# Quantum Statistical Mechanics for Superstable Interactions: Bose-Einstein Statistics 

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#### Abstract

We study quantum systems of interacting Bose particles confined to a bounded region $A$ of $\mathbf{R}^{\nu}$. For any superstable and (strong) lower regular interaction, we obtain uniform bounds on the expectations of exponentials of local number operators for any activity and for any temperature. The method we use here is an improvement over our previous method on the same subject. As a consequence of the bounds, any infinite volume limit states are entire analytic and locally normal. Furthermore under an integrability condition on the interaction, the limit states are modular states. In this case, we use the Green's function method to construct an infinite volume limit Hilbert spece, a strongly continuous time evolution group of unitary operators and an invariant vector. Moreover we prove the existence of the pressure and its independence of boundary conditions.


KEY WORDS: Superstable interaction; Bose-Einstein statistics; local number operators; locally normal states; modular states; Wiener integral formalism; pressure; boundary conditions.

## 1. INTRODUCTION

In this paper we continue our study, initiated in Ref. 10 , of statistical systems of quantum particles obeying Bose-Einstein statistics and interacting via a superstable and lower regular potential. For the classical systems with superstable interactions, Ruelle established uniform bounds on the finite volume correlation functions. ${ }^{(14)}$ Using the bounds he obtained various results on the infinite volume equilibrium states and the pressure. In Ref. 10 we have extended partially Ruelle's results to the quantum systems of Bose particles. Because of technical difficulties we have assumed a condition stronger than superstability on the interactions,

[^0]namely, the strong superstability. Although the results can be applied to physically interesting various interactions such as Lennard-Jones types, ${ }^{(10)}$ the method was not suitable to deal with all superstable interactions. Recently we came to know that Esposito, Nicolo, and Pulvirenti have extended the classical results to the quantum systems obeying Maxwell-Boltzmann statistics. ${ }^{(2)}$ In this paper we modify and improve the method we developed in Ref. 10 to include the basic ideas in Ref. 2 so that the main results in Ref. 10 hold for all superstable interactions. Thus the result of this paper is the answer to the problem (c) we proposed in the Introduction of Ref. 10. The improved method also enables us to prove the existence of the pressure and its independence of boundary conditions.

It may be worthwhile to remark on some results in quantum statistical mechanics. There have been extensive studies on the thermodynamic limit in statistical mechanics of interacting quantum systems, and fairly satisfactory results have been obtained for the thermodynamic functions such as pressure. ${ }^{(1,9,13)}$ On the other hand, the results concerning the equilibrium states for such systems are less satisfactory. In the dilute regime, detailed properties of the infinite volume (thermodynamic) limit states have been obtained for various classes of interactions. ${ }^{(1,3,4,6,16)}$ The charge conjugate invariant systems have also been studied in Ref. 5. It has to be pointed out that our approach is not perturbative, so it works at all values of temperatures and activities. In passing we mention that Ruelle's extimates ${ }^{(14)}$ have been extended to unbounded classical and quantum spin systems. ${ }^{(15,8,11)}$

We present a brief discussion on the main results of this paper below. Let $\rho_{A}$ be the finite volume Gibbs state for a system of interacting Bose particles confined to a bounded region $A$ of the configuration space $\mathbf{R}^{v}$ and let $N_{B}$ be the local number operator for $B \subset A$. We assume that the interaction satisfies the superstability and (strong) lower regularity conditions (see Section 2 for the definitions). We then obtain the bounds of the form ( $\nu<4$ if the potential is not repulsive):

$$
\rho_{A}\left(\exp \left[a N_{B}\right]\right) \leqslant \exp [A(B, a)]
$$

where $A(B, a)$ is a constant depending only on the diameter of $B$ and $|a|$. We will give the exact definitions of the models, conditions on the interactions, and the main results in Section 2. As a consequence of the bounds it follows that any limit point $\rho$ of the state $\rho_{A}$ in the weak* topology is entirely analytic and locally normal in the CCR algebra. Furthermore under an additional assumption on the integrability of the potential, the state $\rho$ is a modular state. In this case we use the Green's function method ${ }^{(1)}$ to construct an infinite volume physical Hilbert space, a strongly
continuous time evolution group of unitary operators and invariant vector(s). We also show the existence of the pressure and its independence of boundary conditions.

We next give a breif description of the basic ingredients of the method we use in this paper. The general strategy is essentially the same as that of Ref. 10. As in Ref. 10 we will use the Ginibre's representation of the partition function in terms of the Wiener integrals, ${ }^{(6,1,5,7)}$ and a modification of Ruelle's method used for classical systems. ${ }^{(14)}$ For quantum systems one has to control (a) quantum statistics and (b) interactions between two group of Wiener paths (W terms). Compared to the classical systems, these problems originate from the fluctuations (deviations) of Wiener paths. As in Ref. 10 we will use the fact that the Wiener measure of a subset of the paths with large deviations is small. On the other hand, the system behaves like a classical one on a subset with small deviations. The basic idea is to decompose the Wiener space $\Omega^{n}$ of $n$ paths into disjoint subsets according to amount of fluctuations. On each subset we utilize the above facts to obtain uniform bounds. The method we developed in Ref. 10 is strong enough to control quantum statistics (B-E statistics), but not strong enough to control W terms for all superstable interactions. This was the reason why we assumed a stronger condition than superstability. On the other hand the method developed in Ref. 2 is sufficiently good to handle interaction terms for all superstable interactions, yet the method cannot be applied directly to solve quantum statistics. Thus we combine and improve two methods to control (a) and (b) for any superstable interactions.

We are unable to obtain pointwise bounds on reduced density matrices. But we believe the method we used can be extended to give such bounds. By establishing decay properties one ought to be able to see whether or not the systems do exhibit the Bose-Einstein condensations. As mentioned in Ref. 10, the Wiener integral formalism does seem to be unsuitable to handle the systems of interacting Fermi particles. Thus it should be desirable to develop a method in the pure operator approach, i.e., the second quantization formalism. That kind of a method should be useful to study composite systems such as matter systems.

We try to make this paper as self-contained as possible. Therefore, as we follow the general procedures of Ref. 10, it is necessary both to introduce some new notations and to repeat some arguments from Ref. 10. To make the paper a reasonable size we freely use some technical results of Ref. 10 without the proofs whenever its proofs do not affect main arguments of the paper.

The contents of the paper are as follows: In Section 2.1 We introduce notations and the definitions of the models. In Section 2.2 we list the assumption on the interactions (Assumption A) and then give our main
estimate (Proposition 2.2.1). Then, the entire analyticity and local normality of any infinite volume limit state follow from the main estimate and the arguments used in Refs. 10 and 1. Under an assumption of the interactions (an integrability assumption) we use the main estimate and the Green's function method to construct an infinite volume limit theory (Theorem 2.2.3).

Sections 3 and 4 are devoted to the proof of the main estimate. In Section 3.1 we review the Wiener integral formalism in statistical mechanics and we then introduce a decomposition of the space of Wiener paths into multually disjoint subsets. With an assumption of one estimate (Theorem 3.1.1) we control quantum statistics in Section 3.2. In Section 4, we prove Theorem 3.1.1 and so we complete the proof of our main estimate. In Section 5 we finally show the existence of the pressure in the infinite volume limit and its independence of boundary conditions.

Before closing the Introduction we would like to make an apology to the readers for some printing (and typing) errors in Ref. 10. Although the correct meanings are obvious if one reads carefully, we make corrections of those errors in the Appendix of this paper.

## 2. THE DEFINITION OF MODELS AND THE MAIN RESULTS

### 2.1. Some Notations and Definitions

We first introduce the Hilbert space and the finite volume Gibbs states for a system of interacting Bose particles confined to a bounded open region $A$ of the configuration space $\mathbf{R}^{v}$. Let

$$
\begin{equation*}
\mathscr{H}_{n}^{(s)}(\Lambda)=\bigotimes_{i=1}^{n} L^{2}\left(\Lambda, d^{v} x_{i}\right) \tag{2.1.1}
\end{equation*}
$$

be the subspace of $L^{2}\left(\Lambda^{n}, d^{n v} x\right)$ formed by the totally symmetric functions of $n$ variables $x_{i} \in A$. The associated Fock space

$$
\begin{equation*}
\mathrm{F}^{(s)}(\Lambda)=\bigoplus_{n=0} \mathscr{H}_{n}^{(s)}(\Lambda) \tag{2.1.2}
\end{equation*}
$$

describe the states of an arbitrary number of particles. The total Hamiltonian is given by

$$
\begin{equation*}
H_{\Lambda}=\oplus_{n \geqslant 0} H_{\Lambda}^{(n)} \tag{2.1.3}
\end{equation*}
$$

in terms of the $n$-particle Hamiltonian $H_{A}^{(n)}$ which has the form

$$
\begin{equation*}
H_{A}^{(n)}=-\frac{1}{2} \sum_{i=1}^{n} A_{A, i}+U\left((x)_{n}\right) \tag{2.1.4}
\end{equation*}
$$

where $\Delta_{A, i}$ is the Laplacian in the variable $x_{i}$ with 0-Dirichlit boundary conditions on the boundary $\partial A$ of $A$ and the interaction energy of $n$ particles at the point $(x)_{n}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ is given by

$$
\begin{equation*}
U\left((x)_{n}\right)=\sum_{1 \leqslant i<j \leqslant n} \Phi\left(x_{i}-x_{j}\right) \tag{2.1.5}
\end{equation*}
$$

where $\Phi$ is a two-body potential between particles. If the interaction satisfies the stability condition

$$
U\left((x)_{n}\right) \geqslant-B n, \quad B>0
$$

for all $n$ and $(x)_{n} \in \Lambda^{n}, H_{A}^{(n)}$ is a self-adjoint operator on $\mathscr{H}_{n}^{(s)}(\Lambda)$, and so the total Hamiltonian $H_{A}$ is a self-adjoint operator on $\left.\mathrm{F}^{(s)}(\Lambda)\right)^{(1,13)}$

Let $B \subset A$ be a bound region in $\mathbf{R}^{\nu}$. We define the local number operators $N_{B}$ by

$$
\begin{equation*}
\left.N_{B}\right|_{\mathscr{H}_{n}^{(n)}(A)} \psi\left(x_{1}, \ldots, x_{n}\right)=\sum_{i=1}^{n} \chi_{B}\left(x_{i}\right) \psi\left(x_{1}, \ldots, x_{n}\right) \tag{2.1.6}
\end{equation*}
$$

where $\chi_{B}$ is the characteristic function of $B$. We note that for any $\psi \in \mathscr{H}_{n}^{(s)}(\Lambda), N_{A} \psi=n \psi$. The partition function and the finite volume Gibbs states are defined by

$$
\begin{equation*}
\Xi_{A}=\operatorname{Tr}_{\mathrm{F}^{(s)}(A)}\left(\exp \left[-\beta\left(H_{A}-\mu N_{A}\right)\right]\right) \tag{2.1.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{A}(A)=\Xi_{A}^{-1} \operatorname{Tr}_{\mathrm{F}^{(s)}(A)}\left(A \exp \left[-\beta\left(H_{A}-\mu N_{A}\right)\right]\right) \tag{2.1.8}
\end{equation*}
$$

respectively, for any $\beta \in \mathbf{R}^{+}, \mu \in \mathbf{R}$, and for any bounded operator $A$ on $F^{(s)}(A)$. The above are well defined under the stability condition for any bounded $A \subset \mathbf{R}^{v}{ }^{(1,13)}$

### 2.2. The assumptions on the Interactions and the Main Results

We now list the assumptions on the interactions. We shall assume that the interaction between particles is given by a pair potential $\Phi$ :

$$
U\left(x_{1}, \ldots, x_{n}\right)=\sum_{1 \leqslant i<j \leqslant n} \Phi\left(x_{i}-x_{j}\right)
$$

where $\Phi$ is a Lebesgue measurable function which satisfies $\Phi(x)=\Phi(-x)$ and which may take real values and the value $+\infty$. For every $r \in Z^{v}$ we define a half-open cube with side 1 ,

$$
\begin{equation*}
Q(r)=\left\{x \in \mathbf{R}:\left(r^{i}-\frac{1}{2}\right) \leqslant x^{i}<\left(r^{i}+\frac{1}{2}\right)\right\} \tag{2.2.1}
\end{equation*}
$$

These cubes form a partition of $\mathbf{R}^{v}$. If $X=\left(x_{1}, \ldots, x_{m}\right) \in\left(\mathbf{R}^{v}\right)^{m}$, we let $n(X, r)$ be the number of points of the sequence $X=\left(x_{1}, \ldots, x_{m}\right)$ which belong to $Q(r)$.

