九州大学学術情報リポジトリ
Kyushu University Institutional Repository

Construction of Gibbs measures for 1－ dimensional continuum fields<br>Hariya，Yuu<br>Research Institute for Mathematical Sciences，Kyoto University

https：／／hdl．hand le．net／2324／3366

出版情報：MHF Preprint Series．MHF2005－19，2005－05－17．九州大学大学院数理学研究院 バージョン：権利関係：

# MHF Preprint Series 

Kyushu University
21st Century COE Program
Development of Dynamic Mathematics with High Functionality

# Construction of Gibbs measures for 1-dimensional continuum fields 

Y. Hariya

MHF 2005-19
( Received May 17, 2005 )

Faculty of Mathematics
Kyushu University
Fukuoka, JAPAN

# Construction of Gibbs measures for 1-dimensional continuum fields 

Yuu Hariya*


#### Abstract

We study 1-dimensional continuum fields of Ginzburg-Landau type under the presence of an external and a long-range pair interaction potentials. The corresponding Gibbs states are formulated as Gibbs measures relative to Brownian motion [17]. In this context we prove the existence of Gibbs measures for a wide class of potentials including a singular external potential as hard-wall ones, as well as a non-convex interaction. Our basic methods are: (i) to derive moment estimates via integration by parts; and (ii) in its finite-volume construction, to represent the hard-wall Gibbs measure on $C\left(\mathbb{R} ; \mathbb{R}^{+}\right)$in terms of a certain rotationally invariant Gibbs measure on $C\left(\mathbb{R} ; \mathbb{R}^{3}\right)$.


## 1 Introduction

This paper studies $\mathbb{R}^{d}$-valued continuum fields $X=\{X(s), s \in \mathbb{R}\}$ over $\mathbb{R}$ with interactions prescribed by a (formal) Hamiltonian given by

$$
\begin{equation*}
H(X)=\frac{1}{2} \int|\nabla X(s)|^{2} d s+\int \varphi(X(s)) d s+\frac{1}{2} \iint \psi(s-t, X(s), X(t)) d s d t \tag{1.1}
\end{equation*}
$$

where $\nabla X=\partial X / \partial s$. This is a continuous counterpart to Ginzburg-Landau random fields over the lattice $\mathbb{Z}$ (see [23]) with an external potential $\varphi$ and a long-range pair interaction potential $\psi$; analogously to those lattice cases, the Gibbs state associated with (1.1) would formally be given by

$$
\begin{equation*}
d \mu(X)=\exp \{-H(X)\} d X / \text { normalization } \tag{1.2}
\end{equation*}
$$

where the reference measure $d X \equiv \prod_{s \in \mathbb{R}} d X(s)$ is "Feynman's measure". By incorporating the first term in (1.1) into the reference measure, the measure like (1.2) can be

[^0]interpreted as a local perturbation from Wiener measure; that is, we consider Brownian motion under the presence of potentials $\varphi$ and $\psi$. The corresponding measure on the path space $C\left(\mathbb{R} ; \mathbb{R}^{d}\right)$ is formulated as a Gibbs measure relative to Brownian motion [17], and defined through the so-called DLR equation (see Definition 2.1). Similarly to the lattice field case, such a measure appears as equilibrium states of random timeevolutions of the fields described by Ginzburg-Landau equations (see [8, 9, 7]); [11] also studies the dynamics by using the Dirichlet form theory.

In the case $\psi=0$, the corresponding Gibbs measure can be realized as a $P(\phi)_{1^{-}}$ stationary Markov process and it has been fairly understood (see, e.g., [20, 22]). On the other hand, in the non-Markovian case $\psi \neq 0$, the first mathematical treatment was done by [17]; their motivation was originated from quantum field theory, particularly from the Euclidean quantization of a certain Hamiltonian operator called Nelson's scalar field model [16]: its physical context is a quantum particle in $\mathbb{R}^{d}$ governed by the Schrödinger operator $H_{0}=-(1 / 2) \Delta+\varphi$ and coupled to a free Bose field. It has been known [16] that, when the coupling is restricted to the time interval $[-T, T]$, the Euclidean quantization involves the following types of measures:

$$
\begin{equation*}
Z_{T}^{-1} \exp \left\{\int_{|s|,|t| \leq T} d s d t \int_{\mathbb{R}^{d}} d k \frac{\tilde{f}(k)}{f(k)} e^{-f(k)|s-t|} \cos (k \cdot(X(s)-X(t)))\right\} d \mathcal{P}^{\varphi}(X) \tag{1.3}
\end{equation*}
$$

where $\mathcal{P}^{\varphi}$ denotes the law on $C\left(\mathbb{R} ; \mathbb{R}^{d}\right)$ of the $P(\phi)_{1}$-process associated with $H_{0}, f$ and $\tilde{f}$ are some non-negative functions, and $Z_{T}$ is the normalization. It is discussed in [15] that the limit measure as $T \rightarrow \infty$ (if it exists) plays an important role in the analysis of spectral structure of the operator. For details, see the references cited above.

Now we raise the question whether the family of measures as (1.3) has a limit as $T \rightarrow$ $\infty$; that is, we discuss the existence of an infinite-volume Gibbs measure for prescribed potentials $\varphi$ and $\psi$. There are several possible ways to show the existence: Ruelle's superstability estimates [21, 13] on lattice fields; monotonicity method developed in [17]; cluster expansions applied in [14]. Superstability estimate relies on the product structure of a reference measure and there are a number of difficulties in adapting this to the continuum case. Monotonicity uses log-concave inequalities, hence requires certain convexity on $\psi$, which excludes interactions as in (1.3). Cluster expansion is applicable to such an interaction when the coupling is sufficiently weak; in [14], the growth order in spatial variables was also restricted to at most quadratic. Recently, [1] has proved the existence under a mild assumption on $\varphi$; his condition is almost as weak as assuming $H_{0}$ to have an $L^{2}$-ground state, and the existence is shown in the case of bounded interactions satisfying a certain pathwise condition.

One of the main purposes of this paper is to establish the existence for a wide class of interactions. We deal with interactions of growth order $q_{0}+1\left(q_{0} \geq 0\right)$ in spatial
variables. We choose external potentials to grow polynomially of order $2 p_{0}\left(p_{0} \geq 1\right)$ at infinity. Roughly speaking, our result reads:

$$
\text { if } p_{0} \geq q_{0} \text {, then the associated Gibbs measure exists. }
$$

The key is to show the localization in its finite-volume construction (see Theorem 3.1). Our method to derive this is different from those quoted above. Its advantages are: (i) it requires no convexity on interactions; (ii) it is applicable to an arbitrary strength of interactions in the case $p_{0}>q_{0}$. Our technique to be developed here involves (deterministic) integration by parts. So we require weak differentiability on potentials, however, many examples including Nelson's model can be treated in our framework.

Another purpose of this paper is to construct hard-wall Gibbs measures with interactions on $C\left(\mathbb{R} ; \mathbb{R}^{+}\right)$. We deal with an external potential such as

$$
\varphi(x)=\infty, \quad x<0
$$

Under the effect of $\varphi$, the path is restricted to stay positive over the real line (the hard wall). Similarly to lattice fields [12, 3], a "repelling" phenomenon is also observed in the present continuum case; indeed, the paths, renormalized by the hard-wall effect, behave as 3-dimensional Bessel process, not as reflecting Brownian motion. This observation reduces the problem to the localization of certain finite-volume Gibbs measures on $C\left(\mathbb{R} ; \mathbb{R}^{3}\right)$, to which our method applies.

This paper is organized as follows: in Section 2, we state the main result; in Section 3, we prepare a key theorem, which we prove in Section 4; in Section 5, we discuss the existence of hard-wall Gibbs measures with interactions.

Throughout this paper we use the notation $(x, y)$ instead of $x \cdot y$ for the inner product on $\mathbb{R}^{d}$. We write $|x|=(x, x)^{1 / 2}$ and $\|x\|=\left(1+|x|^{2}\right)^{1 / 2}$ for $x \in \mathbb{R}^{d}$. $\mathbb{R}^{+}$denotes the set of the non-negative real numbers. We denote by $\langle\cdot\rangle$ the expectation with subscript of a reference measure. Other notation will be introduced as needed.

