# Relaxed Criteria of the Dobrushin-Shlosman Mixing Condition 

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#### Abstract

An interacting particle system (Glauber dynamics) which evolves on a finite subset in the $d$-dimensional integer lattice is considered. It is known that a mixing property of the Gibbs state in the sense of Dobrushin and Shlosman is equivalent to several very strong estimates in terms of the Glauber dynamics. We show that similar, but seemingly much milder estimates are again equivalent to the Dobrushin-Shlosman mixing condition, hence to the original ones found by Stroock and Zegarlinski. This may be understood as the absence of intermediate speed of convergence to equilibrium.


KEY WORDS: Relaxed criteria: the Dobrushin-Shlosman mixing condition; Glauber dynamics; convergence to equilibrium.

## 1. INTRODUCTION

A standard way to describe a mixing property of a Gibbs state is to estimate the difference of the expectations of a local observable with respect to a finite-volume Gibbs state with different boundary conditions. For example, the Dobrushin-Shlosman mixing condition which we will discuss can roughly be stated as follows: the difference mentioned above is exponentially small in the distance between the support of the observable and the sites at which the boundary conditions are different [cf. (2.16)]. On the other hand, since the mixing property of a Gibbs state reflects the rapid relaxation of the Glauber dynamics, the mixing property can be expressed in the following different ways in terms of the Glauber dynamics: (a) estimate of the logarithmic Sobolev constant, (b) estimate of the spectral gap, (c) estimate of the rate of convergence (the difference between the

[^0]semigroup at time $t$ and its equilibrium). Stroock and Zegarlinski ${ }^{(16)}$ succeeded in rephrasing the Dobrushin-Shlosman mixing condition in three ways described above, which we will review as Theorem 3.1 below.

The purpose of this paper is to present relaxed criteria of the DobrushinShlosman mixing condition of types $(a)-(c)$, which are respectively conditions ( 2 a )-(2c) in Theorem 3.2, our main result. Technically, the derivation of the original Dobrushin-Shlosman mixing condition from these relaxed criteria is based on the fact that the exponential decay in the statement of the Dobrushin-Shlosman mixing condition is equivalent to a certain polynomial decay. This point is also made clear as conditions (2d) and (2e) in Theorem 3.2. These relaxed criteria are potentially useful to check the Dobrushin-Shlosman mixing condition, which is often accepted as an assumption to do something with, but is not always easy to verify in practical applications. Since our result is based on the very old and well-known idea that "an exponential mixing follows from a certain polynomial mixing," there are many results with the same sprit in the earlier literature. The relation between these earlier results and ours will be discussed in the series of remarks following the statement of Theorem 3.2.

## 2. BASIC DEFINITIONS

The lattice. We will work on the $d$-dimensional integer lattice $\mathbf{Z}^{d}=$ $\left\{x=\left(x_{i}\right)_{i=1}^{\prime \prime}: x_{i} \in \mathbb{Z}\right\}$, on which we consider the $l^{\prime \prime}$ norm; $|x|=\max _{1 \leqslant i \leqslant d}\left|x_{i}\right|$. For a set $A \subset \mathbf{Z}^{\prime}$, $\operatorname{diam} A$ and $|A|$ stand, respectively, for its diameter and the number of the points it contains. We write $A \subset \subset \mathbf{Z}^{d}$ when $1 \leqslant|A|<\infty$ and define a family $\mathscr{A}$ by

$$
\begin{equation*}
\mathscr{A}=\left\{A ; A \subset \subset \mathbf{Z}^{\prime}\right\} \tag{2.1}
\end{equation*}
$$

The distance between two subsets $\Lambda_{1}$ and $\Lambda_{2}$ of $\mathbf{Z}^{d}$ will be denoted by $d\left(\Lambda_{1}, A_{2}\right)$. For $r \geqslant 1$, the $r$-boundary of a set $\bar{A}$ is defined by

$$
\begin{equation*}
\partial_{r} A=\{x \notin A ; d(x, A) \leqslant r\} \tag{2.2}
\end{equation*}
$$

The value of $r$ will eventually be chosen as an upper bound $r_{0}$ of the range of the interaction we consider [cf. (2.5) below]. For $v \in \mathbf{Z}^{d}$ and an integer $m \geqslant 1$, we define a subset $\mathscr{B}_{r}(m)$ of $\mathscr{A}$ as the totality of $A \subset \subset \mathbf{Z}^{d}$ of the following form:

$$
\begin{equation*}
A=\left\{x \in \mathbf{Z}^{d} ; v_{i}+m a_{i} \leqslant x_{i}<v_{i}+m b_{i} \text { for } i=1, \ldots, d\right\} \tag{2.3}
\end{equation*}
$$

where $a_{i}$ and $b_{i}$ are integers with $a_{i}<b_{i}$. In particular, the class $\mathscr{B}_{B_{1}}(m)$ consists of boxes with every sidelength a positive multiple of $m$. We define
a subset $\mathscr{C}_{r}(m)$ of $\mathscr{B}_{r}(m)$ as the totality of cubes in $\mathscr{B}_{r}(m)$, i.e., the totality of $A \subset \subset \mathbf{Z}^{d}$ of the form (2.3) with $b_{i}-a_{i}(1 \leqslant i \leqslant d)$ identical.

The configurations. We take a finite set $S$ as a single spin space and $\lambda$ will stand for the uniform distribution on $S ; \lambda(d t)=(1 /|S|) \sum_{s \in S} \delta_{s}(d t)$. Configuration spaces are defined as follows:

$$
\begin{aligned}
\Omega_{A} & =\left\{\sigma=\left(\sigma_{x}\right)_{x \in A} ; \sigma_{x} \in S\right\}, \quad \Lambda \subset \mathbf{Z}^{d} \\
\Omega & =\Omega_{\mathbf{Z}^{d}}
\end{aligned}
$$

For $\Lambda \subset \mathbf{Z}^{d}$ and $(\sigma, \omega) \in \Omega^{2}, \sigma_{.1} \cdot \omega_{. .^{\prime}}$ denotes the following configuration:

$$
\left(\sigma_{A} \cdot \omega_{.1}\right)_{x}=\left\{\begin{array}{lll}
\sigma_{x} & \text { if } & x \in \Lambda \\
\omega_{x} & \text { if } & x \notin \Lambda
\end{array}\right.
$$