Assumption A. (a) Superstability. There exist $A>0, B>0$ such that if $\mathscr{R}$ is a finite subset of $Z^{\nu}$ and

$$
x_{1}, \ldots, x_{m} \in \bigcup_{r \in \Re} Q(r), \quad X=\left(x_{1}, \ldots, x_{m}\right)
$$

then

$$
\begin{equation*}
U(X) \geqslant \sum_{r \in \mathscr{R}}\left[A n(X, r)^{2}-B n(X, r)\right] \tag{2.2.2}
\end{equation*}
$$

(b) Strong Lower Regularity. There is a positive decreasing function $\psi$ on $(0,+\infty)$ such that

$$
\begin{equation*}
\int_{0}^{\infty} t^{\nu+\mu-1} \psi(t) d t<\infty \tag{2.2.3}
\end{equation*}
$$

where $\mu>1 / 2$, and for any $x \in \mathbf{R}^{v}$

$$
\Phi(x) \geqslant-\psi(|x|)
$$

Remark. (a) In Ref. 10 we have assumed a stronger condition than superstability, namely, the strong superstability:

$$
U(X) \geqslant \sum_{r \in \mathscr{R}}\left[A n(X, r)^{p^{\prime}}-B n(X, r)\right]
$$

with

$$
\frac{p^{\prime}}{p^{\prime}-1}<1+\frac{2\left(p^{\prime}-1\right)}{v p^{\prime}+(2-v)}
$$

Thus our results in this paper are improvements of those in Ref. 10.
(b) If $\Phi$ is a positive measurable function with the property that there exist constants $c>0$ and $d>0$ such that $\Phi(x)>c$ for $|x| \leqslant d$, then Assumption A is automatically satisfied. ${ }^{(13,14)}$

We first give our main estimate:
Proposition 2.2.1. Let the interaction satisfy Assumption $A$ and let $v<4$ if the potential is not repulsive. For given $\beta \in(0, \infty)$ and $\mu \in \mathbf{R}$, let $\rho_{A}$ be the finite volume Gibbs states defined in (2.1.7) for the interaction. Then for any $B \subset A$ and $a \in \mathbf{R}$ there is a constant $A(B, a)$ such that

$$
\rho_{A}\left(\exp \left\{a N_{B}\right\}\right) \leqslant \exp A(B, a)
$$

holds, where $A(B, a)$ depends only on $\operatorname{diam}(B)$ and $|a|$.

We will produce the proof of the above proposition in Sections 3 and 4.

We next consider some consequences of Proposition 2.2.1. We introduce the CCR algebra of the local observables. For the details we refer to Refs. 1 and 13. For $f \in L^{2}\left(\Lambda, d^{v} x\right)$, let $a(f)$ and $a(f)^{*}$ be the annihilation and creation operators defined on $\mathrm{F}^{(s)}(\Lambda)$. Then $\Phi(f)=(1 / \sqrt{2})\left[a(f)+a(f)^{*}\right]$ for real $f$ has a self-adjoint extension, which we write again $\Phi(f)$. We write the Weyl operator by $W(f)=\exp [i \Phi(f)]$, and let $\mathrm{A}(\Lambda)$ be the $C^{*}$ algebra generated by the Weyl operator $W(f), f \in L^{2}\left(A, d^{v} x\right)$, and let

$$
\begin{equation*}
\mathrm{A}=\overline{\bigcup_{\Lambda \subset \mathbf{R}^{v}} \mathrm{~A}(\Lambda)} \tag{2.2.4}
\end{equation*}
$$

be the quasi local CCR algebra in sense of Ref. 1. To study the infinite volume equilibrium states we follow the Green's function method used in Chapter 6.3 of Ref. 1. Let $\alpha^{A}$ be the time evolution automorphism on $A_{A}$ given by

$$
\begin{equation*}
\alpha_{t}^{\Lambda}(B)=e^{i t H_{A}} B e^{-i t H_{A}} \tag{2.2.5}
\end{equation*}
$$

We then define the finite volume Green's functions by

$$
\begin{equation*}
G_{A}(A, B ; t)=\rho_{A}\left(A \alpha_{t}^{A}(B)\right) \tag{2.2.6}
\end{equation*}
$$

Although $\rho_{A}$ is defined as a state over $A_{A}$, it has an extension to a state on A by the Hahn-Banach theorem which we denote again $\rho_{A}$. The bounds

$$
\begin{equation*}
\left|G_{A}(A, B ; t)\right| \leqslant\|A\|\|B\| \tag{2.2.7}
\end{equation*}
$$

imply that there exists a subnet $\left\{A_{\alpha}\right\}$ such that

$$
\begin{equation*}
G(A, B ; t)=\lim _{A_{\mathrm{x}} \rightarrow \mathbf{R}^{v}} G_{A}(A, B ; t) \tag{2.2.8}
\end{equation*}
$$

for all $A, B \in \mathrm{~A}, t \in \mathbf{R}$. This is a consequence of Tychonoff's theorem. Clearly the value

$$
\begin{equation*}
\rho(A)=G(A, 1 ; 0) \tag{2.2.9}
\end{equation*}
$$

determines a state $\rho$ over the quasilocal algebra $A$.
Theorem 2.2.2. Under the assumptions stated in Proposition 2.2.1, any weak* limit $\rho$ of $\rho_{A}$ defines an entire analytic state over the CCR algebra $A$. The state $\rho$ has finite local particle density and hence is locally normal.

Proof. Using Proposition 2.2.1 and following the method used in the proof of Theorem II. 3.2 of Ref. 10 one may get the bounds of the form

$$
\begin{equation*}
\rho_{A}\left(\prod_{i=1}^{k} N_{B_{i}}\right) \leqslant c k!\prod_{i=1}^{k} N^{+}\left(B_{i}\right) \tag{2.2.10}
\end{equation*}
$$

where $N^{+}(B)$ is the minimal numner of unit cubes which cover $B$. The theorem follows from (2.2.10) and from the method used in the proof of Theorem II.3.3 of Ref. 10. See also Section 6.3 of Ref. 1.

We next construct an infinite volume physical Hilbert space, a strongly continuous time evolution group of unitary operators, and a time translation invariant vector. Let the two-body potential $\Phi$ satisfy the following estimate: For some $\varepsilon>0$

$$
\begin{equation*}
\int d^{v} x\left(1+|x|^{2}\right)^{v+\varepsilon}|\Phi(x)|^{2}<\infty \tag{2.2.11}
\end{equation*}
$$

Let $G(A, B ; t)$ be the Green's functions defined in (2.2.8). From the bound (2.2.10), (2.2.11) and the methods used in the proofs of Proposition 6.3.29 and Theorem 6.3 .31 of Ref. 1 , it follows that $G: A \times A \times \mathbf{R} \rightarrow \mathbf{C}$ satisfies the following properties:
(1) $A, B \rightarrow G(A, B ; t)$ is bilinear for all $t \in \mathbf{R}$.
(2) $t \rightarrow G(A, B ; t)$ is continuous for all $A, B \in \mathrm{~A}$.
(3) $G(A, C B ; 0)=G(A C, B ; 0)$ for all $A, B, C \in \mathrm{~A}$.
(4) $G(1,1 ; 0)=1$.
(5) $\sum_{i, j=1}^{n} G\left(A_{i}^{*}, A_{j} ; t_{j}-t_{i}\right) \geqslant 0$ for any finite sequences $\left\{A_{i}\right\}_{i=1}^{n}$ in A and $\left\{t_{i}\right\}_{i=1}^{n}$ in $\mathbf{R}$.
(6) [Weak KMS conditions] For all $A, B \in \mathrm{~A}$, and for all $\hat{f} \in \mathrm{D}$ $\int d t f(t) G(A, B ; t)=\int d t f(t+i \beta) G(B, A ;-t)$.
The following result follows from Theorems 6.3.27 and 6.3.28 of Ref. 1.
Theorem 2.2.3. Let the interaction satisfy Assumption A and let $v<4$ if the potential is not repulsive. Assume that the potential satisfies the bound in (2.2.11). Let ( $\mathscr{H}_{\rho}, \pi_{\rho}, \Omega_{\rho}$ ) be the cyclic representation of A with respect to the state defined in (2.2.9). Then there exists a Hilbert space $\mathscr{H}$ containing $\mathscr{H}_{\rho}$, a strongly continuous representation $U$ of $\mathbf{R}$ such that
(i) $\mathscr{H}=V_{t \in \mathbf{R}} U(t) \mathscr{H}_{\rho}$;
(ii) $G(A, B ; t)=\left(\pi_{\rho}\left(A^{*}\right) \Omega_{\rho}, U(t) \pi_{\rho}(B) \Omega_{\rho}\right), A, B \in \mathrm{~A}, t \in \mathbf{R}$;
(iii) $\rho$ is a modular state on A , i.e., $\Omega_{\rho}$ is separating $\pi_{\rho}(\mathrm{A})^{\prime \prime}$.

Furthermore (i) and (ii) determine ( $\mathscr{H}, U$ ) uniquely up to unitary equivalence.

Since the proof of the theorem is same as those in the proofs of Theorems 6.3.27, 6.3.28, and 6.3 .31 of Ref. 1 , we do not produce the proof here and refer to Ref. 1.

Finally we discuss the existence of the infinite volume limit pressure and its independence of boundary conditions. Let $\Delta_{\sigma, A}$ be the Laplacian with boundary conditions $\partial \psi / \partial n=\sigma \psi$ where $\sigma \in C^{1}(\partial A)$. Dirichlet boundary conditions are corresponded to $\sigma=\infty$. In this case we have written $A_{\infty, A}$ by $\Delta_{A}$. Neumann boundary conditions are corresponded to $\sigma=0$. Let $\Xi_{\sigma, A}$ be partition functions obtained by replacing $\Delta_{A}$ by $\Delta_{\sigma, A}$ in the definition $\Xi_{A}$ in (2.1.7). Then, it follows from the minimax theorem that

$$
\begin{equation*}
\Xi_{A}=\Xi_{\infty, A} \leqslant \Xi_{\sigma, A} \leqslant \Xi_{0, A} \tag{2.2.12}
\end{equation*}
$$

For the detailed discussion on boundary conditions, we refer to Section 6.3 of Ref. 1.

The finite volume pressure is defined by

$$
\begin{equation*}
P_{\sigma, A}=\frac{1}{|A|} \log \Xi_{\sigma, A} \tag{2.2.13}
\end{equation*}
$$

We now state the results on the pressure:
Theorem 2.2.4. Let the interaction satisfy Assumption A. Assume that the two-body potential $\Phi$ satisfies the following decay property: there is a positive function $\tilde{\varphi}$ on $(0, \infty)$ such that $|\Phi(|x|)| \leqslant \tilde{\varphi}(|x|)$ for $|x| \geqslant b$ for some $b>0$ and

$$
\begin{equation*}
\int_{b}^{\infty} \tilde{\varphi}(t) t^{v-1} d t<\infty \tag{2.2.14}
\end{equation*}
$$

Then the infinite volume limit pressure

$$
P_{\sigma}=\lim _{A \uparrow \mathbf{R}^{v}} P_{\sigma, A}
$$

exists for each $0 \leqslant \sigma \leqslant \infty$ as $\Lambda$ tends to $\mathbf{R}^{\nu}$ in the sense of van Hove. Furthermore the limit $P_{\sigma}$ is independent of boundary conditions $\sigma$.

We postpone the proof of the above theorem to Section 5. We believe that one may be able to remove the condition on the decay proporty of potentials in (2.2.14) by improving the method used in Section 5 (the proof of Proposition 5.1).

## 3. BOUNDS ON LOCAL NUMBER OPERATORS

### 3.1. The Wiener Integral Formalism and the Decompositions

In this section and Section 3.1 we prove Proposition 2.2.1 under an assumption of one estimate (Theorem 3.1.1). The methods we will use are the Wiener integral formalism in statistical mechanics and the method similar to that used in Section IV of Ref. 10. We will modify the method used in Ref. 10 to combine the main ideas of Ref. 2 so that the stronger conditions in Ref. 10 are removed.

We start by introducing the Wiener integral formalism. For the details we refer to Refs. 1,6 , and 7 . We will use the following notation:

$$
\begin{align*}
(x)_{n} & =\left(x_{1}, \ldots, x_{n}\right), \quad x_{i} \in \mathbf{R}^{v} \\
d(x)_{n} & =\prod_{i=1}^{n} d^{v} x_{i} \tag{3.1.1}
\end{align*}
$$

The path space of the Wiener measure can be choosen to be

$$
\Omega=C\left([0, \beta], \mathbf{R}^{\nu}\right)
$$

The Wiener measure $P^{\beta}(x, y ; d \omega)$, conditioned on those paths $\omega \in \Omega$ with $\omega(0)=x, \omega(\tau=\beta)=y$, is a $\sigma$-additive, finite measure on $\Omega$. Let $\chi_{B}^{\beta}$ be the characteristic function of the subset $\{\omega \in \Omega: \omega(\tau) \in B$ for all $\tau \in[0, \beta]\}$. We will drop the superscript $\beta$ from $P^{\beta}(x, y ; d \omega)$ and $\chi_{B}^{\beta}$ if there is no confusion involved. We set

$$
\begin{array}{r}
P_{A}(x, y ; d \omega) \equiv \chi_{A}(\omega) P(x, y ; d \omega) \\
P_{A}\left((x)_{n},(y)_{n} ; d(\omega)_{n}\right) \equiv \prod_{j=1}^{n} P_{A}\left(x_{j}, y_{j} ; d \omega_{j}\right) \tag{3.1.2}
\end{array}
$$

Let $\psi_{A}^{\beta}\left((x)_{n},(y)_{n}\right)$ be the kernel of the operator $\exp \left(-\beta H_{A}^{(n)}\right)$ on $L^{2}\left(A^{n}\right)$, where $H_{A}^{(n)}$ is the $n$-particle Hamiltonian given in (2.1.3). By the Feyn-man-Kac formular ${ }^{(1,6,7)}$
$\psi_{A}^{\beta}\left((x)_{n},(y)_{n}\right)=\int_{\Omega^{n}} P_{A}\left((x)_{n},(y)_{n} ; d(\omega)_{n}\right) \exp \left[-\int_{0}^{\beta} U\left((\omega(\tau))_{n}\right) d \tau\right]$

Let $S_{n}$ be the group of the permutations of $\{1,2, \ldots, n\}$ and let $A$ be the multiplication operator by a function $A\left((x)_{n}\right)$ invariant under any $\pi \in S_{n}$.