## 2 Main result

For given functions $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R} \cup\{\infty\}$ and $\psi: \mathbb{R} \times \mathbb{R}^{d} \times \mathbb{R}^{d} \rightarrow \mathbb{R}$, we define $(\varphi, \psi)$-Gibbs measures following [17]: Let $\mathcal{C}=C\left(\mathbb{R} ; \mathbb{R}^{d}\right)$ endowed with the compact uniform topology. For a finite interval $\Lambda=\left[T_{1}, T_{2}\right] \subset \mathbb{R}$, let $\mathcal{C}_{\Lambda}=C\left(\Lambda ; \mathbb{R}^{d}\right)$ and $\mathcal{C}_{\Lambda}^{*}=C\left(\Lambda^{\mathrm{c}} ; \mathbb{R}^{d}\right)$. For a given $\xi \in \mathcal{C}$, we define the local Hamiltonian $\mathcal{H}_{\Lambda, \xi} \equiv \mathcal{H}_{\Lambda, \xi}^{\varphi, \psi}$ by

$$
\begin{aligned}
\mathcal{H}_{\Lambda, \xi}(X)=\int_{\Lambda} \varphi(X(s)) d s+\frac{1}{2} \int_{\Lambda^{2}} & \psi(s-t, X(s), X(t)) d s d t \\
& +\int_{\Lambda \times \Lambda^{c}} \psi(s-t, X(s), \xi(t)) d s d t, \quad X \in \mathcal{C}_{\Lambda}
\end{aligned}
$$

The functions $\varphi$ and $\psi$ are called the external potential and interaction potential, respectively. We assume $\psi$ is symmetric in the sense that, for all $s \in \mathbb{R}$ and $x, y \in \mathbb{R}^{d}$,

$$
\begin{equation*}
\psi(s, x, y)=\psi(|s|, x, y), \quad \psi(s, x, y)=\psi(s, y, x) \tag{2.1}
\end{equation*}
$$

Let $\pi_{\Lambda}: \mathcal{C} \rightarrow \mathcal{C}_{\Lambda}$ and $\pi_{\Lambda}^{*}: \mathcal{C} \rightarrow \mathcal{C}_{\Lambda}^{*}$ be canonical projections. For a probability measure $\mu$ on $\mathcal{C}$, define $\mu_{\Lambda, \xi}(\cdot)=\mu\left(\pi_{\Lambda} \in \cdot \mid \pi_{\Lambda}^{*}\right)(\xi)$. Here $\mu\left(\cdot \mid \pi_{\Lambda}^{*}\right)$ is the regular conditional probability with respect to the $\sigma$-field $\sigma\left(\pi_{\Lambda}^{*}\right)$.
Definition 2.1. A probability measure $\mu$ on $\mathcal{C}$ is called a $(\varphi, \psi)$-Gibbs measure if its regular conditional probabilities satisfy the DLR equation:

$$
d \mu_{\Lambda, \xi}(X)=\left(Z_{\Lambda, \xi}\right)^{-1} e^{-\mathcal{H}_{\Lambda, \xi}(X)} d \mathcal{W}_{\Lambda, \xi}(X) \quad \mu \text {-a.e. } \xi \in \mathcal{C} .
$$

Here $\mathcal{W}_{\Lambda, \xi}$ denotes the law of Brownian bridge on $\mathcal{C}_{\Lambda}$ with boundary conditions $X\left(T_{1}\right)=$ $\xi\left(T_{1}\right)$ and $X\left(T_{2}\right)=\xi\left(T_{2}\right)$, and $Z_{\Lambda, \xi}=\left\langle e^{-\mathcal{H}_{\Lambda, \xi}(X)}\right\rangle_{\mathcal{W}_{\Lambda, \xi}}$.

We proceed to the setup for the main result. We assume (A.1)-(A.3):
(A.1) Assumptions on $\varphi$. $\varphi$ is bounded from below. There exist $\varphi_{0}$ and $\varphi_{1}$ such that $\varphi=\varphi_{0}+\varphi_{1}$, satisfying (A.1a) and (A.1b), respectively:
(A.1a) $\varphi_{0}$ is a continuous function such that the associated Schrödinger operator $H_{0}=$ $-(1 / 2) \Delta+\varphi_{0}$ acting on $L^{2}\left(\mathbb{R}^{d} ; d x\right)$ has a strictly positive ground state $f_{0}$ of class $C^{2}\left(\mathbb{R}^{d}\right)$ satisfying the following conditions:
(i) (strict log-concavity) there exists an $\alpha>0$ such that

$$
\begin{equation*}
\left(\zeta, \operatorname{Hess}_{u_{0}}(x) \zeta\right) \geq \alpha|\zeta|^{2} \quad \text { for all } \zeta, x \in \mathbb{R}^{d} \tag{2.2}
\end{equation*}
$$

where $u_{0}=-\log f_{0}$ and $\operatorname{Hess}_{u_{0}}=\left(\partial^{2} u_{0} / \partial x_{i} \partial x_{j}\right)_{1 \leq i, j \leq d}$ is the Hessian of $u_{0}$;
(ii) there exists a $p_{0} \geq 1$ such that

$$
\begin{equation*}
0<\liminf _{r \rightarrow \infty} \frac{1}{r^{2 p_{0}}} \inf _{|x|=r} U(x), \quad \quad \limsup _{r \rightarrow \infty} \frac{1}{r^{p_{0}}} \sup _{|x|=r}|\boldsymbol{V}(x)|<\infty, \tag{2.3}
\end{equation*}
$$

where $U=\left(f_{0}\right)^{-2} \operatorname{div}\left(f_{0} \nabla f_{0}\right)$ and $\boldsymbol{V}=\left(f_{0}\right)^{-1} \nabla f_{0}$. Here div denotes the divergence.
(A.1b) $\varphi_{1} \in W_{\text {loc }}^{1,1}\left(\mathbb{R}^{d}\right)$ and there exist $\mathrm{b} \geq 0$ and $0 \leq p_{1}<p_{0}$ such that

$$
\left|\nabla \varphi_{1}(x)\right| \leq \mathrm{b}\|x\|^{p_{1}} \quad \text { for a.e. } x \in \mathbb{R}^{d} .
$$

Here $W_{\mathrm{loc}}^{1,1}\left(\mathbb{R}^{d}\right)$ is the set of functions $f \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{d}\right)$ whose distributional derivatives $\partial f / \partial x_{i}, 1 \leq i \leq d$, belong to $L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{d}\right)$.
(A.2) Assumptions on $\psi$. For each fixed $s \in \mathbb{R}$ and $y \in \mathbb{R}^{d}, \psi(s, \cdot, y) \in W_{\mathrm{loc}}^{1,1}\left(\mathbb{R}^{d}\right)$. There exists a non-negative, integrable function $\psi_{0}$ on $\mathbb{R}$ satisfying (i) and (ii):
(i) there exists a $q_{0} \geq 0$ such that, for a.e. $s \in \mathbb{R}$ and $x, y \in \mathbb{R}^{d}$,

$$
\left|\nabla_{x} \psi(s, x, y)\right| \leq \psi_{0}(s)\left(\|x\|^{q_{0}}+\|y\|^{q_{0}}\right) \quad\left(\nabla_{x}=\left(\partial / \partial x_{1}, \ldots, \partial / \partial x_{d}\right)\right)
$$

(ii) for a.e. $s \in \mathbb{R}$ and $x, y \in \mathbb{R}^{d}, \psi(s, x, y) \geq-\psi_{0}(s)$.
(A.3) $p_{0}$ is strictly larger than $q_{0}$.

Remark 2.1. From (2.2), we see in particular that $|\boldsymbol{V}(x)|-\alpha|x|$ is bounded from below. We assumed $p_{0} \geq 1$ in (2.3) to make our assumption consistent.

Now we state one of our main results, the existence of $(\varphi, \psi)$-Gibbs measures:
Theorem 2.1. Assume (A.1)-(A.3). Then there exists a translation invariant $(\varphi, \psi)$ Gibbs measure $\mu$ satisfying

$$
\left.\left.\langle | X(0)\right|^{2 p_{0}}\right\rangle_{\mu}<\infty
$$

Remark 2.2. In the case $p_{0}=q_{0}$, the existence also holds if $\int_{\mathbb{R}} \psi_{0}(s) d s$ is sufficiently small. See Remark 4.2.

