\sum_{|x| \leq n x} \eta(x)-\frac{1}{2}\right)\left(\sum_{|x| \leq n} \tilde{\eta}(x)\right) \right\rvert\,\right. \\
\eta(-n-1)=a, \eta(n+1)=b]  \tag{5.8}\\
-2 n^{\chi} E_{\nu_{1 / 2}}[f \mid \eta(-n-1)=a, \eta(n+1)=b]
\end{array}
$$

and $\widetilde{\eta}(x)=\eta(x)-E_{\nu_{1 / 2}}[\eta(x) \mid \eta(-n-1)=a, \eta(n+1)=b]$. The error $r_{1}(f, n)$ is less than the bound on the second derivative (5.10) with $h=f(\eta)\left(\sum_{|x| \leq n \times} \eta(x)-\right.$ $1 / 2$ ) plus $2 n^{\chi}$ times the bound on the first derivative (5.9) with $h=\bar{f}$. We now bound the second derivative; estimating the first derivative is similar. See notation $\bar{h}$ and $\bar{\eta}(x)$ above (5.9).

For $|y| \leq \delta$, from the formula for the tilt $\lambda(z)$ in Section 2.5, the derivatives $\lambda^{(k)}(y+1 / 2)$ for $k=1,2,3$ are uniformly bounded. The expectation $E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right]$ in (5.10) is handled as follows. By splitting the sum $\sum_{|x| \leq n} \bar{\eta}(x)$ over indices $|x| \leq n^{2 \chi}, n^{2 \chi}<|x| \leq$ $n-n^{\chi}$ and $|x|>n-n^{\chi}$, and using the uniform exponentially mixing property of $v_{y+1 / 2}$, for $|y| \leq \delta$, and that all variables $|\eta(x)| \leq 1$, one bounds this term as $O\left(\|f\|_{L^{\infty} n^{3 \chi}}\right)$.

Consider now the other term $E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right)^{2} \mid \eta(-n-1)=a, \eta(n+\right.$ $1)=b$ in (5.10). Split the sum over $|x| \leq n$ into sums over $|x| \leq n^{(1+u) \chi}$ and $|x|>n^{(1+u) \chi}$, and square to yield three terms. Bounding the cross term is the most involved, the other two being straightforward. The cross term is

$$
\begin{array}{r}
2 E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n^{(1+u) x}} \bar{\eta}(x)\right)\left(\sum_{n^{(1+u) x}<|x| \leq n} \bar{\eta}(x)\right) \mid\right. \\
\eta(-n-1)=a, \eta(n+1)=b] .
\end{array}
$$

By splitting the sum over $n^{(1+u) \chi}<|x| \leq n$ into sums on $n^{(1+u) \chi}<|x|<$ $n^{(1+2 u) \chi}, n^{(1+2 u) \chi} \leq|x| \leq n-n^{u \chi}$ and $|x|>n-n^{u \chi}$, and using the exponentially mixing property of $v_{y+1 / 2}$, one can bound the cross term $O\left(\|f\|_{\left.L^{\infty} n^{(3+3 u) \chi}\right)}\right.$ which for $u<1 / 3$ gives the desired error bound.

We now relate terms in (5.7) to $\varphi_{f}^{\prime}(1 / 2)$ and $\varphi_{f}^{\prime \prime}(1 / 2)$ [cf. (2.4)], using the exponentially mixing property. It is straightforward that the difference between the first conditional expectation on the right-hand side of (5.7) and $\lambda^{\prime}(1 / 2)^{-1} \varphi_{f}^{\prime}(1 / 2)=0$ is exponentially close. Also, as $\varphi_{f}^{\prime}(1 / 2)=0$, the first conditional expectation in the expression $B$ in (5.8) is exponentially close to $\left(\lambda^{\prime}(1 / 2)\right)^{-1} \varphi_{f}^{\prime \prime}(1 / 2)$, which in turn is exponentially close to the expectation on the right-hand side of (5.6). The other expectation in $B$ is exponentially small.

The cases $i=0$, 2 with respect to equation (5.6), are argued analogously.
Here, for functions $h$ supported on sites in $\Lambda_{n} x$, and notation

$$
\begin{aligned}
& \bar{h}(\eta)=h(\eta)-E_{v_{y+1 / 2}}[h \mid \eta(-n-1)=a, \eta(n+1)=b] \quad \text { and } \\
& \bar{\eta}(x)=\eta(x)-E_{v_{y+1 / 2}}[\eta(x) \mid \eta(-n-1)=a, \eta(n+1)=b]
\end{aligned}
$$

the first derivative is

$$
\begin{align*}
& \frac{d}{d y} E_{v_{y+1 / 2}}[h(\eta) \mid \eta(-n-1)=a, \eta(n+1)=b] \\
& =\lambda^{\prime}\left(y+\frac{1}{2}\right)  \tag{5.9}\\
& \quad \times E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right]
\end{align*}
$$