Then we have

$$
\begin{align*}
& \operatorname{Tr}_{\mathscr{H}_{n}^{(s)}(A)}\left(A \exp \left[-\beta H_{A}^{(n)}\right]\right)=\frac{1}{n!} \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} A\left((x)_{n}\right) \\
& \quad \times \int_{\Omega^{n}} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) \exp \left[-\int_{0}^{\beta} U\left(\left(\omega(\tau)_{n}\right) d \tau\right]\right. \tag{3.1.4}
\end{align*}
$$

where $\pi(x)_{n}=\left(x_{\pi(1)}, \ldots, x_{\pi(n)}\right)$.
We next express the expectation of local number operators in terms of the Wiener integrals. From (2.1.1)-(2.1.4), (2.1.7), and (2.1.8) it follows that

$$
\begin{align*}
\Xi_{A} & =\sum_{n=0}^{\infty} z^{n} \operatorname{Tr}_{\mathscr{H}_{n}^{(s)}(A)}\left(\exp \left[-\beta H_{A}^{(n)}\right]\right)  \tag{3.1.5}\\
\rho_{A}\left(\exp \left(a N_{B}\right)\right) & =\Xi_{A}^{-1} \sum_{n=0}^{\infty} z^{n} \operatorname{Tr}_{\mathscr{H}_{n}^{(s)}(A)}\left(\exp \left(a N_{B}\right) \exp \left[-\beta H_{A}^{(n)}\right]\right)
\end{align*}
$$

where $z=e^{\beta \mu}$ and the term corresponding to $n=0$ equals to 1 . In the rest of this paper we use the following notations;

$$
\begin{align*}
U\left((\omega)_{n}\right) & =\int_{0}^{\beta} U\left((\omega(\tau))_{n}\right) d \tau \\
N_{B}\left((x)_{n}\right) & =\sum_{i=1}^{n} \chi_{B}\left(x_{i}\right) \tag{3.1.6}
\end{align*}
$$

From (3.1.4) it then follows that

$$
\begin{align*}
& \operatorname{Tr}_{\mathscr{H}_{n}^{(s)}(A)}\left(\exp \left[a N_{B}\right] \exp \left[-\beta H_{A}^{(n)}\right]\right)=\frac{1}{n!} \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} \exp \left[a N_{B}\left((x)_{n}\right)\right]  \tag{3.1.7}\\
& \quad \times \int_{\Omega^{n}} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) e^{-U\left((\omega)_{n}\right)}
\end{align*}
$$

By a translation one may assume that $B$ is contained in the ball of diam $(B)$ centered at the origin.

We now introduce a decomposition of the space of $n$ paths into disjoint subsets. To do this we will use the following notations: For any given path configuration $(\omega)_{n}$ and for any given unit cube $Q(r)$ centered at $r \in Z^{v}$, let

$$
\begin{equation*}
n(r, \tau)=\operatorname{card}\left\{\omega_{i}: \omega_{i} \in(\omega)_{n}, \omega_{i}(\tau) \in Q(r)\right\} \tag{3.1.8}
\end{equation*}
$$

That is, $n(r, \tau)$ is the number of $\omega_{i}, i=1,2, \ldots, n$, such that $\omega_{i}(\tau) \in Q(r)$. For given small $\alpha>0$ and for $q \in \mathbf{N}$, let

$$
\begin{align*}
l_{q} & =\left[e^{\alpha q}\right] \\
A_{q} & =\left[-l_{q}-\frac{1}{2}, l_{q}+\frac{1}{2}\right]^{v}  \tag{3.1.9}\\
\left|A_{q}\right| & =\left(2 l_{q}+1\right)^{v}
\end{align*}
$$

where [a] is the integer part of $a \in \mathbf{R}$. For a given $v \in \mathbf{N}$, let us denote

$$
\begin{equation*}
\gamma(v, 2)=1+\frac{2}{v+2} \tag{3.1.10}
\end{equation*}
$$

We now introduce a decomposition of the space of $n$-Wiener paths into a union of multually disjoint subsets

$$
\begin{equation*}
\Omega^{n}=\mathscr{E}_{0} \cup\left(\bigcup_{q \geqslant q_{0}} \mathscr{E}_{q}\right) \tag{3.1.11}
\end{equation*}
$$

as follows: For fixed $1<l<p<\gamma(v, 2)$ which we will choose later and for a fixed large $q_{0} \in Z$, let

$$
\begin{align*}
F_{q}(p)= & \sup _{\tau \in[0, \beta]} \sum_{Q(r) \subset A_{q}} n(r, \tau)^{p}  \tag{3.1.12}\\
E_{q}= & \sum_{Q(r) \in A_{q}} \int_{0}^{\beta} n(r, \tau)^{2} d \tau  \tag{3.1.13}\\
\mathscr{E}_{0}= & \left\{(\omega)_{n} \in \Omega^{n}: F_{q}(p)+E_{q}<\left|\Lambda_{q}\right|^{\prime} \text { for all } q \geqslant q_{0}\right\} \\
\mathscr{E}_{q}= & \left\{(\omega)_{n} \in \Omega^{n}: F_{q}(p)+E_{q} \geqslant\left|\Lambda_{q}\right|^{l} \text { and } F_{q^{\prime}}(p)+E_{q^{\prime}}<\left|A_{q^{\prime}}\right|^{\prime}\right.  \tag{3.1.14}\\
& \left.\quad \text { for all } q_{0} \leqslant q<q^{\prime}\right\}
\end{align*}
$$

The above decomposition is a modification of those used in Refs. 10 and 2 in the following way. In Ref. 10 we have only used the quantity $F_{q}(p)$ to define $\mathscr{E}_{0}$ and $\mathscr{E}_{q}$. On the other hand the authors of Ref. 2 only used the quantity $E_{q}$ for the decomposition. In a sense $F_{q}(p)$ and $E_{q}$ are $L^{\infty}$ norm and $L^{1}$ norm, respectively. As mentioned in the Introduction the method used in Ref. 10 is good enough to control the quantum statistics, but not powerful enough to handle all superstable interactions. Even if the method of Ref. 2 is appropriated to deal with interaction terms, it cannot control the quantum statistics. Thus we conbine two methods as above.

We observe that for $q<q^{\prime}$ and for $\tau \in[0, \beta]$

$$
\begin{align*}
\sum_{Q(r) \subset A_{q^{\prime}}} n(r, \tau) & \leqslant\left|A_{q^{\prime}}\right|\left(\left|A_{q^{\prime}}\right|^{-1} \sum_{Q(r) \subset A_{q^{\prime}}} n(r, \tau)^{p}\right)^{1 / p}  \tag{3.1.15}\\
& \leqslant\left|A_{q^{\prime}}\right|^{1+(l-1) / p} \quad \text { on } \mathscr{E}_{q}
\end{align*}
$$

by Holder's inequality and by the definition of $\mathscr{E}_{q}$ in (3.1.13) and (3.1.14). Also we have that for $q<q^{\prime}$

$$
\begin{equation*}
\sum_{Q(r) \subset A_{q}} \int_{0}^{\beta} n(r, \tau)^{2} d \tau<\left|\Lambda_{q}\right|^{l} \quad \text { on } \quad \mathscr{E}_{q} \tag{3.1.16}
\end{equation*}
$$

The bounds in (3.1.15) and (3.1.16) will be used later repeatly.
Let $\mathscr{E}_{q, k}$ be the subset of $\mathscr{E}_{q}$ in which $k$ paths hit $A_{q+1}$ during the time interval $[0, \beta]$ :
$\mathscr{E}_{q, k}=\left\{(\omega)_{n} \in \mathscr{E}_{q}: \operatorname{card}\left(\left\{\omega_{i}: \omega_{i}(\tau) \in A_{q+1}\right.\right.\right.$ for some $\left.\left.\left.\tau \in[0, \beta]\right\}\right)=k\right\}$
and let $\overline{\mathscr{E}}_{q, k}$ be the subset of $\mathscr{E}_{q}$ in which $\omega_{1}, \ldots, \omega_{k}$ hit $A_{q+1}$ during $[0, \beta]$ :
$\overline{\mathscr{E}}_{q, k}=\left\{(\omega)_{n} \in \mathscr{E}_{q, k}: \omega_{i}(\tau) \in \Lambda_{q+1}\right.$ for some $\left.\tau \in[0, \beta], i=1, \ldots, k\right\}$
Then we have

$$
\begin{equation*}
\mathscr{E}_{q}=\bigcup_{k} \mathscr{E}_{q, k} \tag{3.1.19}
\end{equation*}
$$

From the decompositions in (3.1.12) and (3.1.19) and from (3.1.5) and (3.1.7) it follows that

$$
\begin{equation*}
\rho_{A}\left(\exp \left[a N_{B}\right]\right)=G_{0}+\sum_{n=0}^{\infty} \sum_{q \geqslant q_{0}} \sum_{k} G^{(n)}(q, k) \tag{3.1.20}
\end{equation*}
$$

where $G_{0}$ is the contribution from $\mathscr{E}_{0}$ :

$$
\begin{align*}
G_{0}= & \Xi_{A}^{-1} \sum_{n=0}^{\infty} \frac{z^{n}}{n!} \sum_{\pi \in S_{n}} \int_{\mathcal{A}^{n}} d(x)_{n} \exp \left[a N_{B}\left((x)_{n}\right)\right] \\
& \times \int_{\mathscr{Q}_{0}} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) e^{-U\left((\omega)_{n}\right)} \tag{3.1.21}
\end{align*}
$$

and $G^{(n)}(q, k)$ is the contribution from $\mathscr{E}_{q, k}$;

$$
\begin{align*}
G^{(n)}(q, k)= & \Xi_{\boldsymbol{A}}^{-1} \frac{z^{n}}{n!} \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} \exp \left[a N_{B}\left((x)_{n}\right)\right] \\
& \times \int_{\varepsilon_{q}, k} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) \exp \left[-U\left((\omega)_{n}\right)\right] \tag{3.1.22}
\end{align*}
$$

For a given $B \subset A$ we choose $q_{0}$ sufficiently large so that $B \subset \Lambda_{q_{0}}$. Then by the bound in (3.1.15) we have that for $q_{0} \leqslant q$

$$
\begin{equation*}
N_{B}\left((x)_{n}\right) \leqslant \sum_{Q(r) \in A_{q+1}} n(r, \tau=0) \leqslant\left|A_{q+1}\right|^{1+(i-1) / p} \quad \text { on } \quad \mathscr{E}_{q} \tag{3.1.23}
\end{equation*}
$$

and so

$$
\begin{align*}
G_{0} & \leqslant \exp \left[|a|\left|A_{q_{0}}\right|^{1+(l-1) / p}\right] \\
G^{(n)}(q, k) & \leqslant \Xi_{A}^{-1} z^{n} \exp \left[|a|\left|A_{q+1}\right|^{1+(l-1) / p}\right] \bar{G}^{(n)}(q, k) \tag{3.1.24}
\end{align*}
$$

where

$$
\begin{equation*}
\bar{G}^{(n)}(q, k)=\frac{1}{n!} \int_{A^{n}} d(x)_{n} \int_{\varepsilon_{q, k}} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) e^{-U\left((\omega)_{n}\right)} \tag{3.1.25}
\end{equation*}
$$

We reindex the $k$ paths $\omega_{i}$ 's which hit $\Lambda_{q+1}$ so that $\omega_{i}, i=1,2 .,, . k$ hit $A_{q+1}$. There are $n!/(n-k)!k!$ ways to choose $k$ paths among $n$ paths and so

$$
\begin{align*}
\bar{G}^{(n)}(q, k)= & \frac{1}{(n-k)!k!} \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} \\
& \times \int_{\bar{\sigma}_{q, k}} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) e^{-U\left((\omega)_{n}\right)} \tag{3.1.26}
\end{align*}
$$

Notice that we have $(n-k)$ ! factor instead of $n!$. The total number of terms in the above is card $\left(S_{n}\right)=n!$. This is the main problem arising from quantum statistics.