We give examples of potentials satisfying the above assumptions.
Example 2.1. We can construct examples of $\varphi$ from $f_{0}$. Note that we may assume without loss of generality that $-(1 / 2) \Delta f_{0}+\varphi_{0} f_{0}=0$, by adding a constant to $\varphi_{0}$ if necessary. $\varphi_{0}$ is then expressed as $\varphi_{0}=(1 / 2) \Delta f_{0} / f_{0}$.
(i) Ornstein-Uhlenbeck processes:

$$
\varphi_{\mathrm{ou}}(x)=\frac{1}{2} \sum_{i=1}^{d} a_{i}^{2} x_{i}^{2}, \quad x \in \mathbb{R}^{d}, a_{i}>0
$$

with $f_{0}(x)=\exp \left\{-(1 / 2) \sum_{i=1}^{d} a_{i} x_{i}^{2}\right\}$.
(ii) Double-well potentials:

$$
\varphi_{\mathrm{dw}}(x)=\frac{1}{2} a^{2} x^{4}-\frac{1}{2} b^{2} x^{2}, \quad x \in \mathbb{R}, a, b>0 .
$$

This is obtained by taking

$$
\varphi_{0}(x)=\frac{1}{2} a^{2}\left(x^{4}+x^{2}\right)+R(x), \quad \varphi_{1}(x)=-\frac{1}{2}\left(a^{2}+b^{2}\right) x^{2}-R(x)
$$

with $R(x)=-a\|x\|+(a / 2)\|x\|^{-1}$. In this case we take $f_{0}(x)=\exp \left\{-(a / 3)\|x\|^{3}\right\}$.
Examples of $\psi$ are:
Example 2.2. (i) Nelson's scalar field model: Typically, it is given by

$$
\psi_{\text {nel }}(s, x, y)=-\frac{1}{s^{2}+1+|x-y|^{2}}
$$

We refer to [1]. This example corresponds to the case $q_{0}=0$.
(ii) Non-convex interactions:

$$
\psi_{\mathrm{nc}}(s, x, y)=\psi_{0}(s) v(x-y),
$$

where $v(x)=|x|^{q_{0}+1}+Q(|x|)$ with $Q$ a polynomial whose degree is less than $q_{0}+1$, and $\psi_{0}$ is as in (A.2). Note that $v$ need not be convex.

By Theorem 2.1 we have established the existence for a large class of interactions including Example 2.2; indeed, applying the theorem to the above examples, we conclude: (i) there always exists a Gibbs measure for $\left(\varphi_{\mathrm{ou}}, \psi_{\text {nel }}\right)$ and for $\left(\varphi_{\mathrm{dw}}, \psi_{\text {nel }}\right)$; (ii) there exists a Gibbs measure for $\left(\varphi_{\mathrm{ou}}, \psi_{\mathrm{nc}}\right)$ (resp. for $\left.\left(\varphi_{\mathrm{dw}}, \psi_{\mathrm{nc}}\right)\right)$ if $q_{0}<1$ (resp. $q_{0}<2$ ).

## 3 Key theorem: Localization

To show Theorem 2.1, we prepare a key theorem, namely, the localization of finitevolume Gibbs measures. As was already suggested in (1.3), we will incorporate the effect of the external potential into the reference measure; this means we take a $P(\phi)_{1^{-}}$ process as a reference process, instead of a Brownian motion itself.

For the Schrödinger operator $H_{0}=-(1 / 2) \Delta+\varphi_{0}$, recall that the associated $P(\phi)_{1^{-}}$ process is determined by the stochastic differential equation (SDE)

$$
\begin{equation*}
d X(t)=d W(t)-\nabla u_{0}(X(t)) d t \tag{3.1}
\end{equation*}
$$

with $W$ a $d$-dimensional Brownian motion. Here $u_{0}=-\log f_{0}$ as in (i) of (A.1a). The process $X$ is stationary under the measure $f_{0}(x)^{2} d x$, which we denote by $m_{0}(d x)$. Here we assume that $f_{0}$ is normalized so that $m_{0}(d x)$ is a probability measure. We denote by $\mathcal{P}^{\varphi_{0}}$ the law of this process on $\mathcal{C}$. For the precise definition of $P(\phi)_{1}$-processes and related notion, see $[20,22]$ and references therein; see also [2] for more detailed descriptions.

We take $\mathcal{P}^{\varphi_{0}}$ as a reference measure, and perturb it by $\varphi_{1}$ and $\psi$ : For a finite interval $\Lambda \subset \mathbb{R}$, let $\mathcal{P}_{\Lambda}^{\varphi_{0}}$ denote the restriction of $\mathcal{P}^{\varphi_{0}}$ to the $\sigma$-field $\sigma\left(\pi_{\Lambda}\right)$. Define the finite-volume Gibbs measure $\mu_{\Lambda}^{\varphi, \psi}$ by

$$
\begin{equation*}
d \mu_{\Lambda}^{\varphi, \psi}(X)=\left(Z_{\Lambda}^{\varphi, \psi}\right)^{-1} \exp \left\{-\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}(X)\right\} d \mathcal{P}_{\Lambda}^{\varphi_{0}}(X) \tag{3.2}
\end{equation*}
$$

where

$$
\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}(X)=\int_{\Lambda} \varphi_{1}(X(s)) d s+\frac{1}{2} \int_{\Lambda^{2}} \psi(s-t, X(s), X(t)) d s d t
$$

and $Z_{\Lambda}^{\varphi, \psi}=\left\langle\exp \left\{-\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}(X)\right\}\right\rangle_{\mathcal{P}^{\varphi_{0}}}$. Note that $Z_{\Lambda}^{\varphi, \psi}<\infty$ by assumption. $\mu_{\Lambda}^{\varphi, \psi}$ is a probability measure on $\mathcal{C}_{\Lambda}$ with free boundary condition. If we take a symmetric interval $\Lambda=[-T, T]$, each corresponding subscript $\Lambda$ is replaced by $T$; e.g., $\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}=\mathcal{H}_{T}^{\varphi_{1}, \psi}$, $\mu_{\Lambda}^{\varphi, \psi}=\mu_{T}^{\varphi, \psi}$ and so on.

Note that by (i) of (A.2), the growth order of $\psi$ in spatial variables $x$ and $y$ is less than or equal to $q_{0}+1$. Suppose we had proved that $\left\{\mu_{\Lambda}^{\varphi, \psi}\right\}_{\Lambda \subset \mathbb{R}}$ has the following localization property:

$$
\begin{equation*}
\left.\left.\sup _{\Lambda \subset \mathbb{R}} \max _{t \in \Lambda}\langle | X(t)\right|^{q_{1}}\right\rangle_{\mu_{\Lambda}^{\varphi, \psi}}<\infty \tag{*}
\end{equation*}
$$

for some $q_{1} \geq q_{0}+1$. Then, from the same argument as in [17, Sect. 4], the existence of a $(\varphi, \psi)$-Gibbs measure follows. The next theorem shows that $(*)$ is really the case:

Theorem 3.1. Under (A.1a), (A.1b), (A.2) and (A.3), the localization (*) holds.
Remark 3.1. As we will see in the proof, the localization $(*)$ holds with $q_{1}=2 p_{0}$. Note that $2 p_{0} \geq q_{0}+1$ since we assume $p_{0} \geq 1$ and $p_{0}>q_{0}$.

## 4 Proof of Theorem 3.1

As was mentioned, once Theorem 3.1 is shown, then the existence result Theorem 2.1 can be obtained as a corollary, in the same way as [17]. We thus concentrate on proving Theorem 3.1. The essence of our proof is simple; we only have to use integration by parts once, which is formulated as:

$$
\int_{\mathbb{R}^{d}} \operatorname{div}\left(f_{0} \nabla_{z} f_{0}\right) e^{-\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}\left(X^{z}\right)} d z=-\int_{\mathbb{R}^{d}}\left(f_{0} \nabla_{z} f_{0}, \nabla_{z} e^{-\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}\left(X^{z}\right)}\right) d z
$$

where $X^{z}$ denotes a path given $X(0)=z$. See (4.8) below.
We begin with the following lemma:
Lemma 4.1. Assume (ii) of (A.1a). Then there exist constants $\mathrm{a}_{i}>0, i=1,2,3$, such that, for all $x \in \mathbb{R}^{d}$,

$$
\begin{equation*}
U(x) \geq \mathrm{a}_{1}\|x\|^{2 p_{0}}-\mathrm{a}_{2}, \quad|\boldsymbol{V}(x)| \leq \mathrm{a}_{3}\|x\|^{p_{0}} \tag{4.1}
\end{equation*}
$$

Proof. By (2.3), there exist $c, c^{\prime}>0$ and $r>0$ such that $U(x) \geq c|x|^{2 p_{0}}$ and $|\boldsymbol{V}(x)| \leq$ $c^{\prime}|x|^{p_{0}}$ for all $|x| \geq r$. By the continuity, $U$ is bounded from below and $|\boldsymbol{V}|$ from above on $\{|x| \leq r\}$. Combining these ends the proof.

Remark 4.1. By definition, $\int_{\mathbb{R}^{d}} U(x) m_{0}(d x)=0$, which implies that $U$ has a negative part. So the constant $a_{2}$ above cannot be taken negative; indeed, it must satisfy $a_{2} \geq$ $\mathrm{a}_{1} \int_{\mathbb{R}^{d}}\|x\|^{2 p_{0}} m_{0}(d x)$.