The second derivative is

$$
\begin{align*}
& \frac{d^{2}}{d y^{2}} E_{v_{y+1 / 2}}[h(\eta) \mid \eta(-n-1)=a, \eta(n+1)=b] \\
& =\lambda^{\prime \prime}\left(y+\frac{1}{2}\right) E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right] \\
&  \tag{5.10}\\
& \quad+\left(\lambda^{\prime}\left(y+\frac{1}{2}\right)\right)^{2} \\
& \quad \\
& \quad \times E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right)^{2} \mid \eta(-n-1)=a, \eta(n+1)=b\right] .
\end{align*}
$$

The third derivative is

$$
\begin{aligned}
\frac{d^{3}}{d y^{3}} E_{v_{y+1 / 2}} & {[h(\eta) \mid \eta(-n-1)=a, \eta(n+1)=b] } \\
= & \lambda^{\prime \prime \prime}\left(y+\frac{1}{2}\right) E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right] \\
+ & 3 \lambda^{\prime}\left(y+\frac{1}{2}\right) \lambda^{\prime \prime}\left(y+\frac{1}{2}\right) \\
& \times E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right)^{2} \mid \eta(-n-1)=a, \eta(n+1)=b\right] \\
+ & \left(\lambda^{\prime}\left(y+\frac{1}{2}\right)\right)^{3} \\
& \times E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right)^{3} \mid \eta(-n-1)=a, \eta(n+1)=b\right] \\
- & 3\left(\lambda^{\prime}\left(y+\frac{1}{2}\right)\right)^{3} \\
& \times E_{v_{y+1 / 2}}\left[\bar{h}(\eta)\left(\sum_{|x| \leq n} \bar{\eta}(x)\right) \mid \eta(-n-1)=a, \eta(n+1)=b\right] \\
& \times E_{v_{y+1 / 2}}\left[\left(\sum_{|x| \leq n} \bar{\eta}(x)\right)^{2} \mid \eta(-n-1)=a, \eta(n+1)=b\right] .
\end{aligned}
$$

Step 6. By the exponentially mixing property of $\nu_{1 / 2}$ and the assumption $\varphi_{f}(1 / 2)=\varphi_{f}^{\prime}(1 / 2)=0$ [cf. (2.4)], we have

$$
\lambda^{\prime}(1 / 2)^{2} E_{\nu_{1 / 2}}\left[f(\eta)\left(\sum_{|x| \leq n^{2} x} \eta(x)-1 / 2\right)^{2}\right]=\varphi_{f}^{\prime \prime}(1 / 2)+O\left(n^{-3 / 2}\right)
$$

Also, note the relation $\lambda^{\prime}(1 / 2) \sigma^{2}(1 / 2)=1$ (cf. Section 2.5 ), and by exponential mixing that $\left|\sigma^{2}(1 / 2)-\sigma_{n}^{2}(1 / 2)\right|=O\left(n^{-1}\right)$. Recall the asymptotic behaviors of $\kappa_{0}, \kappa_{1}$ and $\kappa_{2}$ (cf. step 5 of proof of Proposition 5.1). In addition, a factor $n^{4 \chi} y$ is of order $O\left(n^{-(1 / 2-4 \chi)}\right)$ in $L^{4}\left(\nu_{1 / 2}\right)$. Hence, with the parameter $\chi$ chosen small enough, dominant terms may be gathered, as done in the proof of Proposition 5.1, to obtain for all large $n$ that

$$
\begin{aligned}
& \| 1(|y| \leq \delta)\left(E_{\nu_{1 / 2}}[f \mid y, \eta(-n-1)=a, \eta(n+1)=b]\right. \\
& \left.\quad-\frac{\varphi_{f}^{\prime \prime}(1 / 2)}{2}\left[y^{2}-\frac{\sigma_{n}^{2}(1 / 2)}{2 n+1}\right]\right) \|_{L^{4}\left(\nu_{1 / 2}\right)} \\
& \leq C\|f\|_{L^{\infty} n^{-3 / 2+\varepsilon}} .
\end{aligned}
$$

On the other hand, large deviation estimates yield

$$
\begin{aligned}
& \| 1(|y|>\delta)\left(E_{\nu_{1 / 2}}[f \mid y, \eta(-n-1)=a, \eta(n+1)=b]\right. \\
& \left.\quad-\frac{\varphi_{f}^{\prime \prime}(1 / 2)}{2}\left[y^{2}-\frac{\sigma_{n}^{2}(1 / 2)}{2 n+1}\right]\right) \|_{L^{4}\left(\nu_{1 / 2}\right)} \\
& \leq C\|f\|_{L^{\infty} n^{-3 / 2}}
\end{aligned}
$$

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