For any configulation $\left\{\omega_{i}\right\}$ of paths in $\overline{\mathscr{E}}_{q, k}$ we write

$$
\begin{align*}
U_{1}\left((\omega)_{k}\right) & =\int_{0}^{\beta} U_{1}\left((\omega(\tau))_{k}\right) d \tau \\
U_{2}\left((\omega)_{n-k}\right) & =\int_{0}^{\beta} U_{2}\left((\omega(\tau))_{n-k}\right) d \tau  \tag{3.1.27}\\
W\left((\omega)_{k},(\omega)_{n-k}\right) & =\int_{0}^{\beta} W\left((\omega(\tau))_{k},(\omega(\tau))_{n-k}\right) d \tau
\end{align*}
$$

where

$$
\begin{align*}
U_{1}\left((\omega(\tau))_{k}\right) & =\sum_{1 \leqslant i<j \leqslant k} \Phi\left(\omega_{i}(\tau)-\omega_{j}(\tau)\right) \\
U_{2}\left((\omega(\tau))_{n-k}\right) & =\sum_{k+1 \leqslant i<j \leqslant n} \Phi\left(\omega_{i}(\tau)-\omega_{j}(\tau)\right)  \tag{3.1.28}\\
W\left((\omega(\tau))_{k},(\omega(\tau))_{n-k}\right) & =\sum_{i=1}^{k} \sum_{j=k+1}^{n} \Phi\left(\omega_{i}(\tau)-\omega_{j}(\tau)\right)
\end{align*}
$$

We then have

$$
\begin{equation*}
U\left((\omega)_{n}\right)=U_{1}\left((\omega)_{k}\right)+U_{2}\left((\omega)_{n-k}\right)+W\left((\omega)_{k},(\omega)_{n-k}\right) \tag{3.1.29}
\end{equation*}
$$

As in Ref. 10 we define the fluctuations of $k$ paths as follows: we write

$$
\begin{align*}
V(\omega) & =\sup _{\tau_{1}, \tau_{2} \in[0, \beta]}\left|\omega\left(\tau_{1}\right)-\omega\left(\tau_{2}\right)\right|^{2} \\
V\left((\omega)_{k}\right) & =\sum_{i=1}^{k} V\left(\omega_{i}\right) \tag{3.1.30}
\end{align*}
$$

The following is the main technical estimate in proving Prososition 2.2.1:
Theorem 3.1.1. Let the interaction satisfy Assumption A and let $v<4$ if the potential is not repulsive. For a given $b>0$ one can choose $\alpha$ in (4.1.9) small enough, $q_{0}$ large enough, and $1<l<p<\gamma(v, 2)$ such that there exist constant $c>0$ and $\varepsilon>0$ such that

$$
\begin{gathered}
U_{1}\left((\omega)_{k}\right)+b V\left((\omega)_{k}\right)+W\left((\omega)_{k},(\omega)_{n-k}\right) \\
\geqslant c\left(\left|A_{q}\right|^{1+[(l-1)(1+\varepsilon) / p]}+k^{1+\varepsilon}\right)
\end{gathered}
$$

on $\overline{\mathscr{E}}_{q, k}$.
The above theorem is the result corresponding to Theorem IV.2.1 of Ref. 10. We will prove the theorem in Section 4.

### 3.2. The Proof of Proposition 2.2.1; Control of Quantum Statistics

In this section we prove Proposition 2.2.1 under the assumption that Theorem 3.1.1 holds. The proof of Proposition 2.2.1 is almost identical to that in Section IV. 3 of Ref. 10. In Ref. 10 the expression in Theorem 3.1.1 is bounded by $\left|\Lambda_{q+1}\right|^{1+(l-1) \gamma / p}$ for some $\gamma>1$ (see Theorem IV.3.1 of Ref. 10) and also $k$ is bounded by $\left|A_{q+1}\right|^{1+(l-1) / p}$ on $\overline{\mathscr{E}}_{q, k}$. Here we do not have such a bound for $k$. But we have the convergent factor $k^{1+\varepsilon}$ in Theorem 3.1.1. Thus, if one redefine $M_{\gamma}$ in (4.3.26) of Ref. 10 by $M_{\gamma}=k+2\left|A_{q+\alpha^{-1} \log (\gamma+3)}\right|^{1+(1-1) / P}$, every argument used in Section IV. 3 can be applied to our cases here. Therefore in principle the proof of Proposition 2.2 .1 is complete. Since we try to make this paper as self-contained as possible, we will prove Proposition 2.2 .1 in more details by following every necessary steps used in Section IV. 3 of Ref. 10.

From Theorem 3.1.1 and from (3.1.26) and (3.1.29) it follows that

$$
\begin{equation*}
\bar{G}^{(n)}(q, k) \leqslant \exp \left[-c\left(\left|\Lambda_{q}\right|^{1+(l-1)(1+\varepsilon) / p}+k^{1+\varepsilon}\right)\right] \bar{F}^{n)}(q, k) \tag{3.2.1}
\end{equation*}
$$

where

$$
\begin{align*}
\bar{F}^{(n)}(q, k)= & \frac{1}{(n-k)!k!} \sum_{\pi \in s_{n}} \int_{A^{n}} d(x)_{n} \int_{\bar{E}_{q, k}} P_{A}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) \\
& \times \exp \left[b V\left((\omega)_{k}\right)-U_{2}\left((\omega)_{n-k}\right)\right] \tag{3.2.2}
\end{align*}
$$

As we have discussed before, there are $n!$ terms in (3.2.2). Therefore one has to show that many terms in (3.2.2) do not contribute. The main idea is decoupling of the paths in $\left\{\omega_{i}: i=1, \ldots, k\right\}$ from the paths in $\left\{\omega_{i}: i=k+1, \ldots, n\right\}$. On the subset $\overline{\mathscr{E}}_{q, k}$ a path $\omega_{i}, i \leqslant k$, may join to a path $\omega_{j}, k<j$, to form a composite path, i.e., $x_{\pi(i)}=x_{j}$. If the fluctuation of $\omega_{i}$ is small, there will not be many ways to form a composite tragectory. On the other hand, if the fluctuation is large, the corresponding Wiener measure will be small. Thus we further decompose the subset $\bar{E}_{q, k}$ into a union of mutually disjoint subsets corresponding to large and small fluctuations.

We adapt the notations used in Section III. 2 and Section IV. 3 of Ref. 10 . For a unit cube $\Delta$ and for $l \geqslant 0$, we write

$$
\Lambda(A, l)=\left\{x \in \mathbf{R}^{v}: \max _{1 \leqslant i \leqslant v} \inf _{y \in \Delta}\left|x^{i}-y^{i}\right| \leqslant l\right\}
$$

That is, $\Lambda(\Delta, l)$ is the cube with its volume $(2 l+1)^{v}$ containing $\Delta$ at the center. We denote

$$
\begin{gather*}
\mathscr{E}_{\Delta, l}=\left\{\omega \in \Omega ; \omega(\tau) \in \Delta, \omega\left(\tau^{\prime}\right) \in A(\Delta, l+1) \backslash A(\Delta, l) \text { for some } \tau, \tau^{\prime} \in[0, \beta]\right. \\
\text { and } \left.\omega\left(\tau^{\prime \prime}\right) \notin \Lambda(\Delta, l+1)^{c} \text { for all } \tau^{\prime \prime} \in[0, \beta]\right\}  \tag{3.2.3}\\
\mathscr{E}_{4,-1}=\{\omega \in \Omega: \omega(\tau) \in \Delta \text { for all } \tau \in[0, \beta]\}
\end{gather*}
$$

where $A^{c}$ is the complement of $A$. Note that any path in $\mathscr{E}_{\Delta, l} l \geqslant 0$, must visit $\Delta$ and $\Lambda(\Delta, l+1) \backslash \Lambda(\Delta, l)$, but not $\Lambda(\Delta, l+1)^{c}$. We also have that for any positive function $F$ on $\Omega$

$$
\begin{align*}
& \int P_{A}(x, y ; d \omega)\left[1-\chi_{A_{q+1}^{c}}(\omega)\right] F(\omega) \\
& \quad \leqslant \sum_{\Delta \subset A_{q+1}} \sum_{l=-1}^{\infty} \int_{\mathscr{E}_{A, l}} P(x, y ; d \omega) F(\omega) \tag{3.2.4}
\end{align*}
$$

where $\sum_{A \subset \Lambda_{q+1}}$ is the summation over unit cube $A=Q(r) \subset \Lambda_{q+1}$. That is, in order to hit $\Lambda_{q+1}$ the path must hit one of $Q(r) \subset \Lambda_{q+1}$. For the detailed derivation of (3.2.4) we refer to the derivation of (4.3.7) in Ref. 10. Let $\mathscr{E}_{d, t}^{(i)}$ be the subset $\mathscr{E}_{4, l}$ for the $i$ path, $i=1, \ldots, k$. We write

$$
\begin{align*}
& \overline{\mathscr{E}}(q, \\
& \left.\quad k ;\left(\Delta_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right)  \tag{3.2.5}\\
& \quad=\left\{(\omega)_{n} \in \overline{\mathscr{E}}_{q, k} ; \omega_{i} \in \mathscr{E}_{\Delta_{i}, l_{i}}^{(i)}, i=1, \ldots, k\right\}
\end{align*}
$$

Since the fluctuation is bounded by $V(\omega) \leqslant[\operatorname{diag}(\Lambda(A, l+1))]^{2} \leqslant$ $4 v(l+2)^{2}$ on $\mathscr{E}_{\Delta, l}$, we have that on $\overline{\mathscr{E}}\left(q, k ;\left(\Delta_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right)$

$$
\begin{equation*}
V\left((\omega)_{k}\right) \leqslant \sum_{i=1}^{k}\left(l_{i}+2\right)^{2} 4 v \tag{3.2.6}
\end{equation*}
$$

We remember that $\omega_{i}, i \leqslant k$, must visit $A_{q+1}$ at least one time. Thus we have

$$
\begin{align*}
& \int_{\bar{\delta}_{q, k}} P_{A}\left((x)_{n},(y)_{n} ; d(\omega)_{n}\right) \cdots \\
& \quad=\int_{\bar{\delta}_{q, k}} P_{A}\left((x)_{n},(y)_{n} ; d(\omega)_{n}\right) \prod_{i=1}^{k}\left[1-\chi_{A_{q+1}}^{c}\left(\omega_{i}\right)\right] \cdots \tag{3.2.7}
\end{align*}
$$

We apply (3.2.7), (3.2.4), and (3.2.6) (in that order) to (3.2.2) to obtain the bound

$$
\begin{align*}
\bar{F}^{(n)}(q, k) \leqslant & \frac{1}{(n-k)!k!} \sum_{\Delta_{1} \subset A_{q+1}} \cdots \sum_{\Delta_{k} \subset \Lambda_{q+1}} \sum_{l_{1}=-1}^{\infty} \cdots \sum_{l_{k}=-1}^{\infty} \\
& \times \exp \left[4 v b \sum_{i=1}^{b}\left(l_{i}+2\right)^{2}\right] \bar{F}^{(n)}\left(q, k ;\left(\Delta_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right) \tag{3.2.8}
\end{align*}
$$

where

$$
\begin{align*}
\bar{F}^{(n)}\left(q, k ;\left(\Delta_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right)= & \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} \mathbf{K}\left((x)_{n}, \pi(x)_{n}\right) \\
\mathbf{K}\left((x)_{n},(y)_{n}\right) \equiv & \int_{\tilde{\varepsilon}\left(q, k ;\left(A_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right)}  \tag{3.2.9}\\
& \times P_{A}\left((x)_{n},(y)_{n} ; d(\omega)_{n}\right) e^{-U_{2}\left((\omega)_{n-k}\right)}
\end{align*}
$$

We recall that for any configuration of $n$ path $(\omega)_{n}$ in $\overline{\mathscr{E}}\left(q, k ;\left(\Delta_{1}, l_{1}\right), \ldots,\left(\Lambda_{k}, l_{k}\right)\right)$ the path $\omega_{i}, i=1, \ldots, k$, must stay inside of $A\left(\Delta_{i}, l_{i}+1\right)$ by (3.2.3) and (3.2.5). This means that $\omega_{i}, i=1, \ldots, k$, cannot form a composite path with any path $\omega_{l}$ if $\omega_{l}(\tau=0) \in \Lambda\left(A_{i}, l_{i}+1\right)^{c}$. By the bound (3.1.15) there are at most $\left|A_{q+\alpha^{-1} \log \left(l_{i}+3\right)}\right|^{1+(l-1) / p}$ number of paths which can hit $\Lambda\left(\Delta, l_{i}+1\right) \subset A_{q+\alpha^{-1}} \log \left(l_{i}+3\right)$. Thus if $l_{i}, i \leqslant k$, are small, many terms in (3.2.9) will be vanished. If $l_{i}^{\prime}$ 's are large, the corresponding Wiener integral will be small (see Proposition 3.2 .3 in the following). This is the idea of controlling quantum statistics.