From now on we take a symmetric interval $\Lambda=[-T, T]$ for simplicity. Let

$$
M_{T}=\max _{|t| \leq T}\left\langle\|X(t)\|^{2 p_{0}}\right\rangle_{\mu_{T}^{\varphi, \psi}} .
$$

Proposition 4.2. Assume (A.1a), (A.1b), (A.2) and (A.3). Let $\mathrm{a}_{i}, i=1,2,3$, be as in Lemma 4.1. Then it holds that

$$
\begin{equation*}
\mathrm{a}_{1} M_{T}-\mathrm{a}_{2} \leq C_{d, \alpha}\left(\mathrm{~b} M_{T}^{\frac{p_{0}+p_{1}}{2 p_{0}}}+\mathrm{a}_{3} \bar{\psi}_{0} M_{T}^{\frac{p_{0}+q_{0}}{2 p_{0}}}\right) \tag{4.2}
\end{equation*}
$$

for all $T>0$. Here $C_{d, \alpha}=4 \sqrt{d} / \alpha$ and $\bar{\psi}_{0}=\int_{\mathbb{R}} \psi_{0}(s) d s$.

Once this proposition is shown, then the proof of Theorem 3.1 is straightforward:
Proof of Theorem 3.1. Let us consider the following equation in $m>0$ :

$$
\begin{equation*}
\mathrm{a}_{1} \mathrm{~m}-\mathrm{a}_{2}=C_{d, \alpha}\left(\mathrm{bm} \frac{\frac{p_{0}+p_{1}}{2 p_{0}}}{}+\mathrm{a}_{3} \bar{\psi}_{0} \mathrm{~m}^{\frac{p_{0}+q_{0}}{2 p_{0}}}\right) . \tag{4.3}
\end{equation*}
$$

If $p_{0}>q_{0}$, this equation has a unique solution $\mathrm{m}_{1}$. By Proposition 4.2, we then obtain

$$
\begin{equation*}
\max _{|t| \leq T}\left\langle\|X(t)\|^{2 p_{0}}\right\rangle_{\mu_{T}^{\varphi, \psi}} \leq \mathrm{m}_{1} \quad \text { for all } T>0 \tag{4.4}
\end{equation*}
$$

This implies the theorem.
For simplicity, we prove Proposition 4.2 in the case $\varphi_{1} \equiv 0$; that is, we let $\mathbf{b}=0$ in (4.2). We write $\mathcal{H}_{\Lambda}^{\psi}(X)$ for $\mathcal{H}_{\Lambda}^{\varphi_{1}, \psi}(X)$ accordingly. In the following we suppress potentials from the notation if there is no confusion; e.g., $Z_{\Lambda}^{\varphi, \psi}=Z_{\Lambda}, \mu_{T}^{\varphi, \psi}=\mu_{T}$ and so on.

For each fixed $t_{0} \in[-T, T]$, let $I=\left\langle U\left(X\left(t_{0}\right)\right)\right\rangle_{\mu_{T}}$. By the estimate on $U$ in (4.1), we have

$$
\begin{equation*}
\mathrm{a}_{1}\left\langle\left\|X\left(t_{0}\right)\right\|^{2 p_{0}}\right\rangle_{\mu_{T}}-\mathrm{a}_{2} \leq I \tag{4.5}
\end{equation*}
$$

We also have an upper estimate on $I$ as follows:
Proposition 4.3. It holds that

$$
I \leq \mathrm{a}_{3} \bar{\psi}_{0} C_{d, \alpha}\left\langle\left\|X\left(t_{0}\right)\right\|^{2 p_{0}}\right\rangle_{\mu_{T}}^{\frac{1}{2}} M_{T}^{\frac{q_{0}}{2 p_{0}}}
$$

Proposition 4.2 is an immediate consequence of these estimates:
Proof of Proposition 4.2. By (4.5) and Proposition 4.3, we have

$$
\mathrm{a}_{1}\left\langle\left\|X\left(t_{0}\right)\right\|^{2 p_{0}}\right\rangle_{\mu_{T}}-\mathrm{a}_{2} \leq \mathrm{a}_{3} \bar{\psi}_{0} C_{d, \alpha}\left\langle\left\|X\left(t_{0}\right)\right\|^{2 p_{0}}\right\rangle_{\mu_{T}}^{\frac{1}{2}} M_{T}^{\frac{q_{0}}{2 p_{0}}}
$$

Taking the maximum over $\left|t_{0}\right| \leq T$ on both sides leads to (4.2).
It now remains to prove Proposition 4.3, which we will do in a sequence of lemmas.
First we consider the disintegration of $\mathcal{P}^{\varphi_{0}}$ by conditioning on $X(0)$. For this purpose, let $W^{+}$and $W^{-}$, together with a probability measure $P_{W}$, be independent $d$-dimensional Brownian motions starting at 0 . Let $X^{z, \pm}(t) \equiv X^{z}\left(t, W^{ \pm}\right)$be the strong solutions of (3.1) starting at $z$. We set

$$
X^{z}(t)= \begin{cases}X^{z,+}(t), & t \geq 0 \\ X^{z,-}(-t), & t \leq 0\end{cases}
$$

Let $\mathcal{P}^{\varphi_{0}}(\cdot \mid X(0)=z)$ denote the regular conditional probability of $\mathcal{P}^{\varphi_{0}}$ given $X(0)=z$. From the Markov property of $\mathcal{P}^{\varphi_{0}}$, we easily deduce:

Lemma 4.4. The process $\left(\left\{X^{z}(t), t \in \mathbb{R}\right\}, P_{W}\right)$ has the same law as $\mathcal{P}^{\varphi_{0}}(\cdot \mid X(0)=z)$.
Recall from Section 3 that $\mathcal{P}^{\varphi_{0}}(X(0) \in d z)=m_{0}(d z)$. For a $\mathcal{P}^{\varphi_{0}}$-integrable functional $F$ on $\mathcal{C}$, we have, by Lemma 4.4 and by Fubini's theorem,

$$
\begin{align*}
\langle F(X)\rangle_{\mathcal{P}^{\varphi_{0}}} & =\int_{\mathbb{R}^{d}}\langle F(X) \mid X(0)=z\rangle_{\mathcal{P}^{\varphi_{0}}} m_{0}(d z) \\
& =\left\langle\int_{\mathbb{R}^{d}} F\left(X^{z}\right) m_{0}(d z)\right\rangle_{P_{W}} \tag{4.6}
\end{align*}
$$

By definition, $I=Z_{T}^{-1}\left\langle U\left(X\left(t_{0}\right)\right) e^{-\mathcal{H}_{T}^{\psi}(X)}\right\rangle_{\mathcal{P} \varphi_{0}}$. For convenience, we may shift $t_{0}$ to the origin by the stationarity of $\mathcal{P}^{\varphi_{0}}$; moreover, disintegrating $\mathcal{P}^{\varphi_{0}}$ as (4.6), we have

$$
\begin{equation*}
I=Z_{\Lambda}^{-1}\left\langle\int_{\mathbb{R}^{d}} U(z) e^{-\mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)} m_{0}(d z)\right\rangle_{P_{W}} \tag{4.7}
\end{equation*}
$$

where $\Lambda=\left[T_{1}, T_{2}\right]$ with $T_{1}=-\left(T+t_{0}\right)$ and $T_{2}=T-t_{0}$. Recall $m_{0}(d z)=f_{0}(z)^{2} d z$. By integration by parts formula, and by the estimate on $|\boldsymbol{V}|$ in (4.1),

$$
\begin{align*}
\int_{\mathbb{R}^{d}} U(z) e^{-\mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)} m_{0}(d z) & =\int_{\mathbb{R}^{d}}\left(\boldsymbol{V}(z), \nabla_{z} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)\right) e^{-\mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)} m_{0}(d z) \\
& \leq \mathrm{a}_{3} \int_{\mathbb{R}^{d}}\|z\|^{p_{0}}\left|\nabla_{z} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)\right| e^{-\mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)} m_{0}(d z) \tag{4.8}
\end{align*}
$$

We shall estimate $\left|\nabla_{z} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)\right|$ from above by using the following lemma:
Lemma 4.5. For each $z \in \mathbb{R}^{d}$ and $1 \leq i \leq d$, let $Y_{i}^{z}(t)=\left(\partial X_{j}^{z}(t) / \partial z_{i}\right)_{1 \leq j \leq d}$. Under the condition (i) of (A.1a), it holds that, for all $z$ and $i$,

$$
\left|Y_{i}^{z}(t)\right| \leq e^{-\alpha|t|}, \quad t \in \mathbb{R}
$$

Proof. By symmetry, we need only to consider the case $t \geq 0$. Then by definition, $X^{z}(t) \equiv X^{z,+}(t)$ satisfies

$$
\begin{equation*}
X^{z}(t)=z+W(t)-\int_{0}^{t} \nabla u_{0}\left(X^{z}(s)\right) d s, \quad t \geq 0 \tag{4.9}
\end{equation*}
$$

Here we simply write $W$ for $W^{+}$. Differentiating both sides of (4.9) with respect to $z_{i}$ and $t$ successively, we have

$$
\frac{d}{d t} Y(t)=-\operatorname{Hess}_{u_{0}}\left(X^{z}(t)\right) Y(t)
$$

where $Y=Y_{i}^{z}$. By (2.2), it then holds that

$$
\frac{1}{2} \frac{d}{d t}|Y(t)|^{2}=-\left(Y(t), \operatorname{Hess}_{u_{0}}\left(X^{z}(t)\right) Y(t)\right) \leq-\alpha|Y(t)|^{2}
$$

which shows $|Y(t)|^{2} \leq|Y(0)|^{2} e^{-2 \alpha t}=e^{-2 \alpha t}$. This ends the proof.