For $f: \Omega_{A} \rightarrow \mathbf{R}$ we introduce the notations

$$
\begin{aligned}
\nabla_{x} f(\sigma) & =\int_{S^{x}} f\left(\xi_{x} \cdot \sigma_{x^{\prime}}\right) \lambda^{x}\left(d \xi_{x}\right)-f(\sigma), \quad X \subset \subset \mathbf{Z}^{d} \\
\|f\| & =\sup _{\sigma \in \Omega_{1}}|f(\sigma)| \\
\|\|f\| & =\sum_{x \in A}\left\|\nabla_{x} f\right\| \\
\operatorname{osc}_{x}(f) & =\sup _{\left(\pi, \sigma^{\prime}\right) \in \Omega^{2}}\left\{\left|f(\sigma)-f\left(\sigma^{\prime}\right)\right| ; \sigma \equiv \sigma^{\prime} \text { off } x\right\}
\end{aligned}
$$

$$
\Delta_{f}=\left\{x \in A ; f \text { is not a constant with respect to } \sigma_{x}\right\}
$$

The function spaces $\mathscr{C}$ and $\mathscr{C}_{1}\left(A \subset \mathbf{Z}^{d}\right)$ are defined, respectively, by

$$
\begin{aligned}
\mathscr{C} & =\left\{f: \Omega \rightarrow \mathbf{R} ;\left|\Delta_{j}\right|<\infty\right\} \\
\mathscr{C}_{A} & =\left\{f: \Omega \rightarrow \mathbf{R} ; \Delta_{j} \subset A\right\}
\end{aligned}
$$

The interaction and finite-volume Gibbs states. A family $\boldsymbol{\Phi}=$ $\left\{\Phi_{X} \in \mathscr{C}_{X}: X \subset \subset \mathbf{Z}^{d}\right\}$ is called a bounded, finite-range interaction if it satisfies the following:
(Ф-1) There exists $M_{0}<\infty$ such that

$$
\begin{equation*}
\|\Phi\|:=\sup _{x \in \mathbb{Z}^{\prime}} \sum_{X: X \ni: X}\left\|\Phi_{X}\right\| \leqslant M_{0} \tag{2.4}
\end{equation*}
$$

(Ф-2) There exists $r_{0}<\infty$ such that

$$
\begin{equation*}
r(\boldsymbol{\Phi}):=\sup \left\{\operatorname{diam}(X) ; \Phi_{X} \not \equiv 0\right\} \leqslant r_{0} \tag{2.5}
\end{equation*}
$$

$\|\Phi\|$ in (2.4) and $r(\Phi)$ in (2.5) are called the norm and the range of the interaction, respectively.

From here on we fix a bounded, finite-range interaction $\boldsymbol{\Phi}$. For each $\Lambda \subset \subset \mathbf{Z}^{d}$ we define the Hamiltonian $H_{A} \in \mathscr{C}$ by

$$
\begin{equation*}
H_{A}=\sum_{x: X \cap A \neq \varnothing} \Phi_{X} \tag{2.6}
\end{equation*}
$$

For each $A \subset \subset \mathbf{Z}^{d}$ and $\omega \in \Omega$ we define the finite-volume Gibbs state $\mu^{\text {A.(") }}$ as the probability measure on $\Omega_{A}$ in which each configuration $\sigma_{A} \in \Omega_{A}$ appears with probability

$$
\begin{equation*}
\mu^{A \cdot(")}\left(\left\{\sigma_{A}\right\}\right)=\frac{\exp -H_{A}\left(\sigma_{A} \cdot \omega_{A c}\right)}{Z^{A(\omega)}} \tag{2.7}
\end{equation*}
$$

where $Z^{\Lambda, \prime \prime \prime}$ is the normalizing constant.
The stochastic dynamics. We introduce now for the model above the time evolution called Glauber dynamics. We define for each $x \in \mathbf{Z}^{d}$ an operator $A_{x}: \mathscr{C} \rightarrow \mathscr{C}$ by

$$
A_{*} f(\sigma)=\sum_{s \in S} c_{*}(\sigma, s)\left(f\left(\sigma \cdot{ }^{*} s\right)-f(\sigma)\right)
$$

where $\sigma{ }^{*} s$ is a configuration obtained from $\sigma$ by replacing $\sigma_{x}$ by $s$ and the coefficient $c_{x}(\sigma, s)$, which is called the flip rate, is required to satisfy the following.
(R-1) Boundedness: There exist positive constants $\underline{c}$ and $\bar{c}$ such that

$$
\underline{c} \leqslant c_{.}(\sigma, s) \leqslant \bar{c} \quad \text { for all } \quad x \in \mathbf{Z}^{d}, \quad \sigma \in \Omega, \quad s \in S
$$

(R-2) Finite range: There exists $r_{1} \geqslant 0$ such that

$$
\begin{equation*}
\gamma(x, y)=0 \quad \text { if } \quad|x-y|>r_{1} \tag{2.8}
\end{equation*}
$$

where

$$
\begin{equation*}
\gamma(x, y)=\sup _{\left(\sigma, \sigma^{\prime}\right) \in \Omega^{2}}\left\{\sum_{s \in S}\left|c_{N}(\sigma, s)-c_{v}\left(\sigma^{\prime}, s\right)\right| ; \sigma \equiv \sigma^{\prime} \text { off } y\right\} \tag{2.9}
\end{equation*}
$$

(R-3) Detailed balance condition: For all $x \in \mathbf{Z}^{d}, \sigma \in \Omega$, and $\left(s, s^{\prime}\right) \in S^{2}$,

$$
\begin{equation*}
\mu^{\{x\} \cdot \sigma}(\{s\}) c_{x}\left(\sigma^{x} s, s^{\prime}\right)=\mu^{\{x\}, \sigma}\left(\left\{s^{\prime}\right\}\right) c_{x}\left(\sigma^{x} s^{\prime}, s\right) \tag{2.10}
\end{equation*}
$$

Remark 2.1. This is a very technical remark and is relevant only when one would like to identify a numerical constant $M_{1}$ which will appear in Theorem 3.2, Lemma 4.1, and Lemma 4.2. We set

$$
\begin{align*}
M & =\sup _{x} \sum_{,: y \neq x} \gamma(x, y) \\
M_{1} & =\sum_{y \in \in \mathbb{Z}^{d}} \gamma_{1}(y) \\
& =\left(2 r_{1}+1\right)^{d}|S|(\bar{c}-\underline{c}) \tag{2.11}
\end{align*}
$$

where $\gamma(\cdot, \cdot)$ is defined by (2.9) and

$$
\gamma_{1}(y)= \begin{cases}|S|(\bar{c}-\underline{c}) & \text { if }|y| \leqslant r_{1}  \tag{2.12}\\ 0 & \text { otherwise }\end{cases}
$$

The function $\gamma(\cdot, \cdot)$, together with the constant $M$ plays fundamental role in analyzing a particle system (ref. 10, Chapter I). The function $\gamma_{1}(\cdot)$ is introduced as a shift-invariant object which bounds $\gamma(\cdot, \cdot)$ from above. Note that $\gamma(x, y) \leqslant \gamma_{1}(x-y)$ and hence $M \leqslant M_{1}$.