We will use the following abbreviated notations:

$$
\begin{align*}
M\left(l_{i}\right) & =\left|A_{q+\alpha^{-1} \log \left(l_{i}+3\right)}\right|^{1+[l-1) / p]} \\
M_{l_{i}} & =\left[2 M\left(l_{i}\right)\right]+k  \tag{3.2.10}\\
\gamma & =\max \left\{l_{1}, \ldots, l_{k}\right\}
\end{align*}
$$

We remark that the definition of $M_{y}$ differ to that in (4.3.26) of Ref. 10 by the factor $k$. We will use the fact that $M_{\gamma}-k>M(\gamma)$ as in Ref. 10. We write

$$
\begin{align*}
S(k, n)= & \left\{\pi \in S_{n}: \pi \text { is a permutation of }\{k+1, \ldots, n\}\right\} \\
E\left(m, m^{\prime}\right)= & \left\{\pi \in S_{n}: \pi \text { is an interchange of } m\right. \text { with one of } \\
& \left.m, m+1, \ldots, m^{\prime} \text { for } m \leqslant m^{\prime} \leqslant n\right\} \\
E\left(m, m^{\prime}\right)= & E(m, n) \text { if } m^{\prime}>n  \tag{3.2.11}\\
E_{k}\left(l_{1}, \ldots, l_{k}\right)= & \left\{p=p_{k} p_{k-1} \cdots p_{1} ; p_{i} \in E\left(i, M_{l_{i}}\right)\right\} \\
\tilde{E}_{k}\left(l_{1}, \ldots, l_{k}\right)= & \left\{p=p_{1} p_{2} \cdots p_{k} ; p_{i} \in E\left(i, M_{l_{i}}\right)\right\} \\
\widetilde{E}_{n, k}= & \left\{p=p_{1} p_{2} \cdots p_{k} ; p_{i} \in E(i, n) \text { for } i=1, \ldots, k\right\}
\end{align*}
$$

Here we have modified the definition of $E_{k}\left(l_{1}, \ldots, l_{k}\right)$ in (4.3.14) of Ref. 10 by replacing $E\left(i,\left[2 M\left(l_{i}\right)\right]\right)$ by $E\left(i, M_{l_{i}}\right)$. We note that

$$
S_{n}=\tilde{E}_{n, k} S(k, n)
$$

and

$$
\begin{align*}
\operatorname{card}\left(\widetilde{E}_{k}\left(l_{1}, \ldots, l_{k}\right)\right) & =\operatorname{card}\left(E_{k}\left(l_{1}, \ldots, l_{k}\right)\right) \\
& \leqslant \prod_{i=1}^{k} M_{l_{i}} \tag{3.2.12}
\end{align*}
$$

For any $\sigma \in S(k, n)$ we use the notation
$\chi_{\sigma}:$ the characteristic function of the subset

$$
\begin{equation*}
\left\{(x)_{n} \in \mathbf{R}^{v n}:\left\|x_{\sigma(k+1)}\right\| \leqslant\left\|x_{\sigma(k+2)}\right\| \leqslant \cdots \leqslant\left\|x_{\sigma(n)}\right\|\right\} \tag{3.2.13}
\end{equation*}
$$

where $\|x\|==\max _{1 \leqslant i \leqslant \nu}\left|x^{i}\right|$. We then have

$$
\begin{equation*}
\sum_{\sigma \in S(k, n)} \chi_{\sigma}\left((x)_{n}\right)=\sum_{\sigma \in S(k, n)} \chi_{e}\left(\sigma(x)_{n}\right)=1 \tag{3.2.14}
\end{equation*}
$$

where $e \in S(k, n)$ is the identity element. The following is the result corresponding to Lemma 4.3.1 of Ref. 10.

Lemma 3.2.1. Let $\mathbf{K}\left((x)_{n},(y)_{n}\right)$ be defined as in (3.2.9). Then

$$
\begin{aligned}
& \bar{F}^{(n)}\left(q, k ;\left(A_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right) \\
& \quad \leqslant k!\sum_{p \in \mathbb{E}_{k}\left(l_{1}, \ldots, l_{k}\right)} \sum_{\sigma, \sigma^{\prime} \in S(k, n)} \int_{A_{n}} d(x)_{n} \chi_{e}\left((x)_{n}\right) \mathbf{K}\left(\sigma^{\prime}(x)_{n}, p \sigma(x)_{n}\right)
\end{aligned}
$$

Proof. The lemma follows from the method used in the proof of Lemma 4.3.1 of Ref. $10, \Lambda(\Delta, l) \subset A_{q+\alpha^{-1} \log (l+3)}$ for $\Delta \subset A_{q+1}$, and from the fact that $M_{l_{i}}-k>M\left(l_{i}\right)$. For the details, see the proof of Lemma 4.3.1 of Ref. 10.

We now briefly review the partial symmetry spaces introduced in Section 3.2 of Ref. 10. Let $A_{n}$ be a subset of $S_{n}$. The subset $A_{n}$ is not necessary subgroup of $S_{n}$. Let $\left\{f_{i}\right\}$ be an orthonormal base for $L^{2}\left(A, d^{v} x\right)$ and let $\mathscr{H}_{A_{n}}^{(s)}(A)$ be the closed subspace of $L^{2}\left(A^{n}\right)$ spanned by the following vectors:

$$
\operatorname{card}\left(A_{n}\right)^{-1} \sum_{\pi \in A_{n}} f_{i_{1}}\left(x_{\pi(1)}\right) f_{i_{2}}\left(x_{\pi(2)}\right) \cdots f_{i_{n}}\left(x_{\pi(n)}\right)
$$

Let $P\left(A_{n}\right)$ be the projection operator onto $\mathscr{H}_{A_{n}}^{(s)}(A)$. For any $f \in L^{2}\left(\Lambda^{n}\right)$ we have

$$
\left(P\left(A_{n}\right) f\right)(x)_{n}=\operatorname{card}\left(A_{n}\right)^{-1} \sum_{\pi \in A_{n}} f\left(\pi(x)_{n}\right)
$$

If $A$ be an operator of the trace class which admits an integral kernel $A\left((x)_{n},(y)_{n}\right)$, it follows that

$$
\begin{align*}
& \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(P\left(A_{n}\right) A P\left(A_{n}\right)\right) \\
& \quad \leqslant \operatorname{card}\left(A_{n}\right)^{-2} \sum_{\pi, \pi^{\prime} \in A_{n}} \int_{A^{n}} d(x)_{n} A\left(\pi^{-1}(x)_{n}, \pi^{\prime-1}(x)_{n}\right) \tag{3.2.15}
\end{align*}
$$

For more detaily we refer to Ref. 10.
Let $P(S(k, n))$ and $P\left(E_{k}\left(l_{1}, \ldots, l_{k}\right)\right)$ be the projection operators onto $\mathscr{H}_{(S(k, n)}^{(s)}(A)$ and $\mathscr{H}_{E_{k k}\left(l_{1}, \ldots, l_{k}\right)}^{(A)}(A)$, respectively. For any $\sigma \in S(k, n)$, let $\chi_{\sigma}$ be the (projection) operator defined by

$$
\left(\chi_{\sigma} f\right)(x)_{n}=\chi_{e}\left(\sigma(x)_{n}\right) f\left((x)_{n}\right)
$$

for any $f \in L^{2}\left(\mathbf{R}^{v n}\right)$. From (3.2.13) we have

$$
\begin{equation*}
\sum_{\sigma \in S(k, n)} \chi_{\sigma}=1 \tag{3.2.16}
\end{equation*}
$$

In order to avoid notational complications we use the following abbreviated notations:

$$
\begin{align*}
P_{1} & =P(S(k, n)) \\
P_{2} & =P\left(E_{k}\left(l_{1}, \ldots, l_{k}\right)\right)  \tag{3.2.17}\\
C_{1,2} & =\operatorname{card}(S(k, n))^{2} \operatorname{card}\left(E_{k}\left(l_{1}, \ldots, l_{k}\right)\right)
\end{align*}
$$

Let $\mathbf{K}$ be the operator on $L^{2}\left(A^{n}\right)$ defined by its kernel $\mathbf{K}\left((x)_{n},(y)_{n}\right)$, where $\mathbf{K}\left((x)_{n},(y)_{n}\right)$ has been defined in (3.2.9). Using (3.2.15) the following is easy to derive:

$$
\begin{align*}
\operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} \mathbf{K} P_{1} P_{2} \chi_{e}\right)= & C_{1,2}^{-1} \sum_{p \in \widetilde{E}_{k}\left(l_{1}, \ldots, l_{k}\right)} \sum_{\sigma, \sigma^{\prime} \in S(k, n)} \\
& \times \int_{\mathcal{A}^{n}} d(x)_{n} \chi_{e}\left((x)_{n}\right) \mathbf{K}\left(\sigma^{\prime}(x)_{n}, p \sigma(x)_{n}\right) \tag{3.2.18}
\end{align*}
$$

and so by Lemma 3.2.1 we have

$$
\begin{equation*}
\bar{F}^{(n)}\left(q, k ;\left(\Delta_{1} . l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right) \leqslant k!C_{1,2} \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} \mathbf{K} P_{1} P_{2} \chi_{e}\right) \tag{3.2.19}
\end{equation*}
$$

We use the identity (3.2.16) to obtain

$$
\begin{equation*}
\operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} \mathbf{K} P_{1} P_{2} \chi_{e}\right)=\sum_{\sigma^{\prime \prime} \in S(k, n)} \operatorname{Tr}_{\dot{L}^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} \mathbf{K} P_{1} \chi_{\sigma^{\prime \prime}} P_{2} \chi_{e}\right) \tag{3.2.20}
\end{equation*}
$$

Let $\gamma$ and $M_{\gamma}$ be defined as in (3.2.10), and let

$$
\begin{equation*}
A(k, n) \equiv\left\{\sigma^{\prime \prime} \in S(k, n): \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} \mathbf{K} P_{1} \chi_{\sigma^{\prime \prime}} P_{2} \chi_{e}\right) \neq 0\right\} \tag{3.2.21}
\end{equation*}
$$

Then the following result corresponds to Lemma IV.3.2 of Ref. 10:

## Lemma 3.2.2:

$$
\operatorname{card}(A(k, n)) \leqslant \operatorname{card}\left(\tilde{E}\left(l_{1}, \ldots, l_{k}\right)\right) M_{\gamma}!/\left(M_{\gamma}-k\right)!
$$

Proof. The lemma follows from the method same as that used in Lemma IV.3.2 of Ref. 10 and from the fact that $M_{\gamma}-k>M(\gamma)$. For the details we refer to Ref. 10.

Using the definition of $A(k, n)$ and (3.2.19), (3.2.20) we obtain

$$
\begin{align*}
\bar{F}^{(n)}\left(q, k ;\left(\Delta_{1}, l_{1}, \ldots,\left(\Delta_{k}, l_{k}\right)\right) \leqslant\right. & k!C_{1,2} \sum_{\sigma^{\prime \prime} \in A(k, n)} \operatorname{Tr}_{L^{2}\left(A^{n}\right)} \\
& \times\left(\chi_{e} P_{1} \mathbf{K} P_{1} \chi_{\sigma^{\prime \prime}} P_{2} \chi_{e}\right) \tag{3.2.22}
\end{align*}
$$

The next step is to get a bound for (3.2.22). To do this, we define

$$
\begin{equation*}
\mathbf{K}_{\Delta, l}(x, y)=\int_{\mathscr{E}(A, l)} P(x, y ; d \omega) \tag{3.2.23}
\end{equation*}
$$

Let $\mathbf{K}_{A, l}$ be the operator corresponding to the kernel $\mathbf{K}_{\Delta, l}(x, y)$. Let $P_{A(\Delta, l)}$ be the projection operator onto $L^{2}(A(\Delta, l)) \subset L^{2}\left(\mathbf{R}^{v}\right)$, and let

$$
\begin{equation*}
\widetilde{\mathbf{K}}_{\Delta, l}=P_{A(\Delta, l+1)} \exp \left(\frac{\beta}{4} \Delta\right) P_{A(A, l+1)} \tag{3.2.24}
\end{equation*}
$$

We then have the following result (Proposition III.2.2 of Ref. 10):
Lemma 3.2.3. For $l=-1,0,1, \ldots$ there are constants $c_{1}>0$ and $c_{2}>0$ such that
(a) $\mathbf{K}_{A, l}(x, y) \leqslant c_{1} \exp \left(-\frac{l^{2}}{16 \beta}\right) \overline{\mathbf{K}}_{\Delta, l}(x, y)$
(b) $\operatorname{Tr}_{L^{2}\left(\mathbf{R}^{v}\right)}\left(\overline{\mathbf{K}}_{A, l}\right) \leqslant c_{2}|\Lambda(A, l+1)|$
where $\overline{\mathbf{K}}_{\Delta, l}(x, y)$ is the kernel of $\overline{\mathbf{K}}_{\Delta, l}$ and $|A(\Delta, l+1)|$ is the volume of $A(\Delta, l+1)$.

For the proof of the above lemma we refer to Ref. 10, the above result implies that for large $l$ (for large fluctuation) the contribution of the Wiener measure is small.