By Lemma 4.5, we obtain the following estimate on $\left|\nabla_{z} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)\right|$.
Lemma 4.6. Under the assumptions (A.1a)(i) and (A.2)(i), we have

$$
\begin{equation*}
\left|\nabla_{z} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)\right| \leq \sqrt{d} \int_{\Lambda^{2}}\left(\left\|X^{z}(s)\right\|^{q_{0}}+\left\|X^{z}(t)\right\|^{q_{0}}\right) \Psi_{\alpha}(d s, d t) \tag{4.10}
\end{equation*}
$$

where we set $\Psi_{\alpha}(d s, d t)=\psi_{0}(s-t) e^{-\alpha|s|} d s d t$.
Proof. Let $Y_{i}^{z}$ be as in Lemma 4.5. By the symmetry (2.1) of $\psi$, we see

$$
\frac{\partial}{\partial z_{i}} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)=\int_{\Lambda^{2}}\left(\nabla_{x} \psi\left(s-t, X^{z}(s), X^{z}(t)\right), Y_{i}^{z}(s)\right) d s d t
$$

for each $1 \leq i \leq d$. Then by (i) of (A.2) and Lemma 4.5, we have

$$
\left|\frac{\partial}{\partial z_{i}} \mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)\right| \leq \int_{\Lambda^{2}} \psi_{0}(s-t)\left(\left\|X^{z}(s)\right\|^{q_{0}}+\left\|X^{z}(t)\right\|^{q_{0}}\right) e^{-\alpha|s|} d s d t
$$

for all $1 \leq i \leq d$. Now the assertion follows readily.
We denote by $J\left(X^{z}\right)$ the RHS of (4.10). Combining this lemma and (4.8), we have

$$
\int_{\mathbb{R}^{d}} U(z) e^{-\mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)} m_{0}(d z) \leq \mathrm{a}_{3} \int_{\mathbb{R}^{d}}\|z\|^{p_{0}} J\left(X^{z}\right) e^{-\mathcal{H}_{\Lambda}^{\psi}\left(X^{z}\right)} m_{0}(d z)
$$

Plugging this estimate into (4.7), we obtain

$$
\begin{align*}
I & \leq \mathrm{a}_{3}\left\langle\|X(0)\|^{p_{0}} J(X)\right\rangle_{\mu_{\Lambda}} \\
& =\sqrt{d} \mathrm{a}_{3} \int_{\Lambda^{2}}\left\langle\|X(0)\|^{p_{0}}\left(\|X(s)\|^{q_{0}}+\|X(t)\|^{q_{0}}\right)\right\rangle_{\mu_{\Lambda}} \Psi_{\alpha}(d s, d t) \tag{4.11}
\end{align*}
$$

Since $p_{0}>q_{0}$, the following is immediate from the (generalized) Hölder inequality:
Lemma 4.7. Let $M_{\Lambda}=\max _{t \in \Lambda}\left\langle\|X(t)\|^{2 p_{0}}\right\rangle_{\mu_{\Lambda}}$. We have, for all $s \in \Lambda$,

$$
\left\langle\|X(0)\|^{p_{0}}\|X(s)\|^{q_{0}}\right\rangle_{\mu_{\Lambda}} \leq\left\langle\|X(0)\|^{2 p_{0}}\right\rangle_{\mu_{\Lambda}}^{\frac{1}{2}} M_{\Lambda}^{\frac{q_{0}}{2 p_{0}}}
$$

Now we are in a position to prove Proposition 4.3:
Proof of Proposition 4.3. By (4.11) and Lemma 4.7, we see that

$$
I \leq 2 \sqrt{d} \mathrm{a}_{3}\left\langle\|X(0)\|^{2 p_{0}}\right\rangle_{\mu_{\Lambda}}^{\frac{1}{\mu_{\Lambda}}} M_{\Lambda}^{\frac{q_{0}}{2 p_{0}}} \Psi_{\alpha}\left(\Lambda^{2}\right)
$$

Note that $\Psi_{\alpha}\left(\Lambda^{2}\right) \leq \Psi_{\alpha}\left(\mathbb{R}^{2}\right)=(2 / \alpha) \bar{\psi}_{0}$ by definition. Now we shift the origin to $t_{0}$ and obtain the proposition.

Remark 4.2. Since Lemma 4.7 remains true, the inequality (4.2) also holds in the case $p_{0}=q_{0}$. In this case the equation (4.3) has a unique solution $\mathrm{m}_{2}$ if $\mathrm{a}_{1}>\mathrm{a}_{3} \bar{\psi}_{0} C_{d, \alpha}$. Then (4.4) holds with $\mathrm{m}_{2}$, from which the existence of the associated Gibbs measures follows.

## 5 Construction of hard-wall Gibbs measures with interactions

In this section we discuss the existence of hard-wall Gibbs measures; that is, we consider an external potential $\varphi^{+}$such that $\varphi^{+}(x)=\infty$ for $x<0$.

In the following, boldfaced letters denote elements in $\mathbb{R}^{3}$; e.g., $\boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{3}$. We often write $\rho(\boldsymbol{x})$ for $|\boldsymbol{x}|$. For a function $f$ on $\mathbb{R}^{+},(f \circ \rho)(\boldsymbol{x})$ means $f(\rho(\boldsymbol{x}))(\equiv f(|\boldsymbol{x}|))$ as usual.

We summarize the assumptions on potentials:
(A.1') Assumptions on $\varphi^{+} . \varphi^{+}$is bounded from below. There exist $\varphi_{0}^{+}$and $\varphi_{1}^{+}$such that $\varphi^{+}=\varphi_{0}^{+}+\varphi_{1}^{+}$, taking values $\infty$ on $\{x<0\}$ and satisfying (A.1'a) and (A.1'b), respectively:
(A.1'a) $\varphi_{0}^{+} \circ \rho$ satisfies (A.1a) with $d=3$; that is, $\varphi_{0}^{+}$is continuous on $\{x \geq 0\}$, and the Schrödinger operator $\mathrm{H}_{0}^{\rho}=-(1 / 2) \Delta+\varphi_{0}^{+} \circ \rho$ acting on $L^{2}\left(\mathbb{R}^{3} ; d \boldsymbol{x}\right)$ has a strictly positive ground state $f_{0}$ of class $C^{2}\left(\mathbb{R}^{3}\right)$ satisfying (i) and (ii) of (A.1a).
(A. $\left.1^{\prime} \mathrm{b}\right) \varphi_{1}^{+} \circ \rho$ satisfies (A.1b) with $d=3$.
(A.2') Assumptions on $\boldsymbol{\psi}$. The function $\psi^{\rho}(s, \boldsymbol{x}, \boldsymbol{y}):=\psi(s, \rho(\boldsymbol{x}), \rho(\boldsymbol{y}))$ satisfies (A.2) with $d=3$.
(A.3') $p_{0}$ is strictly larger than $q_{0}$.

Remark 5.1. By definition, we see $\mathrm{f}_{0}$ in (A.1'a) is radially symmetric; that is, there exists a $g_{0}$ on $\mathbb{R}^{+}$such that $\mathrm{f}_{0}(\boldsymbol{x})=g_{0}(|\boldsymbol{x}|)$.

Now we state the existence of hard-wall Gibbs measures:
Theorem 5.1. Assume (A.1')-(A.3'). Then there exists a $\left(\varphi^{+}, \psi\right)$-Gibbs measure.
The existence of hard-wall Gibbs measures with interactions had been an open problem; [17] required a certain symmetry on external potentials, which excluded hard-wall ones. We have established the existence of such measures by Theorem 5.1.