For fixed $A \subset \subset \mathbf{Z}^{d}$ and $\omega \in \Omega$ we define an operator $A^{A, \omega}: \mathscr{C}_{A} \rightarrow \mathscr{C}_{A}$ and the associated semigroup $\left(T_{1}^{A, 1 ")}\right)_{t \geqslant 0}$ by

$$
\begin{aligned}
A^{A,(\prime \prime} f(\sigma) & =\sum_{x \in A} A_{x} f\left(\sigma_{A} \cdot \omega_{A}\right), \quad f \in \mathscr{C}_{A} \\
T_{A}^{A,(")} & =\exp t A^{A,(\prime)}
\end{aligned}
$$

We then see from (2.10) that for all $\{f, g\} \subset \mathscr{C}_{A}$

$$
\begin{align*}
& \mathscr{E}^{\text {A., (") }}(f, g) \stackrel{\text { der }}{=}-\mu^{A \cdot(")}\left(f A^{A \cdot(\prime \prime \prime} g\right) \\
&= \frac{1}{2} \sum_{x \in A} \sum_{s \in S} \int \mu^{A,(\prime \prime \prime}(d \sigma) c_{x}(\sigma, s)\left[f\left(\sigma \cdot^{*} s\right)-f(\sigma)\right] \\
& \times\left[g\left(\sigma \cdot^{*} s\right)-g(\sigma)\right] \tag{2.13}
\end{align*}
$$

We have defined $A^{A^{\prime}(\prime)}$ and $T_{1}^{A_{1},(,)}$ as operators acting on $\mathscr{C}_{A}$. But we will extend their domain of definition to $\mathscr{C}$ by applying them to $" \sigma \mapsto f\left(\sigma_{A} \cdot \omega_{A}\right)$."

The spectral gap, the logarithmic Sobolev constant, and the DobrushinShlosman mixing condition. We now introduce a couple of quantities and a notion which plays crucial roles in analyzing the long-time behavior of the Glauber dynamics. We define the inverse spectral $\operatorname{gap} \gamma_{S G}(\lambda, \omega)$ as the smallest $\gamma$ for which the following inequality is true for all $f \in \mathscr{C}$ :

Similarly, we define the logarithmic Sobolev constant $\gamma_{\mathrm{LS}}(\Lambda, \omega)$ as the smallest $\gamma$ for which the following inequality is true for all $f \in \mathscr{C}$ :

$$
\begin{equation*}
\mu^{\operatorname{Lc\omega }}\left(f^{2} \log \frac{f^{2}}{\mu^{1+\prime \prime \prime}\left(f^{2}\right)}\right) \leqslant 2 \gamma^{\mathscr{E}}{ }^{\text {A/w }}(f, f) \tag{2.15}
\end{equation*}
$$

For $\mathscr{F} \subset \mathscr{A}$ an interaction $\Phi$ is said to satisfy the Dobrushin-Shlosman mixing condition over $\mathscr{F}$ if there exist constant $C_{i} \in(0, \infty)(i=1,2)$ such that for all $\Lambda \in \mathscr{F}, y \notin \Lambda$, and $f \in \mathscr{C}_{.1}$,

$$
\begin{equation*}
\operatorname{osc}_{4}\left(\mu^{t \cdot} f\right) \leqslant C_{1}\|f\| \exp -\frac{d\left(y, \Delta_{f}\right)}{C_{2}} \tag{2.16}
\end{equation*}
$$

In the sequel we will refer to the above condition as $\operatorname{DSM}(\mathscr{F})$.
It is convenient to introduce the following definition. We call a family $\{c(f)>0\}_{f \in \%}$ an admissible coefficient if it satisfies

$$
\begin{equation*}
\sup \left\{c(f) ;\|f\|+\operatorname{diam} \Delta_{f} \leqslant m\right\}<\infty \quad \text { for all } \quad m>0 \tag{2.17}
\end{equation*}
$$

What (2.17) requires is that $c(f)$ should have an upper bound in terms of $\left\|f^{\prime}\right\|$ and $\operatorname{diam} \Delta_{f}$, which are independent of where $\Delta_{f}$ is. Typically, const. $\|f\|$ or const. $\|\|f\|\|$ appears as the admissible coefficient $c(f)$.

Remark 2.2. We could have worked in the continuous spin setting, in which the discrete spin space $S$ is replaced by a smooth, compact Riemannian manifold (in fact, $C^{2}$ is enough). What we have defined above can easily be modified so that they make sense in the continuous spin setting and the all results as well as their proofs in this paper remain valid up to the value of constants and boring technicalities.

## 3. THE RESULT

First we recall the following result:
Theorem 3.1. ${ }^{(16)}$ Let $\mathscr{F}$ be either $\mathscr{A}$ or $\mathscr{B}_{n}(m)$ for arbitrarily fixed $v \in \mathbf{Z}^{d}$ and $m \geqslant 1$ [cf. (2.1), (2.3)]. Then each of the following conditions is equivalent to the $\operatorname{DSM}(\mathscr{F})$ [cf. (2.16)]:
(1a) The logarithmic Sobolev constant [cf. (2.15)] satisfies

$$
\begin{equation*}
\sup _{\epsilon \in \mathcal{F}, \omega \in \Omega} \gamma_{\mathrm{Ls}}(\Lambda, \omega)<\infty \tag{3.1}
\end{equation*}
$$

(1b) The inverse spectral gap [cf. (2.14)] satisfies

$$
\begin{equation*}
\sup _{\Lambda \in \mathscr{F},(\omega \in \Omega} \gamma_{\mathrm{sG}}(\Lambda, \omega)<\infty \tag{3.2}
\end{equation*}
$$

(1c) There exist a constant $C>0$ and an admissible coefficient $\{c(f)>0\}_{f \in \%}[c f .(2.17)]$ such that for all $f \in \mathscr{C}$ and $t>0$

$$
\begin{equation*}
\sup _{A \in \sqrt{F},(1) \in \Omega}\left\|T_{1}^{A \cdot(\prime)} f-\mu^{A \cdot(\prime \prime \prime} f\right\| \leqslant c(f) \exp -\frac{t}{C} \tag{3.3}
\end{equation*}
$$