We now start to get a bound for (3.2.22). From (3.2.5) the following is obvious:

$$
\begin{equation*}
\overline{\mathscr{E}}\left(q, k ;\left(A_{1}, l_{1}\right), \ldots,\left(A_{k}, l_{k}\right)\right) \subset\left({\underset{i=1}{k}}_{\left.\mathscr{E}_{\Delta_{i}, l_{i}}^{(i)}\right) \times \Omega^{n-k}, ~}^{\text {in }}\right. \tag{3.2.25}
\end{equation*}
$$

Using (3.2.15) [and the method to get (3.2.20)] it is easy to sheck that ${ }^{(10)}$

$$
\begin{align*}
& \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} \mathbf{K} P_{1} \chi_{\sigma^{\prime}} P_{2} \chi_{e}\right) \leqslant C_{1,2}^{-1} \sum_{p \in \widetilde{E}_{k}\left(l_{1}, \ldots, i_{k}\right)} \sum_{\sigma, \sigma \in S(k, n)} \\
& \quad \times \int_{A^{n}} d(x)_{n} \chi_{e}\left((x)_{n}\right) \mathbf{K}\left(\sigma^{\prime}(x)_{n}, p \sigma(x)_{n}\right) \chi_{e}\left(p \sigma^{\prime \prime}(x)_{n}\right) \tag{3.2.26}
\end{align*}
$$

From the definition of $\mathbf{K}\left((x)_{n},(y)_{n}\right)$ in (3.2.9) and from (3.2.25) it follows that

$$
\begin{align*}
\mathbf{K}\left((x)_{n},(y)_{n}\right) \leqslant & \int_{\left(\mathrm{X}_{i=1}^{k} \mathscr{E}_{\delta_{b} i_{i}}^{(i)}\right)} P\left((x)_{k},(y)_{k} ; d(\omega)_{k}\right) \\
& \times\left(\exp \left[-\beta H_{A}^{n-k}\right]\right)\left((x)_{n-k},(y)_{n-k}\right) \tag{3.2.27}
\end{align*}
$$

We use Lemma 3.2.3(a) to obtain

$$
\begin{align*}
\mathbf{K}\left((x)_{n},(y)_{n}\right) \leqslant & \prod_{i=1}^{k}\left[c_{1} \exp \left(-\frac{l_{i}^{2}}{16 \beta}\right)\right] Y\left((x)_{n},(y)_{n}\right)  \tag{3.2.28}\\
Y\left((x)_{n},(y)_{n}\right) \equiv & \left(\prod_{i=1}^{k} \overline{\mathbf{K}}_{a_{i, l}, x_{i}}\left(x_{i}, y_{i}\right)\right) \\
& \times\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right)\left((x)_{n-k},(y)_{n-k}\right) \tag{3.2.29}
\end{align*}
$$

Let $Y$ be the operator corresponding to the kernel $Y\left((x)_{n},(y)_{n}\right)$ :

$$
\begin{equation*}
Y \equiv\left(\prod_{i=1}^{k} \widetilde{\mathbf{K}}_{\Delta_{i, L_{i}}}\right) \exp \left[-\beta H_{A}^{(n-k)}\right] \tag{3.2.30}
\end{equation*}
$$

which is positive operator on $L^{2}\left(\mathbf{R}^{v n}\right)$. From (3.2.22), (3.2.26), and (3.2.28) it follows that

$$
\begin{align*}
& \bar{F}^{(n)}\left(q, k ;\left(A_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right) \\
& \leqslant \\
& k!C_{1,2}\left(\prod_{i=1}^{n} c_{1} \exp \left(-\frac{l_{i}^{2}}{16 \beta}\right)\right)  \tag{3.2.31}\\
& \quad \times \sum_{\sigma^{\prime \prime} \in A(k, n)} \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} Y P_{1} \chi_{\sigma^{\prime \prime}} P_{2}\right)
\end{align*}
$$

We now use the positivity of $Y$ and the abstract Hölder's inequality ${ }^{(12)}$ to conclude that

$$
\begin{align*}
& \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} Y P_{1} \chi_{\sigma^{\prime \prime}} P_{2}\right) \\
& \leqslant\left[\operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} Y P_{1}\right)\right]^{1 / 2}\left[\operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{\sigma^{\prime}} P_{1} Y P_{1}\right)\right]^{1 / 2} \tag{3.2.32}
\end{align*}
$$

Here we have used the fact that $\operatorname{Tr}(P A P) \leqslant \operatorname{Tr}(A)$ for any projection $P$ and positive $A$. A direct computation gives us that for any $\sigma^{\prime \prime} \in(S(k, n))$ :

$$
\begin{equation*}
\operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{\sigma^{\prime \prime}} P_{1} Y P_{1}\right)=[(n-k)!]^{-1} \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(P_{1} Y P_{1}\right) \tag{3.2.33}
\end{equation*}
$$

See the below of (4.3.36) of Ref. 10 for the detailed derivation of the above equality. We substitute (3.2.33) into (3.2.32) and then we use the definitions of $P_{1}$ and $Y$ in (3.2.17) and (3.2.30), respectively, to obtain

$$
\begin{align*}
& \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(\chi_{e} P_{1} Y P_{1} \chi_{\sigma^{\prime \prime}} P_{2}\right) \\
& \quad \leqslant \frac{1}{(n-k)!} \operatorname{Tr}_{L^{2}\left(A^{n}\right)}\left(P_{1} Y P_{1}\right) \\
& \quad=\frac{1}{(n-k)!}\left(\prod_{i=1}^{k} \operatorname{Tr}_{L^{2}(A)}\left(\overline{\mathbf{K}}_{A_{i, l_{i}}}\right)\right) \operatorname{Tr}_{\mathscr{H}_{n-k}^{(s)}(A)}\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right) \\
& \quad \leqslant \frac{1}{(n-k)!} c_{2}^{k} \prod_{i=1}^{k}\left(2 l_{i}+2\right)^{v} \operatorname{Tr}_{\mathscr{H}_{\substack{(s)}}^{(s)}(A)}\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right) \tag{3.2.34}
\end{align*}
$$

for any $\sigma^{\prime \prime} \in S(k, n)$. Here we have used Lemma 3.2.3(b) to get the third inequality. We now combine (3.2.31) and (3.2.34) to obtain

$$
\begin{align*}
& \bar{F}^{(n)}\left(q, k ;\left(\Delta_{1}, l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right) \\
& \quad \leqslant M\left(q, k ; l_{1}, \ldots, l_{k}\right) \operatorname{Tr}_{\mathscr{H}_{n-k}^{(s)}(\Lambda)}\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right) \tag{3.2.35}
\end{align*}
$$

where

$$
\begin{align*}
& M\left(q, k ; l_{1}, \ldots, l_{k}\right) \\
& =\operatorname{card}(A(k, n)) k!C_{1,2} \frac{1}{(n-k)!} c_{2}^{k} \prod_{i=1}^{k}\left(2 l_{i}+2\right)^{v} \\
& \quad \times \exp \left(-\frac{l_{i}^{2}}{16 \beta}\right) \tag{3.2.36}
\end{align*}
$$

From the bound in Lemma 3.2.2

$$
\begin{equation*}
\operatorname{card}(A(k, n)) \leqslant \operatorname{card}\left(\widetilde{E}_{k}\left(l_{1}, \ldots, l_{k}\right) M_{\gamma}^{k}\right) \tag{3.2.37}
\end{equation*}
$$

and from the definition of $M_{\gamma}$ in (3.2.10) it is easy to check that there is a constant $c$ such that

$$
\begin{equation*}
M_{\gamma}^{k} \leqslant \exp \left(c k\left[\log k+\log \left|\Lambda_{q+1}\right|+\log (\gamma+3)\right]\right) \tag{3.2.38}
\end{equation*}
$$

From (3.2.12) and (3.2.10) we also have that for some $c>0$

$$
\begin{align*}
& \operatorname{card}\left(\tilde{E}_{k}\left(l_{1}, \ldots, l_{k}\right)\right) \\
& \quad \leqslant \exp \left[c k\left(\log k+\log \left|\Lambda_{q+1}\right|\right] \prod_{i=1}^{k} \exp \left[c\left(l_{i}+2\right)\right]\right. \tag{3.2.39}
\end{align*}
$$

We combine the definition of $C_{1,2}$ in (3.2.17), and (3.2.37)-(3.2.39) to
(3.2.36) and we use the fact that $k \log \gamma \leqslant k \log k+\gamma \log \gamma$ for $\gamma>1$ to conclude that there are constants $c_{1}$ and $c_{2}$ such that

$$
\begin{align*}
\left.M(q, k) ; l_{1}, \ldots, l_{k}\right) \leqslant & c_{1}^{k}(n-k)!\left[\prod_{i=1}^{k} \exp \left(-\frac{l_{i}^{2}}{32 \beta}\right)\right] \\
& \times \exp \left[c_{2} k\left(\log k+\log \left|\Lambda_{q+1}\right|\right)\right] \tag{3.2.40}
\end{align*}
$$

We substitute (3.2.40) into (3.2.35) and (3.2.35) into (3.2.8) to obtain the following bound:

$$
\begin{align*}
\bar{F}^{(n)}(q, & k) \\
\leqslant & c_{1}^{k} \exp \left[c_{2} k\left(\log k+\log \left|A_{q+1}\right|\right)\right] \\
& \times \operatorname{Tr}_{\mathscr{H}_{n-k}^{(s)}(-1)}\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right) \\
& \times \prod_{i=1}^{k}\left[\sum_{L_{i} \in \Lambda_{q+1}} \sum_{i_{i}=-1}^{\infty} \exp \left[-\frac{1}{32 \beta} l_{i}^{2}+4 v b\left(l_{i}+2\right)^{2}\right]\right] \\
\leqslant & d_{i}^{k} \exp \left[d_{2}\left(k \log k+\left|A_{q+1}\right| \log \left|A_{q+1}\right|\right] \operatorname{Tr}_{\mathscr{H}_{n-k}^{(s)}-k}\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right)\right. \tag{3.2.41}
\end{align*}
$$

Here we have choose $b$ such that $4 v b<1 / 32 \beta$ and we have used the fact that $k(\log k+\log \gamma) \leqslant \operatorname{const}(k \log k+\gamma \log \gamma)$ to get the second inequality. Next we note that for any $\varepsilon>0$ and for any $c_{1}, c_{2}, c_{3}, c>0$ there is a constant $d_{1}>0$ such that
$c_{1} k+c_{2}\left(k \log k+\left|A_{q+1}\right| \log \left|A_{q+1}\right|\right)-c\left(k^{1+\varepsilon}+\left|A_{q}\right|^{1+\varepsilon}\right) \leqslant d_{1}$
We finally combine (3.2.41) and (3.2.1) and then use (3.2.42) to obtain

$$
\begin{aligned}
\bar{G}^{(n)}(q, k) \leqslant & \exp \left[-c^{\prime}\left|A_{q+1}\right|^{1+[(l-1)(1+\varepsilon) / P]}-c^{\prime} k^{1+\varepsilon}+d_{1}\right] \\
& \times \operatorname{Tr}_{\mathscr{\varkappa}_{n-k}^{(s)}(A)}\left(\exp \left[-\beta H_{A}^{(n-k)}\right]\right.
\end{aligned}
$$

for some constant $c^{\prime}, d_{1}>0$. We substitute the above bound to (3.1.24) to conclude that

$$
\sum_{n=0}^{\infty} \sum_{q \geqslant 40} \sum_{k} G^{(n)}(q, k) \leqslant e^{D} \Xi_{A}^{-1} \sum_{n^{\prime}=0}^{\infty} z^{n^{\prime}} \operatorname{Tr}_{\mathscr{H}_{n^{\prime}}^{(s)}}\left(\exp \left[-\beta H_{A}^{\left(n^{\prime}\right)}\right]\right)=e^{D}
$$

for some $D>0$. Proposition 2.2.1 follows from (3.1.20), (3.1.24), and the above bound. This complete the proof of Proposition 2.2.1.

## 4. CONTROL OF INTERACTIONS: PROOF OF THEOREM 3.1.1

In this section we complete the proof of Proposition 2.2 .1 by proving Theorem 3.1.1. We start with some technical estimates. Let $F_{q}(p)$ and $\gamma(v, 2)$ be defined as in (3.1.12) and (3.1.10), respectively. The following result corresponds to Proposition IV.4.1 and Proposition IV.4.3 of Ref. 10.