To prove this theorem, we do the finite-volume construction in the same manner as Section 3: For $\varphi_{0}^{+}$, let $H_{0}^{+}$be the corresponding Schrödinger operator with Dirichlet boundary condition:

$$
\left\{\begin{array}{l}
H_{0}^{+} f=-\frac{1}{2} f^{\prime \prime}(x)+\varphi_{0}^{+}(x) f(x), \quad x>0 \\
f(0)=0
\end{array}\right.
$$

Let $g_{0}$ be as in Remark 5.1. As we see in Lemma 5.3, a ground state $f_{0}^{+}$of $H_{0}^{+}$is given by $f_{0}^{+}(x)=2 \sqrt{\pi} x g_{0}(x)$. The constant $2 \sqrt{\pi}$ is chosen so that $\int_{\mathbb{R}^{+}}\left|f_{0}^{+}\right|^{2} d x=\int_{\mathbb{R}^{3}}\left|\mathrm{f}_{0}\right|^{2} d \boldsymbol{x}$. The $P(\phi)_{1}$-process associated with $H_{0}^{+}$is determined by the SDE

$$
\begin{equation*}
d X(t)=d W(t)+\frac{d t}{X(t)}-v_{0}^{\prime}(X(t)) d t \tag{5.1}
\end{equation*}
$$

with $W$ a 1-dimensional Brownian motion. Here we set $v_{0}=-\log g_{0}$. The process $X$ is stationary under the probability measure $f_{0}^{+}(x)^{2} d x$ (we assume $f_{0}^{+}$is normalized). We denote by $\mathcal{P}^{\varphi_{0}^{+}}$the law on $C\left(\mathbb{R} ; \mathbb{R}^{+}\right)$of this process. For $\mathcal{P}^{\varphi_{0}^{+}}$and potentials $\varphi_{1}^{+}, \psi$, we define the finite-volume Gibbs measure $\mu_{\Lambda}^{\varphi^{+}, \psi}$ through (3.2). Similarly to Theorem 2.1, the existence result Theorem 5.1 follows from the localization result of $\left\{\mu_{\Lambda}^{\varphi^{+}, \psi}\right\}_{\Lambda \subset \mathbb{R}}$ :

Theorem 5.2. Assume (A.1'a), (A.1'b), (A.2') and (A.3'). Then the localization (*) for $\left\{\mu_{\Lambda}^{\varphi^{+}, \psi}\right\}_{\Lambda \subset \mathbb{R}}$ holds.

A key to Theorem 5.2 is the following identity in law:

$$
\begin{equation*}
\left(X, \mathcal{P}^{\varphi_{0}^{+}}\right) \stackrel{(d)}{=}\left(|\boldsymbol{X}|, \mathcal{P}^{\varphi_{0}^{+} \circ \rho}\right), \tag{5.2}
\end{equation*}
$$

where, on the RHS, $\mathcal{P}^{\varphi_{0}^{+} \circ \rho}$ denotes the law of the $P(\phi)_{1}$-process associated with $\mathrm{H}_{0}^{\rho}=$ $-(1 / 2) \Delta+\varphi_{0}^{+} \circ \rho$, and $\boldsymbol{X} \in C\left(\mathbb{R} ; \mathbb{R}^{3}\right)$. Here is an example:
Example 5.1. When $\varphi_{0}^{+}(x)=\frac{1}{2} a^{2} x^{2}, x \geq 0$, the $P(\phi)_{1}$-process $X$ associated with $H_{0}^{+}$ has the following explicit representation:

$$
\{X(t), t \in \mathbb{R}\} \stackrel{(d)}{=}\left\{\frac{1}{\sqrt{2 a}} e^{-a t} R^{(3)}\left(e^{2 a t}\right), t \in \mathbb{R}\right\}
$$

where $\left\{R^{(3)}(s), s \geq 0\right\}$ is a 3 -dimensional Bessel process starting at 0 . In view of the DLR equation, this can be seen from the fact that a Brownian bridge conditioned to be positive has the same law as a 3-dimensional Bessel bridge.

Once the identity (5.2) is verified, then Theorem 5.2 is straightforward:
Proof of Theorem 5.2. For $\varphi_{0}^{+} \circ \rho, \varphi_{1}^{+} \circ \rho$ and $\psi^{\rho}$, let $\mu_{\Lambda}^{\rho}$ be the finite-volume Gibbs measure on $C\left(\Lambda ; \mathbb{R}^{3}\right)$ defined via (3.2):

$$
d \mu_{\Lambda}^{\rho}(\boldsymbol{X})=\left(Z_{\Lambda}^{\rho}\right)^{-1} e^{-\mathcal{H}_{\Lambda}^{\rho}(\boldsymbol{X})} d \mathcal{P}_{\Lambda}^{\varphi_{0}^{+} \circ \rho}(\boldsymbol{X})
$$

Here $\mathcal{H}_{\Lambda}^{\rho}=\mathcal{H}_{\Lambda}^{\varphi_{1}^{+} \odot \rho, \psi^{\rho}}$ and $Z_{\Lambda}^{\rho}$ is the normalization. Note that $\mathcal{H}_{\Lambda}^{\rho}(\boldsymbol{X})=\mathcal{H}_{\Lambda}^{\varphi_{1}^{+}, \psi}(|\boldsymbol{X}|)$ by definition. From this and (5.2), we easily see that, for all $t \in \Lambda$,

$$
\begin{equation*}
\left.\left\langle X(t)^{2 p_{0}}\right\rangle_{\mu_{\Lambda}^{\varphi+}, \psi}=\left.\langle | \boldsymbol{X}(t)\right|^{2 p_{0}}\right\rangle_{\mu_{\Lambda}^{\rho}} . \tag{5.3}
\end{equation*}
$$

By the assumptions on potentials, we may use Theorem 3.1 to see that the localization $(*)$ for $\left\{\mu_{\Lambda}^{\rho}\right\}_{\Lambda \subset \mathbb{R}}$ holds with $q_{1}=2 p_{0}$. Combining this with (5.3) ends the proof.

The identity (5.2) is an immediate consequence of the following lemma:

Lemma 5.3. (i) Let $\left\{T_{t}^{+}\right\}_{t \geq 0}$ and $\left\{\mathrm{T}_{t}^{\rho}\right\}_{t \geq 0}$ be the semi-group generated by $H_{0}^{+}$and that generated by $\mathrm{H}_{0}^{\rho}$, respectively. Then it holds that, for $x>0$ and $\boldsymbol{x} \in \mathbb{R}^{3}$ with $|\boldsymbol{x}|=x$,

$$
\begin{equation*}
\left(T_{t}^{+} f\right)(x)=x\left(\mathbf{T}_{t}^{\rho} \frac{f \circ \rho}{\rho}\right)(\boldsymbol{x}), \quad t \geq 0 . \tag{5.4}
\end{equation*}
$$

In particular, a ground state of $H_{0}^{+}$is given by $x g_{0}(x)$.
(ii) Under $\mathcal{P}^{\varphi_{0}^{+} \circ \rho}$, the process $\{|\boldsymbol{X}(t)|, t \in \mathbb{R}\}$ satisfies the $\operatorname{SDE}$ (5.1). Moreover, $|\boldsymbol{X}(0)|$ has the same law as $\mathcal{P}^{\varphi_{0}^{+}}(X(0) \in \cdot)$.