Remark 3.1. In the original statement of this theorem [ref. 16, Theorem 1.8, part (c)] only the family $\mathscr{A}$ is under consideration. The statements for $\mathscr{B}_{r}(m)$ is implicit in the argument in that paper. The point is that both $\mathscr{A}$ and $\mathscr{B}_{r}(m)$ are closed under intersection and contain arbitrarily large cubes. The DSM condition restricted to boxes is easier to check than $\operatorname{DSM}(\mathscr{A})$ is. In fact, $\operatorname{DSM}(\mathscr{A})$ is known to be true for the Ising ferromagnet with $\beta<\beta_{c} / 2$ or $|h|>2 d$ [ref. 5; see also ref. 3, (2.32) for the latter case], while the validity of $\operatorname{DSM}\left(\mathscr{B}_{r}(2)\right)$ extends to $|h|>d-1$. As mentioned in ref. 6, Remark 2.3, there is also a difference between $\operatorname{DSM}\left(\mathscr{B}_{r}(1)\right)$ and $\operatorname{DSM}\left(\mathscr{B}_{r}(m)\right)$ with a large $m$ if $d \geqslant 3$. The DSM condition based on cubes, together with an analogous statement to Theorem 3.1 in that context, is studied extensively in refs. 11-13 and 14. In this connection, it should be mentioned that ref. 11, Theorem 4.1, contains the following nice observation: if $\operatorname{DSM}\left(\bigcup_{r \in \mathbb{Z}^{d}} \bigcup_{m \geqslant m_{n}} \mathscr{C}_{r}(m)\right)$ holds for some $m_{0} \geqslant 1$, then $\operatorname{DSM}\left(\bigcup_{v \in \mathbf{Z}^{i}} \bigcup_{m \geqslant m_{1}} \mathscr{B}_{r}(m)\right)$ holds for large enough $m_{1} \geqslant 1$.

Next, we state the main result of this paper.
Theorem 3.2. Let $\mathscr{F}$ be either $\mathscr{A}$ or $\mathscr{B}_{r}(m)$ for arbitrarily fixed $v \in \mathbf{Z}^{d}$ and $m \geqslant 1$ [cf. (2.1), (2.3)]. Then each of the following conditions is equivalent to $\operatorname{DSM}(\mathscr{F})$ [cf. (2.16)]:
(2a) The logarithmic Sobolev constant [cf. (2.15)] satisfies

$$
\begin{equation*}
\varlimsup_{D \rightarrow \infty} \sup _{A \in \mathscr{F}, \ldots \in \Omega}\left\{\frac{\gamma_{\mathrm{Ls}}(\Lambda, \omega)}{\psi(\operatorname{diam} \Lambda)} ; \operatorname{diam} A \geqslant D\right\}<\frac{1}{40(d-1) r_{1} M_{1}} \tag{3.4}
\end{equation*}
$$

where $\psi(t)=t / \log t$. The constants $r_{1}$ and $M_{1}$ are given, respectively, by (2.8) and (2.11).
(2b) The inverse spectral gap [cf. (2.14)] satisfies

$$
\begin{equation*}
\varlimsup_{D \rightarrow \chi_{-}} \sup _{A \in \mathcal{F} \cdot(1) \in S}\left\{\frac{\gamma_{\mathrm{SG}}(\Lambda, \omega)}{\psi(\operatorname{diam} A)} ; \operatorname{diam} A \geqslant D\right\}<\frac{1}{40(d-1) r_{1} M_{1}} \tag{3.5}
\end{equation*}
$$

where $\psi(t)=t / \log t$. The constants $r_{1}$ and $M_{1}$ are given, respectively, by (2.8) and (2.11).
(2c) There exists an admissible coefficient $\{c(f)>0\}_{f \in \%}[c f$. (2.17)] and a nonincreasing function $\phi:(0, \infty) \rightarrow(0, \infty)$ such that $\phi(t)=$ $o\left(t^{-2(d-1)}\right)$ as $t \rightarrow \infty$ and such that

$$
\begin{equation*}
\sup _{A \in \tilde{F}, \ldots \in S 2} \mu^{A \cdot \omega}\left|T_{1}^{A \cdot \omega} f-\mu^{A \cdot \omega} f\right| \leqslant c(f) \phi(t) \quad \text { for all } f \in \mathscr{C} \tag{3.6}
\end{equation*}
$$

(2d) There exist $C_{1} \in(2, \infty)$, an admissible coefficient $\{c(f)>0\}_{f \in \%}$ and a nonincreasing function $\phi:(0, \infty) \rightarrow(0, \infty)$ such that $\phi(t)=$ $o\left(t^{-2(d-1)}\right)$ as $t \rightarrow \infty$ and such that

$$
\begin{equation*}
\operatorname{osc}_{y}\left(\mu^{\mu \cdot} f\right) \leqslant c(f) \phi\left(d\left(y, \Delta_{f}\right)\right) \tag{3.7}
\end{equation*}
$$

for all $\Lambda \in \mathscr{\mathscr { Y }}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{0}} \Lambda$ which satisfy $\operatorname{diam} \Lambda \leqslant C_{1} d\left(y, \Delta_{f}\right)$.
(2e) There exists a nonincreasing function $\phi:(0, \infty) \rightarrow(0, \infty)$ such that $\phi(t)=o\left(t^{-(d-1)}\right)$ as $t \rightarrow \infty$ and such that

$$
\begin{equation*}
\operatorname{osc}_{\cdot}\left(\mu^{A \cdot} f\right) \leqslant\|f\| \phi\left(d\left(y, \Delta_{j}\right)\right) \tag{3.8}
\end{equation*}
$$

for all $\Lambda \in \mathscr{F}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{0}} \Lambda$.
Remark 3.2. It follows from Theorem 3.2 that the condition (2a) or even (2b) implies (1c) in Theorem 3.1. This is reminiscent of ref. 8, Theorem 5.3. Although (2a) and (2b) require control over boxes, not only over cubes, the conclusion (1c) is quite rewarding as compared with that of the above-quoted reference. The bound $1 /\left\{40(d-1) r_{1} M_{1}\right\}$ on the righthand sides of (3.4), (3.5) is technical and is in no way sharp.