Lemma 4.1. For a given path configuration $(\omega)_{n} \in \overline{\mathscr{E}}_{q}, k$, let us assume that

$$
d_{1} F_{q}(p) \geqslant\left|A_{q}\right|^{\prime} \quad \text { for some } \quad d_{1}>0
$$

and let $1<l<p<\gamma(\nu, 2)$. Then for any given $a_{1}>0$ and $b_{1}>0$ there exist constants $c_{1}>0$ and $c_{2}>0$ such that

$$
a_{1} U_{1}\left((\omega)_{k}\right)+b_{1} V\left((\omega)_{k}\right) \geqslant c_{1}\left|A_{q+1}\right|^{1+(l-1) \gamma(v, 2) / p}-c_{2} k
$$

on $\overline{\mathscr{E}_{q, k}}$.
Proof. By the assumption in the lemma there exists $\tau_{0} \in[0, \beta]$ such that

$$
d_{1} \sum_{Q(\gamma) \in A_{q}} n\left(\gamma, \tau_{0}\right)^{p} \geqslant\left|A_{q}\right|^{\prime}
$$

We fix $\tau_{0} \in[0, \beta]$ so that the above inequality holds. Let

$$
\begin{equation*}
B(\tau) \equiv a_{1} U_{1}\left((\omega(\tau))_{k}\right)+b_{1} V\left((\omega(\tau))_{k}\right) \tag{4.1}
\end{equation*}
$$

where

$$
\begin{equation*}
V\left((\omega(\tau))_{k}\right)=\frac{1}{\beta} \sum_{i=1}^{k}\left|\omega_{i}(\tau)-\omega_{i}\left(\tau_{0}\right)\right|^{2} \tag{4.2}
\end{equation*}
$$

and $U_{1}\left((\omega(\tau))_{k}\right)$ has been defined in (3.1.28). We denote that $\tilde{n}(r, \tau)$ is the number of $\omega_{i}, i=1, \ldots, k$ such that $\omega_{i}(\tau) \in Q(r), \bar{\omega}_{l}$ be the path hitting $\Lambda_{q}$ at $\tau=\tau_{0}$ [i.e., $\left.\bar{\omega}_{( }\left(\tau_{0}\right) \in A_{q}\right]$, and $\bar{n}(r, \tau)$ be the number of $\bar{\omega}_{l}$ 's such that $\bar{\omega}_{l}(\tau) \in Q(r)$. By the definition we have

$$
\begin{equation*}
d_{1} \sum_{Q(r) \subset A_{q}} \bar{n}\left(r, \tau_{0}\right)^{p} \geqslant\left|A_{q}\right|^{2} \tag{4.3}
\end{equation*}
$$

and by the superstability condition [Assumption $\mathrm{A}(\mathrm{a})$ ]

$$
\begin{align*}
B(\tau) \geqslant & a_{1} \sum_{Q(r) \subset A}\left[A \tilde{n}(r, \tau)^{2}-B \tilde{n}(r, \tau)\right]+b_{1} V\left((\omega)_{k}\right) \\
\geqslant & \sum_{Q(r) \subset \Lambda}\left[a_{1} \tilde{n}(r, \tau)^{2}+\frac{b_{1}}{\beta}\right. \\
& \left.\times \sum_{\bar{\omega}_{\hat{\prime}}(\tau) \in Q(r)}\left|\bar{\omega}_{i}(\tau)-\bar{\omega}_{i}\left(\tau_{0}\right)\right|^{2}\right]-a_{1} B k \tag{4.4}
\end{align*}
$$

We now employ the method used in the proof of Proposition IV.4.1 of Ref. 10. Let $\bar{n}_{r^{\prime}}(r, \tau)$ be the number of $\bar{\omega}_{i}$ 's such that $\bar{\omega}_{i}(\tau) \in Q(r)$ and $\bar{\omega}_{i}\left(\tau_{0}\right) \in Q\left(r^{\prime}\right)$. Then

$$
\bar{n}(r, \tau)^{2} \geqslant \sum_{Q\left(r^{\prime}\right) \subset A_{q}} \bar{n}_{r^{\prime}}(r, \tau)^{2}
$$

and so

$$
\begin{aligned}
B(\tau) \geqslant & \sum_{Q\left(r^{\prime}\right) \in A_{q}}\left[a_{1} A \sum_{Q(r) \subset A} \bar{n}_{r^{\prime}}(r, \tau)^{2}\right. \\
& \left.+\frac{b_{1}}{\beta} \sum_{\substack{\bar{\omega}_{i}(\tau) \in Q(r) ; \\
\bar{\omega}_{i}\left(\tau_{0}\right) \in Q\left(r^{\prime}\right)}}\left|\bar{\omega}_{i}(\tau)-\bar{\omega}_{i}\left(\tau_{0}\right)\right|^{2}\right]-a_{1} B k
\end{aligned}
$$

We have arrived at the expression which is the same as that in (4.4.12) of Ref. 10.

Consider the case in which for a given $q>0$

$$
\begin{aligned}
\operatorname{card}\left(\left\{\bar{\omega}_{i}: \bar{\omega}_{i}\left(\tau_{0}\right) \in Q\left(r^{\prime}\right) \text { and }\left|\bar{\omega}_{i}(\tau)-\bar{\omega}_{i}\left(\tau_{0}\right)\right|\right.\right. & \left.\left.\leqslant \bar{n}\left(r^{\prime}, \tau_{0}\right)^{q}\right\}\right) \\
& \geqslant\left[\frac{1}{2} \bar{n}\left(r^{\prime}, \tau_{0}\right)\right]
\end{aligned}
$$

Since the number of $Q(r)^{\prime}$ 's such that $\left|r-r^{\prime}\right| \leqslant \bar{n}\left(r^{\prime}, \tau_{0}\right)^{q}$ is not larger than $\left[2 \bar{n}\left(r^{\prime}, \tau_{0}\right)^{q}+1\right]^{\nu}$ we use Hölder's inequality to obtain

$$
\sum_{Q(r) \subset 1} \bar{n}_{r}(r, \tau)^{2} \geqslant c_{1}^{\prime} \bar{n}\left(r^{\prime}, \tau_{0}\right)^{v q} \bar{n}\left(r^{\prime}, \tau_{0}\right)^{(1-v q)^{2}}-c_{2}^{\prime}
$$

for some constant $c_{1}^{\prime}$ and $c_{2}^{\prime}$. On the other hand if

$$
\begin{aligned}
\operatorname{card}\left(\left\{\bar{\omega}_{i}: \bar{\omega}_{i}\left(\tau_{0}\right) \in Q\left(r^{\prime}\right) \text { and }\left|\bar{\omega}_{i}(\tau)-\bar{\omega}_{i}\left(\tau_{0}\right)\right|\right.\right. & \left.\left.\leqslant \bar{n}\left(r^{\prime}, \tau_{0}\right)^{q}\right\}\right) \\
& <\left[\frac{1}{2} \bar{n}\left(r^{\prime}, \tau_{0}\right)\right]
\end{aligned}
$$

then

$$
\sum_{\bar{\omega}_{i}(\tau) \in Q(r): \bar{w}_{i}\left(\tau_{0}\right) \in Q\left(r^{\prime}\right)}\left|\bar{\omega}(\tau)-\bar{\omega}_{i}\left(\tau_{0}\right)\right|^{2} \geqslant c_{1}^{\prime} \bar{n}\left(r^{\prime}, \tau_{0}\right)^{1+2 q}-c_{2}^{\prime}
$$

for some constant $c_{1}^{\prime}$ and $c_{2}^{\prime}$. Choosing $q$ such that $v q+(1-v q) 2=1+2 q$ (i.e., $q=1 /[2 v+(2-v)]$ ) we obtain

$$
\begin{equation*}
B(\tau) \geqslant c \sum_{Q(r) \subset A_{q}} \bar{n}\left(r, \tau_{0}\right)^{\gamma(v, 2)}-a_{1} B k-c_{2}\left|A_{q}\right| \tag{4.5}
\end{equation*}
$$

For the detailed derivation of (4.5) from (4.4) we refer to the proof of Proposition IV.4.1 of Ref. 10.

We note that

$$
\begin{equation*}
a_{1} U_{1}\left((\omega)_{k}\right)+b_{1} V\left((\omega)_{k}\right) \geqslant \int_{0}^{\beta} B(\tau) d \tau \tag{4.6}
\end{equation*}
$$

We use Hölder's inequality to obtain

$$
\left[\left|A_{q}\right|^{-1} \sum_{Q(r) \subset A_{q}} \bar{n}\left(r, \tau_{0}\right)^{\gamma}\right]^{p / \gamma} \geqslant\left|A_{q}\right|^{-1} \sum_{Q(r) \subset A_{q}} \bar{n}\left(r, \tau_{0}\right)^{p}
$$

and so by (4.3)

$$
\begin{equation*}
\sum_{Q(r) \subset A_{q}} \bar{n}\left(r, \tau_{0}\right)^{\gamma} \geqslant\left|A_{q}\right|^{1+(l-1) \gamma / p} \tag{4.7}
\end{equation*}
$$

The lemma now follows from (4.5), (4.6), and (4.7).
Proposition 4.2. For any $a_{1}>0$ and $b_{1}>0$ one can choose $q_{0}$ large enough such that for $1<l<p<\gamma(v, 2)$ there exist constant $c_{1}>0$ and $c_{2}>0$ such that

$$
a_{1} U\left((\omega)_{k}\right)+b_{1} V\left((\omega)_{k}\right) \geqslant c_{1}\left|A_{q+1}\right|^{1+(l-1)(1+\varepsilon) / p}-c_{2} k
$$

on $\overline{\mathscr{E}}_{q, k}$.
Proof. If $\frac{1}{2} F_{q}(p) \geqslant\left|A_{q}\right|^{l}$, the proposition follows from Lemma 4.1 and the fact that $1<\gamma(v, 2)$. On the other hand if $\frac{1}{2} F_{q}(p)<\left|\Lambda_{q}\right|^{l}$ on $\overline{\mathscr{E}}_{q, k} \subset \mathscr{E}_{q}$, we have that

$$
E_{q}=\sum_{Q(r) \subset A_{a}} \int_{0}^{\beta} n(r, \tau)^{2} d \tau \geqslant \frac{1}{2}\left|A_{q}\right|^{l}
$$

by the definition of $\mathscr{E}_{q}$ in (3.1.14). Thus by Assumption $\mathrm{A}(\mathrm{a})$

$$
\begin{align*}
a_{1} U\left((\omega)_{k}\right) & \geqslant a_{1} \sum_{Q(r) \subset A} \int_{0}^{\beta} d \tau\left[A n(r, \tau)^{2}-B n(r, \tau)\right] \\
& \geqslant c_{1}\left|A_{q}\right|^{l}-c_{2} k \tag{4.8}
\end{align*}
$$

on $\mathscr{E}_{q, k}$ for some constant $c_{1}$ and $c_{2}$. We choose $l$ and $p$ such that $1<l<p<\gamma(v, 2)$. Then $1+[(l-1) / p]<l$ and so the proposition follows from (4.8) and and the definition of $\Lambda_{q}$ in (3.1.9). This completes the proof of the proposition.

Lemma 4.3. For any $a_{2}>$ and $b_{2}>0$ there exist $c_{3}>0$ and $\varepsilon>0$ such that for sufficiently large $q_{0}, 1<l<p<\gamma(v, 2)$, and for sufficiently small $\alpha$

$$
a_{2} U_{1}\left((\omega)_{k}\right)+b_{2} V\left((\omega)_{k}\right) \geqslant c_{3} k^{1+\varepsilon}
$$

on $\overline{\mathscr{E}}_{q, k}$.
Proof. If $k \leqslant c\left|\Lambda_{q}\right|^{1+(l-1) / p}$ for some constant, the lemma follows from Proposition 4.2. Thus we next consider the case in which

$$
k \geqslant c\left|A_{q}\right|^{1+(l-1) / p}
$$

for some $c>0$. Let

$$
k^{\prime}=\operatorname{card}\left(\left\{\omega_{i}: \omega_{i}(\tau) \in A_{q+\log k} \text { for all } \tau \in[0, \beta], i=1, \ldots, k\right\}\right.
$$

We consider the case in which $\frac{1}{2} k \leqslant k^{\prime}$. Let $\tilde{n}(r, \tau)$ be the number of $\omega_{i} \in\left\{\omega_{1}, \ldots, \omega_{k}\right\}$ such that $\omega_{i}(\tau) \in Q(r)$. Then by the superstability

$$
\begin{align*}
a_{2} U_{1}\left((\omega)_{k}\right) & \geqslant a_{2} \int_{0}^{\beta} \sum_{Q(r) \in A}\left[A \tilde{n}(r, \tau)^{2}-B \tilde{n}(r, \tau)\right] d \tau  \tag{4.9}\\
& \geqslant a_{2} \int_{0}^{\beta}\left[A \sum_{Q(r) \in \Lambda_{q+\log k}} \tilde{n}(r, \tau)^{2}\right] d \tau-c^{\prime} k
\end{align*}
$$

for some constant $c^{\prime}$. We use Hölder's inequality to obtain

$$
\begin{aligned}
\sum_{Q(r) \in \Lambda_{q+\log k}} \tilde{n}(r, \tau)^{2} & \geqslant\left|A_{q+\log k}\right|\left[\left|A_{q+\log k}\right|^{-1} \sum_{Q(r) \in \Lambda_{q+\log k}} \tilde{n}(r, \tau)\right]^{2} \\
& \geqslant \frac{1}{4} k^{2}\left|\Lambda_{q+\log k}\right|^{-1}
\end{aligned}
$$

if $k^{\prime} \geqslant \frac{1}{2} k$. Thus from (4.9) and the definition of $A_{q}$ it follows that

$$
\begin{equation*}
a_{2} U_{1}\left((\omega)_{k}\right) \geqslant c^{\prime \prime}\left|A_{q}\right|^{-1} k^{-\alpha \nu} k^{2}-c^{\prime} k \tag{4,10}
\end{equation*}
$$

For small $\alpha$, the above is bigger than $k^{1+\varepsilon}$ for some $\varepsilon>0$. Next we consider the case in which $k^{\prime}<\frac{1}{2} k$. This means that more than $\frac{1}{2} k$ paths hit $A_{q+1}$ and $\Lambda_{q+\log k}^{c}$ during the time interval $[0, \beta]$. Thus

$$
\begin{align*}
b_{2} V\left((\omega)_{k}\right) & \geqslant c \frac{1}{2} k\left[l_{q+\log k}-l_{q+1}\right]^{2} \\
& \geqslant c^{\prime} k\left[e^{\alpha \log (k-1)}-1\right] l_{q+1}^{2}  \tag{4.11}\\
& \geqslant c^{\prime \prime} k^{l+\varepsilon}
\end{align*}
$$

Thus in either case we have the bound

$$
a_{2} U_{1}\left((\omega)_{k}\right)+b_{2} V\left((\omega)_{k}\right) \geqslant c k^{1+\varepsilon}
$$

This proves the lemma.
Corollary 4.4. For any $a_{1}>0$ and $b_{1}>0$ one can choose $\alpha$ small enough, $q_{0}$ large enough such that for $1<l<p<\gamma(v, 2)$ there exist constant $c_{1}>0$ and $\varepsilon>0$ such that

$$
a_{1} U_{1}\left((\omega)_{k}\right)+b_{1} V\left((\omega)_{k}\right) \geqslant c_{1}\left[\left|A_{q+1}\right|^{1+(l-1)(1+\varepsilon) / p}+k^{1+\varepsilon}\right]
$$

on $\overline{\mathscr{E}}_{q, k}$
Proof. The above is a consequence of Proposition 4.2 and Lemma 4.3.