Remark 5.2. The identity (5.4) is related to the $h$-transform of Doob, cf. [19, Chap. VIII]. Proof. (i) Let $\mathcal{W}_{x}$ be the law of 1-dimensional Brownian motion starting at $x$. Let $\tau_{0}=\inf \{t \geq 0 ; X(t)=0\}$. Then $T_{t}^{+}$is expressed as

$$
\left(T_{t}^{+} f\right)(x)=\left\langle f\left(X\left(t \wedge \tau_{0}\right)\right) \exp \left\{-\int_{0}^{t} \varphi_{0}^{+}\left(X\left(s \wedge \tau_{0}\right)\right) d s\right\}\right\rangle_{\mathcal{W}_{x}}, \quad x>0
$$

where $\wedge$ denotes the minimum. By the conditional equivalence between absorbing Brownian motion and 3 -dimensional Bessel process (see, e.g., [10, Lemma 5.2.8]), this is rewritten as, for $\boldsymbol{x} \in \mathbb{R}^{3}$ with $|\boldsymbol{x}|=x$,

$$
\begin{equation*}
\left(T_{t}^{+} f\right)(x)=x\left\langle\frac{f(|\boldsymbol{X}(t)|)}{|\boldsymbol{X}(t)|} \exp \left\{-\int_{0}^{t} \varphi_{0}^{+}(|\boldsymbol{X}(s)|) d s\right\}\right\rangle_{\mathcal{W}_{\boldsymbol{x}}^{(3)}} . \tag{5.5}
\end{equation*}
$$

Here $\mathcal{W}_{\boldsymbol{x}}^{(3)}$ is the law of 3-dimensional Brownian motion starting at $\boldsymbol{x}$. Recalling $\rho(\boldsymbol{x})=$ $|\boldsymbol{x}|$, we see that (5.5) shows (5.4). (ii) Noting $\nabla \log \mathrm{f}_{0}=-|\boldsymbol{x}|^{-1} v_{0}^{\prime}(|\boldsymbol{x}|) \boldsymbol{x}$, we see that the process ( $\boldsymbol{X}, \mathcal{P}^{\varphi_{0}^{+} \circ \rho}$ ) satisfies the following SDE:

$$
d \boldsymbol{X}(t)=d \boldsymbol{W}(t)-\frac{v_{0}^{\prime}(|\boldsymbol{X}(t)|)}{|\boldsymbol{X}(t)|} \boldsymbol{X}(t) d t
$$

where $\boldsymbol{W}$ is a 3-dimensional Brownian motion. By Itô's formula, we have

$$
d|\boldsymbol{X}(t)|=d \widetilde{W}(t)+\frac{d t}{|\boldsymbol{X}(t)|}-v_{0}^{\prime}(|\boldsymbol{X}(t)|) d t, \quad \widetilde{W}(t)=\int_{0}^{t}\left(\frac{\boldsymbol{X}(s)}{|\boldsymbol{X}(s)|}, d \boldsymbol{W}(s)\right)
$$

Since $\widetilde{W}$ is a Brownian motion, the former assertion is proved. For the latter, note that $\mathcal{P}^{\varphi_{0}^{+} \circ \rho}(\boldsymbol{X}(0) \in d \boldsymbol{x})=\mathrm{f}_{0}(\boldsymbol{x})^{2} d \boldsymbol{x}$, hence that $|\boldsymbol{X}(0)|$ is distributed as $f_{0}^{+}(x)^{2} d x$. This ends the proof.

Remark 5.3. Lemma 5.3 suggests that, in the Markovian case, namely, the case $\psi=0$, the existence of Gibbs measures for the hard-wall external potential $\varphi^{+}$is reduced to that of the ground state for the Schrödinger operator defined by

$$
\begin{equation*}
-\frac{1}{2} \Delta+\varphi^{+}(|\boldsymbol{x}|) \quad \text { on } L^{2}\left(\mathbb{R}^{3} ; d \boldsymbol{x}\right) \tag{5.6}
\end{equation*}
$$

Now let us consider the case where $\varphi^{+}$is a "single well" such as

$$
\varphi^{+}(x)=-\frac{\beta}{\left(1+x^{2}\right)^{\gamma}}, \quad x \geq 0
$$

for $\beta, \gamma>0$ (the attractive potential). The parameters $\beta$ and $\gamma$ control the "depth" and "width" of the well, respectively. Let $\mathrm{H}_{\beta} \equiv \mathrm{H}_{\beta, \gamma}$ be the Schrödinger operator defined by (5.6). Note that the essential spectrum of $\mathrm{H}_{\beta}$ is equal to $[0, \infty)$ (see, e.g., [18]). Let $\lambda_{\beta}$ be the bottom of the spectrum of $\mathrm{H}_{\beta}$. If $0<\gamma<1$, then $\lambda_{\beta}<0$ for all $\beta>0$ and $\lambda_{\beta}$ is a simple eigenvalue; that is, the associated ground state exists. On the other hand, if $\gamma \geq 1$, then there exists a $\beta=\beta_{\mathrm{c}}$ (the threshold) such that $\lambda_{\beta} \equiv 0$ for $\beta \leq \beta_{\mathrm{c}}$ and $\lambda_{\beta}<0$ for $\beta>\beta_{\mathrm{c}}$. See [6, Chap. 8]. This phenomenon may be regarded as a counterpart to wetting transitions in lattice models $[4,5]$. We will return to this somewhere else.

Acknowledgments. The author wishes to thank Professor H. Osada and Professor T. Funaki for their useful comments and their encouragement. He is also grateful to Professor J. Lőrinczi, who kindly sent him a preprint of [14] before its publication.

## References

[1] Betz, V.: Existence of Gibbs measures relative to Brownian motion. Markov Process. Related Fields 9, 85-102 (2003)
[2] Betz, V., Lőrinczi, J.: A Gibbsian description of $P(\phi)_{1}$-processes, preprint (1999)
[3] Bolthausen, E., Deuschel, J.-D., Zeitouni, O.: Entropic repulsion of the lattice free field. Commun. Math. Phys. 170, 417-443 (1995)
[4] Bolthausen, E., Deuschel, J.-D., Zeitouni, O.: Absence of a wetting transition for a pinned harmonic crystal in dimensions three and larger. J. Math. Phys. 41, 12111223 (2000)
[5] Caputo, A., Velenik, Y.: A note on wetting transition for gradient fields. Stoch. Proc. Appl. 87, 107-113 (2000)
[6] Davies, E.B.: Spectral Theory and Differential Operators. Cambridge Studies in Advanced Mathematics no.42: Cambridge University Press, 1995
[7] Funaki, T.: The reversible measures of multi-dimensional Ginzburg-Landau type continuum model. Osaka J. Math. 28, 463-494 (1991)
[8] Hohenberg, P.C., Halperin, B.I.: Theory of Dynamic critical Phenomena. Rev. Mod. Phys. 49, 435-479 (1977)
[9] Iwata, K.: Reversible measures of a $P(\phi)_{1}$-time evolution, in "Prob. Meth. in Math. Phys. (eds. K. Itô and N. Ikeda), Proc. of Taniguchi Symp.", 195-209 (1985)
[10] Knight, F.B.: Essentials of Brownian Motion and Diffusion. Mathematical Surveys no.18: AMS 1981
[11] Hariya, Y., Osada, H.: Diffusion processes on path spaces with interactions. Rev. Math. Phys. 13, 199-220 (2001)
[12] Lebowitz, J.L., Maes, C.: The effect of an external field on an interface, entropy repulsion. J. Stat. Phys. 46, 39-49 (1987)
[13] Lebowitz, J.L., Presutti, E.: Statistical mechanics of systems of unbounded spins. Commun. Math. Phys. 50, 195-218 (1976)
[14] Lőrinczi, J., Minlos, J.A.: Gibbs measures for Brownian paths under the effect of an external and a small pair potential. J. Statist. Phys. 105, 605-647 (2001)
[15] Lőrinczi, J., Minlos, J.A., Spohn, H.: The infrared behavior in Nelson's model of a quantum particle coupled to a massless scalar field. Ann. Henri Poincaré 3, 269-295 (2002)
[16] Nelson, E.: Interaction of nonrelativistic particles with a quantized scalar field. J. Math. Phys. 5, 1190-1197 (1964)
[17] Osada, H., Spohn, H.: Gibbs measures relative to Brownian motion. Ann. Probab. 27, 1183-1207 (1999)
[18] Reed, M., Simon, B.: Methods of Modern Mathematical Physics, vol.IV: Analysis of Operators. New York: Academic Press, 1978
[19] Revuz, D., Yor, M.: Continuous Martingales and Brownian Motion. 2nd ed. Berlin: Springer-Verlag, 1994
[20] Rosen, J., Simon, B.: Fluctuations in $P(\phi)_{1}$ processes. Ann. Probab. 4, 155-174 (1976)
[21] Ruelle, D.: Superstable interactions in classical statistical mechanics. Commun. Math. Phys. 18, 127-159 (1970)
[22] Simon, B.: Functional Integration and Quantum Physics. New York: Academic Press, 1979
[23] Spohn, H.: Large Scale Dynamics of Interacting Particles. Berlin: Springer-Verlag 1991

# List of MHF Preprint Series, Kyushu University <br> 21st Century COE Program <br> Development of Dynamic Mathematics with High Functionality 

MHF2003-1 Mitsuhiro T. NAKAO, Kouji HASHIMOTO \& Yoshitaka WATANABE
A numerical method to verify the invertibility of linear elliptic operators with applications to nonlinear problems

MHF2003-2 Masahisa TABATA \& Daisuke TAGAMI
Error estimates of finite element methods for nonstationary thermal convection problems with temperature-dependent coefficients

MHF2003-3 Tomohiro ANDO, Sadanori KONISHI \& Seiya IMOTO
Adaptive learning machines for nonlinear classification and Bayesian information criteria