Remark 3.3. It follows from Theorem 3.2 that (2c) implies (1c), which roughly says the following: if the rate of convergence of a Glauber dynamics is bounded by a certain negative power of the time $t$, then the convergence is necessarily exponentially fast in $t$. As is mentioned in the Introduction, this is not the first time a result of this kind has been obtained. For example, it is shown in ref. 15, Theorem 3.6, that $L^{*}$-convergence as fast as $t^{-(2 d+a)}$ implies exponentially fast $L^{\infty}$-convergence. Also, it is mentioned in ref. 9, Remark 4.3, that $L^{2}$-convergence as fast as $t^{-(2 d+a)}$ implies exponentially fast $L^{2}$-convergence. On the other hand, our
result $[(2 \mathrm{c}) \Rightarrow(1 \mathrm{c})]$ says that $L^{1}$-convergence faster than $t^{-2(d-1)}$ in the sense of (2c) is enough for exponentially fast $L^{\alpha}$-convergence, which improves two results mentioned above. It should also be mentioned that in the case $S=\{-1,+1\}$ and the flip rate is attractive, (1c) is also implied by the following condition:

$$
\begin{equation*}
\lim _{t \rightarrow \infty} t^{d} \sup _{A \in \mathscr{F} \cdot, \ldots \in \Omega} \sup _{x \in A}\left\{\left(T_{1}^{A,(\omega)} \sigma_{x}\right)(+)-\left(T_{1}^{A,(\prime)} \sigma_{x}\right)(-)\right\}=0 \tag{3.9}
\end{equation*}
$$

where ( $\pm$ ) stand for configurations with all spins equal to $\pm 1$. This follows from the proof of ref. 1, Theorem 4 , with a slight adjustment in order to circumvent the lack of shift invariance. In $d=2$ condition (2c) is milder than (3.9) (see ref. 1, Lemma 2.1).

Remark 3.4. Condition (2e) says that the exponential decay in the definition of $\operatorname{DSM}(\mathscr{F})$ is equivalent to a certain polynomial decay. Technically, condition (2d) plays the role of a junction between a statement in terms of dynamics, like ( 2 b ) or (2c), and that in terms of equilibrium, like (2e) (cf. Lemmas 4.2 and 4.3). Let us compare condition (2d) with a similar condition in ref. 4 . In ref. 4 it is proved that $\operatorname{DSM}(\mathscr{F})$ has the following relaxed criterion, which is called condition (IIIb) in that paper: there exists a nonincreasing function $\phi:(0, \infty) \rightarrow(0, \infty)$ such that $\phi(t)=o\left(t^{-2(d-1)}\right)$ as $t \rightarrow \infty$ and such that

$$
\begin{equation*}
\operatorname{osc}_{y}\left(\mu^{A \cdot} \cdot f\right) \leqslant\|f\| \sum_{x \in A_{f}} \phi(|x-y|) \tag{3.10}
\end{equation*}
$$

for all $\Lambda \in \mathscr{F}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{v}} \Lambda$. As can easily be seen, condition (2d) relaxes condition (IIIb) even further.

## 4. PROOF OF THEOREM 3.2

Lemma 4.1. Let $A \subset \subset \mathbf{Z}^{d}$ and $\omega \in \Omega$ be fixed. Suppose that $\phi(t)$ is a nonincreasing function such that $\lim _{t \rightarrow \infty} \phi(t)=0$ and that

$$
\begin{equation*}
\mu^{A,(\prime \prime}\left|T_{1}^{A, \cdots} f-\mu^{A,(\prime \prime} f\right| \leqslant c(f) \phi(t) \quad \text { for all } f \in \mathscr{C} \quad \text { and } \quad t>0 \tag{4.1}
\end{equation*}
$$

with some coefficient $c(f)$. Then there exists $C \in(0, \infty)$ which depends only on the flip rate such that for all $\{f, g\} \subset \mathscr{C}$

$$
\begin{align*}
\mu^{A,(\omega)}(f ; g) & \stackrel{\text { der }}{=} \mu^{A \cdot \omega \prime}(f g)-\mu^{A(1)} f \mu^{A,(\prime \prime \prime} g \\
& \leqslant c(f, g) \phi\left(\frac{d\left(\Delta_{f}, \Delta_{g}\right)}{10 M_{1} r_{1}}\right)+C\|f\| \cdot\|g\| \exp -\frac{d\left(\Delta_{f}, \Delta_{g}\right)}{10 r_{1}} \tag{4.2}
\end{align*}
$$

where $c(f, g)=c(f g)+\|f\| c(g)+c(f)\|g\|$. The constants $r_{1}$ and $M_{1}$ are given respectively by (2.8) and (2.11).

Proof. We first bound the left-hand side of (4.2) by three terms:

$$
\begin{align*}
& \mu^{A \cdot \omega}\left|T_{1}^{A, \omega}(f g)-T_{1}^{A \cdot \omega} f \cdot T_{1}^{A \cdot \omega} g\right| \\
& \leqslant \mu^{A, \omega}\left|T_{1}^{A . \omega \prime}(f g)-T_{s}^{A . c( }(f g)\right|  \tag{4.3}\\
& +\mu^{A, c}\left|T_{s}^{A \cdot c)}(f g)-T_{s}^{A \cdot c \prime} f \cdot T_{s}^{A \cdot c)} g\right|  \tag{4.4}\\
& +\mu^{A, \omega}\left|T_{s}^{A, \omega} f \cdot T_{s}^{A, \omega} g-T_{1}^{A, \omega} f \cdot T_{1}^{A, \omega} g\right| \tag{4.5}
\end{align*}
$$

The condition (4.1) is available to bound the first and the third terms:

$$
\begin{align*}
(4.3) & \leqslant \mu^{A,(1)}\left|T_{1}^{A \cdot(\omega)}(f g)-\mu^{A \cdot \omega}(f g)\right|+\mu^{A \cdot \omega}\left|T_{s}^{A, \omega}(f g)-\mu^{A \cdot \omega}(f g)\right| \\
& \leqslant c(f g)(\phi(t)+\phi(s)) \tag{4.6}
\end{align*}
$$

Similarly,

$$
\begin{equation*}
(4.5) \leqslant(\|f\| c(g)+c(f)\|g\|)(\phi(t)+\phi(s)) \tag{4.7}
\end{equation*}
$$

We now apply Proposition 4.18 of ref. 10 , p. 40 , to (4.4) to see that

$$
\begin{equation*}
(4.4) \leqslant C\|f\| \cdot\|g\| \| \exp \left(4 M_{1} s-\delta \rho\right) \tag{4.8}
\end{equation*}
$$

where $\delta=1 / 2 r_{1}, \rho=d\left(\Delta_{f}, \Delta_{g}\right)$, and $C>0$ is a constant which depends only on the flip rate (cf. Remark 4.1 below). At this point we take $s=\delta \rho / 5 M_{1}$. We then have by (4.6)-(4.8) that

$$
\begin{aligned}
& \mu^{A_{1},(\prime)}\left|T_{1}^{A, w}(f g)-T_{1}^{A \cdot \omega} f \cdot T_{1}^{A \cdot \omega} g\right| \\
& \leqslant \\
& \quad(c(f g)+\|f\| c(g)+c(f)\|g\|)\left(\phi(t)+\phi\left(\delta \rho / 5 M_{1}\right)\right) \\
& \quad+C\|f\| \cdot\|g\| \| \exp -\frac{\delta \rho}{5}
\end{aligned}
$$

which proves (4.2) by letting $t \rightarrow \infty$. QED
Remark 4.1. Since $\left(T_{t}^{A .1 "}\right)_{t \geqslant 0}$ is not a shift-invariant particle system, (4.8) does not come from a direct application of ref. 10, Proposition 4.18, but from a modification of the proof of that proposition, to be precise. The modification is simply to replace $\gamma(x, y)$ by its upper bound $\gamma_{1}(x-y)$ (cf. Remark 2.1). The choice $\delta=1 / 2 r_{1}$ is possible for the following
reason. Going back to the proof of the proposition, $\delta$ has been chosen so that

$$
\begin{equation*}
\sum_{x \in \mathbf{Z}^{d}} \gamma_{1}(x) \exp (\delta|x|) \leqslant 2 M_{1} \tag{4.9}
\end{equation*}
$$

Since the left-hand side of (4.9) is not greater than $\exp \left(\delta r_{1}\right) M_{1}$, our choice $\delta=1 / 2 r_{1}$ is more than enough.

Lemma 4.2. Let $C_{0} \in[0, \infty), C_{1} \in(0, \infty)$, and $\mathscr{F} \subset \mathscr{A}$ be arbitrary. Suppose that there exists $C_{2} \in(0, \infty)$ such that

$$
\begin{equation*}
\overline{\lim _{D \rightarrow \infty}} \sup _{A \in \mathscr{F},(\omega \in \Omega}\left\{\frac{\gamma_{\mathrm{sG}}(A, \omega)}{\psi(\operatorname{diam} \Lambda)} ; \operatorname{diam} \Lambda \geqslant D\right\}<C_{2} \tag{4.10}
\end{equation*}
$$

where $\psi(t)=t / \log t$. Then there exists $C_{3} \in(0, \infty)$ such that

$$
\begin{equation*}
\operatorname{osc}_{y}\left(\mu^{A \cdot} \cdot f\right) \leqslant C_{3}\|f\| d\left(y, \Delta_{f}\right)^{-1 /\left(s c_{1} c_{2} M_{1} r_{1}\right)} \tag{4.11}
\end{equation*}
$$

for all $\Lambda \in \mathscr{F}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{0}} \Lambda$ which satisfy

$$
\begin{equation*}
\operatorname{diam} \Lambda \leqslant C_{1} d\left(y, \Delta_{f}\right)+C_{0} \tag{4.12}
\end{equation*}
$$

The constants $r_{1}$ and $M_{1}$ are given by (2.8) and (2.11), respectively.
Proof. By (4.10), there exist $\theta \in(0,1)$ and $D \geqslant 1$ such that

$$
\begin{equation*}
\frac{\sup _{(j, \Omega \in \Omega} \gamma_{\mathrm{sG}}(\Lambda, \omega)}{C_{2} \psi(\operatorname{diam} A)} \leqslant \theta \leqslant \frac{C_{1} D}{C_{1} D+C_{0}} \tag{4.13}
\end{equation*}
$$

for any $\Lambda \in \mathscr{F}$ with $\operatorname{diam} A \geqslant D$. To prove (4.11), we may assume $\rho:=d\left(y, \Delta_{f}\right) \geqslant D+r_{0}$, since the left-hand side of (4.11) is not greater than $\|f\| \|$ in any case. Let $\left(\omega, \omega^{\prime}\right) \in \Omega^{2}$ be such that $\omega^{\prime} \equiv \omega$ off $y$. We set $g=d \mu^{A, \omega^{\prime}} / d \mu^{A,(")}$. Note that

$$
\|g\| \leqslant \exp (4\|\Phi\|) \quad \text { and } \quad \Delta_{g} \subset A \cap \partial_{r_{0}}\{y\}
$$

Using these and Lemma 4.1, we have

$$
\begin{align*}
\left|\left(\mu^{A . w^{\prime}}-\mu^{A . \omega}\right) f\right|= & \left|\mu^{A \cdot \omega}(g ; f)\right| \\
& \leqslant c_{1}\left\|\left|f \|| | \exp -\frac{\rho-r_{0}}{10 M_{1} r_{1} \gamma_{\mathrm{sG}}(\Lambda, \omega)}\right.\right.  \tag{4.14}\\
& +c_{2}\|\mid f\| \exp -\frac{\rho-r_{0}}{10 r_{1}} \tag{4.15}
\end{align*}
$$

Here and in what follows $c_{i}(i=1,2, \ldots)$ stand for constants which are independent of $\Lambda, f$, and $y$. Since $\operatorname{diam} A \geqslant \rho-r_{0} \geqslant D$, we can take advantage of (4.13) and (4.12) to bound (4.14) from above as follows:

$$
\begin{align*}
(4.14) & \leqslant c_{3}\|f\| \| \exp \left(-\frac{\rho}{10 M_{1} r_{1} \theta C_{2}\left(C_{1} \rho+C_{0}\right)} \log \left(C_{1} \rho+C_{0}\right)\right) \\
& \leqslant c_{3}\||f|\| \exp \left(-\frac{1}{10 C_{1} C_{2} M_{1} r_{1}} \log \left(C_{1} \rho+C_{0}\right)\right) \\
& \leqslant c_{3}\||f|\| \rho^{-1 /\left(5 C_{1} C_{2} M_{1} r_{1}\right)} \tag{4.16}
\end{align*}
$$

Similarly, we have that

$$
\begin{align*}
(4.15) & \leqslant c_{4} \mid\|f\| \| \exp -\frac{\rho}{10 r_{1}} \\
& \leqslant c_{5} \mid\|f\| \| \rho^{-1 /\left(10 c_{1} c_{2} M_{1} r_{1}\right)} \tag{4.17}
\end{align*}
$$

Putting (4.16) and (4.17) together, we conclude (4.11). QED
Lemma 4.3. Let $C_{1} \in(2, \infty)$ and $n>1$ be such that $n(4 n-3) / 2(n-1)^{2} \leqslant C_{1}$ and $\mathscr{F}=\mathscr{B}_{v}(m)$ for arbitrarily fixed $v \in \mathbf{Z}^{d}$ and $m \geqslant 1$. Suppose that there exist an admissible coefficient $\{c(f)>0\}_{f \in \mathscr{K}}$ and a nonincreasing function $\phi:(0, \infty) \rightarrow(0, \infty)$ such that

$$
\begin{equation*}
\operatorname{osc}_{y}\left(\mu^{\prime \cdot} \cdot f\right) \leqslant c(f) \phi\left(d\left(y, \Delta_{f}\right)\right) \tag{4.18}
\end{equation*}
$$

for all $A \in \mathscr{F}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{0}} A$ which satisfy $\operatorname{diam} A \leqslant C_{1} d\left(y, \Delta_{f}\right)$. Then there exists $C_{2} \in(0, \infty)$ such that