MHF2003-4 Kazuhiro YOKOYAMA
On systems of algebraic equations with parametric exponents
MHF2003-5 Masao ISHIKAWA \& Masato WAKAYAMA
Applications of Minor Summation Formulas III, Plücker relations, Lattice paths and Pfaffian identities

MHF2003-6 Atsushi SUZUKI \& Masahisa TABATA
Finite element matrices in congruent subdomains and their effective use for large-scale computations

MHF2003-7 Setsuo TANIGUCHI
Stochastic oscillatory integrals - asymptotic and exact expressions for quadratic phase functions -

MHF2003-8 Shoki MIYAMOTO \& Atsushi YOSHIKAWA
Computable sequences in the Sobolev spaces
MHF2003-9 Toru FUJII \& Takashi YANAGAWA
Wavelet based estimate for non-linear and non-stationary auto-regressive model
MHF2003-10 Atsushi YOSHIKAWA
Maple and wave-front tracking - an experiment
MHF2003-11 Masanobu KANEKO
On the local factor of the zeta function of quadratic orders
MHF2003-12 Hidefumi KAWASAKI
Conjugate-set game for a nonlinear programming problem

MHF2004-1 Koji YONEMOTO \& Takashi YANAGAWA
Estimating the Lyapunov exponent from chaotic time series with dynamic noise

MHF2004-2 Rui YAMAGUCHI, Eiko TSUCHIYA \& Tomoyuki HIGUCHI
State space modeling approach to decompose daily sales of a restaurant into time-dependent multi-factors

MHF2004-3 Kenji KAJIWARA, Tetsu MASUDA, Masatoshi NOUMI, Yasuhiro OHTA \& Yasuhiko YAMADA
Cubic pencils and Painlevé Hamiltonians
MHF2004-4 Atsushi KAWAGUCHI, Koji YONEMOTO \& Takashi YANAGAWA
Estimating the correlation dimension from a chaotic system with dynamic noise

MHF2004-5 Atsushi KAWAGUCHI, Kentarou KITAMURA, Koji YONEMOTO, Takashi YANAGAWA \& Kiyofumi YUMOTO
Detection of auroral breakups using the correlation dimension
MHF2004-6 Ryo IKOTA, Masayasu MIMURA \& Tatsuyuki NAKAKI
A methodology for numerical simulations to a singular limit
MHF2004-7 Ryo IKOTA \& Eiji YANAGIDA
Stability of stationary interfaces of binary-tree type
MHF2004-8 Yuko ARAKI, Sadanori KONISHI \& Seiya IMOTO
Functional discriminant analysis for gene expression data via radial basis expansion

MHF2004-9 Kenji KAJIWARA, Tetsu MASUDA, Masatoshi NOUMI, Yasuhiro OHTA \& Yasuhiko YAMADA
Hypergeometric solutions to the $q$ - Painlevé equations
MHF2004-10 Raimundas VIDŪNAS
Expressions for values of the gamma function
MHF2004-11 Raimundas VIDŪNAS
Transformations of Gauss hypergeometric functions
MHF2004-12 Koji NAKAGAWA \& Masakazu SUZUKI
Mathematical knowledge browser
MHF2004-13 Ken-ichi MARUNO, Wen-Xiu MA \& Masayuki OIKAWA
Generalized Casorati determinant and Positon-Negaton-Type solutions of the Toda lattice equation

MHF2004-14 Nalini JOSHI, Kenji KAJIWARA \& Marta MAZZOCCO
Generating function associated with the determinant formula for the solutions of the Painlevé II equation

MHF2004-15 Kouji HASHIMOTO, Ryohei ABE, Mitsuhiro T. NAKAO \& Yoshitaka WATANABE
Numerical verification methods of solutions for nonlinear singularly perturbed problem

MHF2004-16 Ken-ichi MARUNO \& Gino BIONDINI
Resonance and web structure in discrete soliton systems: the two-dimensional Toda lattice and its fully discrete and ultra-discrete versions

MHF2004-17 Ryuei NISHII \& Shinto EGUCHI
Supervised image classification in Markov random field models with Jeffreys divergence

MHF2004-18 Kouji HASHIMOTO, Kenta KOBAYASHI \& Mitsuhiro T. NAKAO
Numerical verification methods of solutions for the free boundary problem
MHF2004-19 Hiroki MASUDA
Ergodicity and exponential $\beta$-mixing bounds for a strong solution of Lévydriven stochastic differential equations

MHF2004-20 Setsuo TANIGUCHI
The Brownian sheet and the reflectionless potentials
MHF2004-21 Ryuei NISHII \& Shinto EGUCHI
Supervised image classification based on AdaBoost with contextual weak classifiers

MHF2004-22 Hideki KOSAKI
On intersections of domains of unbounded positive operators
MHF2004-23 Masahisa TABATA \& Shoichi FUJIMA
Robustness of a characteristic finite element scheme of second order in time increment

MHF2004-24 Ken-ichi MARUNO, Adrian ANKIEWICZ \& Nail AKHMEDIEV
Dissipative solitons of the discrete complex cubic-quintic Ginzburg-Landau equation

MHF2004-25 Raimundas VIDŪNAS
Degenerate Gauss hypergeometric functions
MHF2004-26 Ryo IKOTA
The boundedness of propagation speeds of disturbances for reaction-diffusion systems

MHF2004-27 Ryusuke KON
Convex dominates concave: an exclusion principle in discrete-time Kolmogorov systems

Multiple attractors in host-parasitoid interactions: coexistence and extinction
MHF2004-29 Kentaro IHARA, Masanobu KANEKO \& Don ZAGIER
Derivation and double shuffle relations for multiple zeta values
MHF2004-30 Shuichi INOKUCHI \& Yoshihiro MIZOGUCHI
Generalized partitioned quantum cellular automata and quantization of classical CA

MHF2005-1 Hideki KOSAKI
Matrix trace inequalities related to uncertainty principle
MHF2005-2 Masahisa TABATA
Discrepancy between theory and real computation on the stability of some finite element schemes

MHF2005-3 Yuko ARAKI \& Sadanori KONISHI
Functional regression modeling via regularized basis expansions and model selection

MHF2005-4 Yuko ARAKI \& Sadanori KONISHI
Functional discriminant analysis via regularized basis expansions
MHF2005-5 Kenji KAJIWARA, Tetsu MASUDA, Masatoshi NOUMI, Yasuhiro OHTA \& Yasuhiko YAMADA
Point configurations, Cremona transformations and the elliptic difference Painlevé equations

MHF2005-6 Kenji KAJIWARA, Tetsu MASUDA, Masatoshi NOUMI, Yasuhiro OHTA \& Yasuhiko YAMADA
Construction of hypergeometric solutions to the $q$ - Painlevé equations
MHF2005-7 Hiroki MASUDA
Simple estimators for non-linear Markovian trend from sampled data:
I. ergodic cases

MHF2005-8 Hiroki MASUDA \& Nakahiro YOSHIDA
Edgeworth expansion for a class of Ornstein-Uhlenbeck-based models
MHF2005-9 Masayuki UCHIDA
Approximate martingale estimating functions under small perturbations of dynamical systems

MHF2005-10 Ryo MATSUZAKI \& Masayuki UCHIDA
One-step estimators for diffusion processes with small dispersion parameters from discrete observations

MHF2005-11 Junichi MATSUKUBO, Ryo MATSUZAKI \& Masayuki UCHIDA
Estimation for a discretely observed small diffusion process with a linear drift

MHF2005-12 Masayuki UCHIDA \& Nakahiro YOSHIDA
AIC for ergodic diffusion processes from discrete observations
MHF2005-13 Hiromichi GOTO \& Kenji KAJIWARA
Generating function related to the Okamoto polynomials for the Painlevé IV equation

MHF2005-14 Masato KIMURA \& Shin-ichi NAGATA
Precise asymptotic behaviour of the first eigenvalue of Sturm-Liouville problems with large drift

MHF2005-15 Daisuke TAGAMI \& Masahisa TABATA
Numerical computations of a melting glass convection in the furnace
MHF2005-16 Raimundas VIDŪNAS
Normalized Leonard pairs and Askey-Wilson relations
MHF2005-17 Raimundas VIDŪNAS
Askey-Wilson relations and Leonard pairs
MHF2005-18 Kenji KAJIWARA \& Atsushi MUKAIHIRA
Soliton solutions for the non-autonomous discrete-time Toda lattice equation
MHF2005-19 Yuu HARIYA
Construction of Gibbs measures for 1-dimensional continuum fields


[^0]:    *Research Institute for Mathematical Sciences, Kyoto University, Sakyo-ku, Kyoto 606-8502, Japan. E-mail: hariya@kurims.kyoto-u.ac.jp