$$
\begin{equation*}
\operatorname{osc}_{y}\left(\mu^{A \cdot} \cdot f\right) \leqslant C_{2}\|f\| d\left(y, \Delta_{f}\right)^{d-1} \phi\left(d\left(y, \Delta_{f}\right)-2 m-r_{0}\right) \tag{4.19}
\end{equation*}
$$

for all $\Lambda \in \mathscr{F}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{0}} \Lambda$ which satisfy $d\left(y, \Delta_{f}\right) \geqslant n\left(2 m+r_{0}\right)$.
For $\mathscr{F}=\mathscr{A}$ the same statement with $m=1$ is true.
Proof. Take $A \in \mathscr{F}, f \in \mathscr{C}_{A}$, and $y \in \partial_{r_{0}} \Lambda$, which satisfy $\rho:=d\left(y, \Delta_{f}\right) \geqslant$ $n\left(2 m+r_{0}\right)$. We are going to apply our assumption to $A \cap \Gamma$ rather than $\Lambda$ itself, where $\Gamma$ is an element of $\mathscr{C}_{v}(m)$ we now define. We define $\Gamma_{k} \in \mathscr{C}_{v}(m)$ ( $k=1,2, \ldots$ ) by

$$
\Gamma_{k}=\left\{x: d\left(x, \Gamma_{0}\right) \leqslant k m\right\}
$$

where $\Gamma_{0}$ is an element in $\mathscr{C}_{v}(m)$ such that diam $\Gamma_{0}=m$ and $y \in \Gamma_{0}$. We set $\Gamma=\Gamma_{K}$, where $K=\max \left\{k \geqslant 1: \Gamma_{k} \cap \Delta_{f}=\varnothing\right\}$. We then have by definition that

$$
\begin{equation*}
K m \leqslant \rho \leqslant(K+2) m \tag{4.20}
\end{equation*}
$$

and hence that

$$
\begin{equation*}
K \geqslant \frac{\rho}{m}-2 \geqslant 2(n-1) \tag{4.21}
\end{equation*}
$$

For fixed $\left(\omega, \omega^{\prime}\right) \in \Omega^{2}$ with $\omega \equiv \omega^{\prime}$ off $y$ we set $g=d \mu^{A . m^{\prime}} / d \mu^{\prime \prime \prime}$. Note that

$$
\begin{equation*}
\|g\| \leqslant \exp (4\|\Phi\|), \quad \Delta_{g} \subset \Lambda \cap \partial_{r_{0}}\{y\} \tag{4.22}
\end{equation*}
$$

and that for $z \in \Lambda \cap \partial_{r_{0}} \Gamma$

$$
\begin{equation*}
\operatorname{diam}(\Gamma \cap \Lambda) \leqslant C_{1} d\left(z, \Delta_{g}\right) \tag{4.23}
\end{equation*}
$$

The first two observations are standard. The third one can be seen as follows. We have by (4.22) and (4.20) that

$$
\begin{align*}
d\left(z, \Delta_{g}\right) & \geqslant|z-y|-r_{0} \\
& \geqslant K m-r_{0} \\
& \geqslant\left(\frac{\rho}{m}-2\right) m-r_{0} \\
& =\rho-2 m-r_{0} \\
& \geqslant \frac{n-1}{n} \rho \tag{4.24}
\end{align*}
$$

By (4.20), (4.21), and (4.24), we get

$$
\begin{aligned}
\operatorname{diam}(\Gamma \cap A) & \leqslant(2 K+1) m \\
& \leqslant\left(2+\frac{1}{K}\right) \rho \\
& \leqslant\left(2+\frac{1}{2(n-1)}\right) \frac{n}{n-1} d\left(z, \Delta_{g}\right) \\
& \leqslant C_{1} d\left(z, \Delta_{g}\right)
\end{aligned}
$$

The proof of (4.19) comes down to the following estimate:

$$
\begin{equation*}
\left\|\mu^{\Gamma \cap A \cdot}(g-1)\right\| \leqslant C_{2} \rho^{d-1} \phi\left(\rho-2 m-r_{0}\right) \tag{4.25}
\end{equation*}
$$

In fact,

$$
\begin{aligned}
\left|\left(\mu^{A, m^{\prime}}-\mu^{A \cdot(\prime \prime}\right) f\right| & =\left|\mu^{A, w_{2}}[f(g-1)]\right| \\
& \leqslant\|f\| \cdot\left\|\mu^{\Gamma \cap A \cdot}(g-1)\right\|
\end{aligned}
$$

On the other hand, (4.25) can be seen as follows:

$$
\begin{aligned}
\left|\mu^{\Gamma \cap A \cdot \xi}(g-1)\right| & =\left|\mu^{A, c_{1}}\left(\mu^{\Gamma \cap A, \xi} g-\mu^{\Gamma \cap A \cdot} g\right)\right| \\
& \leqslant \sup _{(\xi, \eta) \in \Omega^{2}}\left\{\left|\mu^{\Gamma \cap A \cdot \xi} g-\mu^{\Gamma \cap A \cdot \eta} g\right| ; \xi \equiv \eta \text { outside } A\right\} \\
& \leqslant\left|A \cap \partial_{r_{0}} \Gamma\right| \sup _{=\in A \cap \partial_{r_{1}} \Gamma} \operatorname{osc}_{z}\left(\mu^{\Gamma \cap A} \cdot g\right) \\
& \leqslant c_{1} \rho^{d-1} c(g) \sup _{=\in A \cap \partial_{r_{1}} \Gamma} \phi\left(d\left(z, A_{g}\right)\right) \\
& \leqslant c_{2} \rho^{d-1} \phi\left(\rho-2 m-r_{0}\right)
\end{aligned}
$$

The first equality comes from an identity: $1=\mu^{A, " \prime \prime} g=\mu^{A \cdot \prime \prime \prime}\left(\mu^{I \cap A \cdot} g\right)$. The inequality in the fourth line is an application of (4.18) to $\Gamma \cap A \in \mathscr{F}$, which is allowed by (4.23) under our present assumption. To proceed to the last line, we used (4.22). QED

Lemma 4.4. Let $\mathscr{F}$ be either $\mathscr{A}$ or $\mathscr{B}_{r}(m)$ for arbitrarily fixed $v \in \mathbf{Z}^{d}$ and $m \geqslant 1$. Suppose that (2e) in Theorem 3.2 holds and set $\alpha_{x, r}=\phi(|x-y|)$. Then,

$$
\begin{equation*}
\lim _{L \rightarrow \infty} \frac{1}{L^{d}} \sum_{x \in[0, L)^{d}} \sum_{y \in c_{r_{1},}[0, L)^{d}} \alpha_{x, y}=0 \tag{4.26}
\end{equation*}
$$

and for all $X \in \mathscr{F}$ and $y \notin X$

$$
\begin{equation*}
\left\|\nabla_{.} \mu^{x \cdot} \cdot f-\mu^{x \cdot} \cdot \nabla_{y} f\right\| \leqslant \sum_{x \in X}\left\|\nabla_{x} f\right\| \alpha_{x, y} \tag{4.27}
\end{equation*}
$$

Remark 4.2. A mixing property defined by the set of conditions (4.26) and (4.27) is essentially the same as the condition (GS2) in ref. 7 and the condition $\operatorname{DSM}(\mathrm{Y})$ in ref. 16. The condition turns out to be equivalent to $\operatorname{DSM}(\mathscr{F})$ in the end (cf. the proof of Theorem 3.2).

Proof. It is easy to see (4.26). In fact,

$$
\begin{aligned}
\sum_{x \in\left[0 . L y^{\prime}\right.} \sum_{y \in \partial_{r_{0}}\left[0, L y^{d}\right.} \alpha_{x, y} & =\sum_{x \in[0, L)^{d}} \sum_{y \in \partial_{r_{1}[ }[0, L)^{d}} \phi(|x-y|) \\
& \leqslant \sum_{m=1}^{L+r_{0}} \sum_{y \in \partial_{r_{0}}\left[0, L y^{d}\right.} \sum_{x:|x-x|=m} \phi(m) \\
& \leqslant c_{1} L^{d-1} \sum_{m=1}^{L+r_{0}} m^{d-1} \phi(m)
\end{aligned}
$$

and thus

$$
\begin{aligned}
\frac{1}{L^{d}} \sum_{x \in[0, L)^{d}} \sum_{y \in \partial_{r_{0}}\left[0, L y^{\prime}\right.} \alpha_{x, y} & \leqslant \frac{c_{1}}{L} \sum_{m=1}^{L+r_{0}} m^{d-1} \phi(m) \\
& \rightarrow 0 \quad \text { as } L \rightarrow \infty
\end{aligned}
$$

We next prove (4.27). Take $X \in \mathscr{F}$ and an alignment $\left\{x_{m}\right\}_{m=1}^{n}$ of points in $X$. We set $I_{m}=\left\{x_{j}\right\}_{j=1}^{m-1}$ and $J_{m}=\left\{x_{j}\right\}_{j=m}^{n}$. Note first that

$$
\nabla_{X} f(\sigma)=\sum_{m=1}^{n} \int \lambda^{l_{m}}\left(d \xi_{t_{m}}\right) \nabla_{r_{m}} f\left(\xi_{l_{m}} \cdot \sigma_{l_{m}}\right)
$$

Using this expression, we have that

$$
\begin{align*}
& \nabla_{y} \mu^{X \cdot \omega} f-\mu^{X, \omega} \nabla_{y} f \\
&=\int \lambda(d s)\left(\mu^{\lambda \cdot(\cdot) \cdot s}-\mu^{X \cdot(\prime)}\right)\left(d \sigma_{X}\right) f\left(\sigma_{X} \cdot(\omega \cdot s)_{X^{x}}\right) \\
&=-\int \lambda(d s)\left(\mu^{X \cdot \omega \cdot s}-\mu^{X \cdot(\prime)}\right)\left(d \sigma_{X}\right) \nabla_{X} f\left(\sigma_{X} \cdot(\omega \cdot s)_{X^{x}}\right) \\
&=-\sum_{m=1}^{n} \int \lambda(d s) \int\left(\mu^{X \cdot(1) \cdot s}-\mu^{X \cdot(1)}\right)\left(d \sigma_{X}\right) \overline{f m}\left(\sigma_{X} \cdot(\omega \cdot s)_{X^{x}}\right) \tag{4.28}
\end{align*}
$$

where $\omega \cdot s$ stands for a configuration obtained from $\sigma$ by replacing $\sigma_{y}$, by $s$ and

$$
\overline{f_{m}}(\sigma)=\int \lambda^{I_{m}}\left(d \xi_{I_{m}}\right) \nabla_{x_{m}} f\left(\xi_{I_{m}} \cdot \sigma_{I_{m}}\right)
$$

At this point, fix $y$ and choose the alignment so that $\left|x_{f}-y\right|$ is nondecreasing in $j=1,2, \ldots$. Since $\left\|\overline{f_{m}}\right\| \leqslant\left\|\nabla_{x_{m}} f\right\|$ and " $\sigma_{x} \mapsto \overline{f_{m}}\left(\sigma_{X^{x}} \cdot(\omega \cdot s)_{X^{x}}\right) " \in \mathscr{C}_{J_{m}}$, we have by (4.28) and (2e) that

$$
\begin{aligned}
\left\|\nabla_{y} \mu^{x \cdot f} f-\mu^{x \cdot} \cdot \nabla_{y} f\right\| & \leqslant \sum_{m=1}^{n}\left\|\nabla_{x_{m}} f\right\| \phi\left(d\left(y, J_{m}\right)\right) \\
& =\sum_{m=1}^{n}\left\|\nabla_{x_{m}} f\right\| \phi\left(\left|x_{m}-y\right|\right) \\
& =\sum_{x \in X}\left\|\nabla_{x} f\right\| \alpha_{x \cdot y}
\end{aligned}
$$

This proves (4.27). QED
Proof of Theorem 3.2. In view of Theorem 3.1, all conditions in Theorem 3.2 are obviously necessary for the Dobrushin-Shlosman mixing condition (2.16) to be true. To show that (2a), (2b), (2d), and (2e) are sufficient, we prove the following sequence of implications:

$$
\begin{equation*}
(2 \mathrm{a}) \Rightarrow(2 \mathrm{~b}) \Rightarrow(2 \mathrm{~d}) \Rightarrow(2 \mathrm{e}) \Rightarrow(1 \mathrm{~b}) \tag{4.2}
\end{equation*}
$$

The condition (1b) is equivalent to the Dobrushin-Shlosman mixing condition by Theorem 3.1. The condition (2a) immediately implies ( 2 b ), since $\gamma_{\mathrm{SG}} \leqslant \gamma_{\mathrm{LS}}[$ ref. 2, p. 224, (6.1.7)]. Implication from (2b) to (2d) is a consequence of Lemma 4.2, and (2d) implies (2e) by Lemma 4.3. Now suppose that (2e) holds. This means that we may assume the conclusion of Lemma 4.4. Then, by the argument used in Section 2 of ref. 16, we obtain (1b) in Theorem 3.1. See the derivation of Corollary 2.8 in that paper. This completes the proof of (4.29). On the other hand, Lemma 4.1 says that ( 2 c ) implies (2d), and hence also (1b) by (4.29). QED

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