Proposition 4.5. For any $b_{2}>0$, one can choose $\alpha$ small enough and $q_{0}$ large enough such that for $1<l<p<\gamma(v, 2)$

$$
U_{1}\left((\omega)_{k}\right)+b_{2} V\left((\omega)_{k}\right) \geqslant \frac{9}{10} A\left|A_{q}\right|^{l}
$$

on $\overline{\mathscr{E}}_{4, k}$, where $A$ is the constant in Assumption $\mathrm{A}(\mathrm{a})$.
Proof. If $20 F_{q}(p) \geqslant\left|A_{q}\right|^{l}$, we choose $p$ such that $l<1+$ $(l-1) \gamma(v, 2) / p(1<\gamma(v, 2) / p)$. We then apply Lemma 4.1 and Corollary 4.4 to get the bound in the proposition. On the orther hand, if $20 F_{q}(p)<\left|\Lambda_{q}\right|^{\prime}$, then $E_{q} \geqslant \frac{19}{20}\left|A_{q}\right|^{l}$. Then it follows that

$$
U_{1}\left((\omega)_{k}\right) \geqslant \frac{19}{20} A\left|A_{q}\right|^{I}-B \beta k
$$

and so the bound in the proposition follows from the above bound and Corollary 4.4. This proves the proposition.

We next control the $W\left((\omega)_{k},(\omega)_{n-k}\right)$ term. We will employ the method similar to that used in Ref. 2. For a given path configuration $(\omega)_{n}$ in $\overline{\mathscr{E}}_{q, k}$ we divide $(\omega)_{k}=\left\{\omega_{1}, w_{2}, \ldots, \omega_{k}\right\}$ into the following two classes:

$$
\begin{equation*}
(\omega)_{k}=(\zeta) \cup(\xi) \tag{4.12}
\end{equation*}
$$

where
( ) : the paths completely contained in $A_{4+3}$
$(\xi)$ : the paths which cross the boundary $\partial A_{q+3}$ of $A_{q+3}$
The foilowing results correspond to Lemma 2.2 and Lemma 2.3 of Ref. 2, respectively.

Lemma 4.6. If ( $\tilde{\omega})$ and $(\hat{\omega})$ are sets of paths of $(\omega)_{n}$ contained in
$\Lambda_{q+1+a}(a \geqslant 0)$ and $\Lambda_{q+1}^{c}$, respectively, then there exist $\alpha$ small enough and $q_{0}$ large enough such that for $q_{0} \leqslant q$

$$
W((\tilde{\omega}),(\hat{\omega})) \geqslant-\frac{1}{16} A\left|A_{q+1}\right|^{l}
$$

on $\mathscr{E}_{q}$, where $A$ is the constant in Assumption A(a).
Lemma 4.7. Let ( $\zeta$ ) be the path contained in $A_{s}, s>q+1$. Then there is a $q_{0}$ large enough such that for each $q \geqslant q_{0}$

$$
W\left((\zeta),(\omega)_{n-k}\right) \geqslant-d_{1}|\zeta|\left|A_{s}\right|^{1 / 2}
$$

on $\mathscr{E}_{q}$ for some constants $d_{1} d_{2}>0$, where $|\zeta|$ is $\operatorname{card}((\zeta))$.
We leave the proofs of Lemma 4.6 and Lemma 4.7 to the end of this section. Using the above lemmas we now prove Theorem 3.1.1.

Proof of Theorem 3.1.1. From Corollary 4.4, Proposition 4.5, and Lemma 4.6 it follows that

$$
\begin{align*}
& U\left((\omega)_{k}\right)+b V\left((\omega)_{k}\right)+W\left((\omega)_{k},(\omega)_{n-k}\right) \\
&= \frac{9}{10} U_{1}\left((\omega)_{k}\right)+\frac{1}{10} U_{1}\left((\omega)_{k}\right)+\frac{b}{2} V\left((\omega)_{k}\right)+W\left((\zeta),(\omega)_{n-k}\right) \\
&+\frac{b}{2} V\left((\omega)_{k}\right)+W\left((\xi),(\omega)_{n-k}\right) \\
& \geqslant\left(\frac{9}{10}\right)^{2} A\left|A_{q}\right|^{l}+c^{\prime}\left[\left|A_{q+1}\right|^{1+(l-1)(1+\varepsilon) / p}+k^{1+\varepsilon}\right]-\frac{1}{16} A\left|A_{q+1}\right|^{l} \\
&+\frac{b}{2} V\left((\omega)_{k}\right)+W\left((\xi),(\omega)_{n-k}\right) \tag{4.14}
\end{align*}
$$

Note that, if the interaction is repulsive, $W\left((\xi),(\omega)_{n-k}\right) \geqslant 0$. In this case the theorem follows from the above inequality. We consider the general case. Let $\left(\xi^{(s)}\right) \subset(\omega)_{k}$ be the paths which are contained in $A_{s}$ and cross the boundary $\partial \Lambda_{s-1}$ of $\Lambda_{s-1}$. Obviously we have that by Lemma 4.7

$$
\begin{aligned}
& \frac{b}{2} V\left((\omega)_{k}\right)+W\left((\xi),(\omega)_{n-k}\right) \\
& \quad \geqslant \sum_{s=q+4}^{\infty}\left[\frac{b}{2} V\left(\left(\xi^{(s)}\right)\right)+W\left(\left(\xi^{(s)}\right),(\omega)_{n-k}\right)\right] \\
& \quad \geqslant \sum_{s=q+4}^{\infty}\left[\frac{b}{2}\left|\xi^{(s)}\right|\left(l_{s}-l_{q+1}\right)^{2}-d\left|\xi^{(s)}\right|\left|A_{s}\right|^{\mid / 2}\right] \\
& \quad \geqslant \sum_{s=q+4}^{\infty}\left|\xi^{(s)}\right|\left[\frac{b}{2}\left(l_{s}-l_{q+1}\right)^{2}-3^{v} d l_{s}^{v / 2}\right]
\end{aligned}
$$

For $v \leqslant 3$ we choose $l>1$ such that $\frac{1}{2} v l<2$ [and also $1<\gamma(v, 2) / p$ so that Proposition 4.5 holds $]$. Notice that $l_{s}=\left[e^{\alpha s}\right]$. The expressions in the bracket in the last inequality become positive for sufficiently large $q_{0}\left(q_{0} \leqslant q\right)$. Thus for sufficiently large $q_{0}$ we have

$$
\begin{equation*}
\frac{b}{2} V\left((\omega)_{k}\right)+W\left((\xi),(\omega)_{n-k}\right) \geqslant 0 \tag{4.15}
\end{equation*}
$$

Thus for small $\alpha$ the theorem follows from (4.14) and (4.15). This completes the proof of Theorem 3.1.1.

In the rest of this section we prove Lemma 4.6 and Lemma 4.7. We first prove Lemma 4.7.

Proof of Lemma 4.7. We use the method similar to that used in the proof of Lemma 2.3 of Ref. 2. By the lower regularity in Assumption A(b)

$$
\begin{aligned}
& -W\left((\zeta(\tau)),(\omega(\tau))_{n-k}\right) \\
& \quad \leqslant \sum_{Q(r) \subset A_{s}} \sum_{Q\left(r^{\prime}\right) \subset A_{q+1}^{c}} \psi\left(r, r^{\prime}\right) n(r,(\zeta(\tau))) n\left(r^{\prime},(\omega(\tau))\right)
\end{aligned}
$$

where $n(r,(\zeta(\tau)))$ is the number of the paths $\zeta_{i} \in(\zeta)$ with $\zeta_{i}(\tau) \in Q(r)$, and $n\left(r^{\prime},(\omega(\tau))\right)$ is the number of the paths $\omega_{l} \in(\omega)_{n-k}$ with $\omega_{l}(\tau) \in Q\left(r^{\prime}\right)$. We use Hölder's inequality to obtain

$$
\begin{align*}
-W & \left((\zeta(\tau)),(\omega(\tau))_{n-k}\right) \\
& \leqslant \sum_{k=0}^{\infty} \psi(k) \sum_{Q(r) \in A_{s}} n(r,(\zeta(\tau))) \sum_{\substack{Q\left(r^{\prime}\right) \subset A_{q+1}^{c} \\
d\left(r, r^{\prime}\right)=k}} n\left(r^{\prime},(\omega(\tau))\right)  \tag{4.16}\\
\leqslant & v_{1} \sum_{k=0}^{\infty} \psi(k)(k+1)^{(v-1) / 2} \sum_{Q(r) \subset A_{s}} n(r,(\zeta(\tau))) \\
& \times\left[\sum_{\substack{Q\left(r^{\prime}\right) \subset A_{q}^{c} \\
d\left(r, r^{\prime}\right)=k}} n\left(r^{\prime},(\omega(\tau))\right)^{2}\right]^{1 / 2}
\end{align*}
$$

for some constant $v_{1}$ depending on $v$. But

$$
\begin{equation*}
\sup _{Q(r) \in A_{s}} \sum_{\substack{Q\left(r^{\prime}\right) \subset \Lambda_{q+1}^{c} \\ d\left(r, r^{\prime}\right)=k}} n\left(r^{\prime},(\omega(\tau))\right)^{2} \leqslant \sum_{Q\left(r^{\prime}\right) \in \Lambda_{s}+r(k)} n\left(r^{\prime},(\omega(\tau))\right)^{2} \tag{4.17}
\end{equation*}
$$

where $r(k)$ is the smallest of integers $r$ such that the set $\left\{Q\left(r^{\prime}\right): d\left(r, r^{\prime}\right)=k\right.$ for all $\left.Q(r) \subset \Lambda_{s}\right\}$ is contained in $\Lambda_{s+r}$. Thus we substitute (4.17) into (4.16). By integrating

$$
\begin{equation*}
-W\left((\zeta),(\omega)_{n-k}\right) \leqslant v_{1}|\zeta| \sum_{k=0}^{\infty} \psi(k)(k+1)^{(\nu-1) / 2}\left|\Lambda_{s+r(k)}\right|^{l / 2} \tag{4.18}
\end{equation*}
$$

on $\mathscr{E}_{q}$. By Lemma B. 2 of Ref. 2

$$
\begin{equation*}
r(k) \leqslant 1+\frac{1}{\alpha} \log (k+2) \tag{4.19}
\end{equation*}
$$

We also note that by (4.19) and the definition of $\Lambda_{q}$

$$
\begin{align*}
\left(\frac{\left|\Lambda_{s+r(k)}\right|}{\left|\Lambda_{s}\right|}\right)^{l / 2} & \leqslant \exp [\alpha v l r(k)]  \tag{4.20}\\
& \leqslant c(k+1)^{v / 2}
\end{align*}
$$

for some constant depending on $v$. From (4.18) and (4.20) it follows that

$$
-W\left((\zeta),(\omega)_{n-k}\right) \leqslant c^{\prime}|\zeta|\left|\Lambda_{s}\right|^{l / 2} \sum_{k=0}^{\infty} \psi(k)(k+1)^{(\nu l+\nu-1) / 2}
$$

If we choose $l=1+\delta$ for sufficiently small $\delta$, then lemma follows from the lower regularity (2.2.3) and the above bound.

We finally prove Lemma 4.6 by using the method similar to that used in the proof of Lemma 2.2. of Ref. 2, and so we complete the proof of Theorem 3.1.1.

Proof of Lemma 4.6. We note that by Assumption A(b)
$-W((\tilde{\omega}(\tau)),(\hat{\omega}(\tau)))$

$$
\begin{align*}
& \leqslant \frac{1}{2} \sum_{Q(r) \subset \Lambda_{q+1+\alpha}} \sum_{Q\left(r^{\prime}\right) \subset \Lambda_{q+1}^{c}} \psi\left(r, r^{\prime}\right)\left[n(r,(\omega(\tau)))^{2}+n\left(r^{\prime},(\omega(\tau))\right)^{2}\right]  \tag{4.21}\\
& =T_{1}+T_{2}+T_{3}+T_{4}+T_{5}
\end{align*}
$$

where

$$
\begin{align*}
& T_{1}=\frac{1}{2} \sum_{Q(r) \subset A_{q+1+a \backslash A_{q}} Q\left(r^{\prime}\right) \subset A_{q+2+a \backslash A_{q+1}}} \psi\left(r, r^{\prime}\right) n(r)^{2} \\
& T_{2}=\frac{1}{2} \sum_{Q(r) \subset A_{q} Q\left(r^{\prime}\right) \subset A_{q+2+a \backslash \Lambda_{q+1}}} \psi\left(r, r^{\prime}\right) n(r)^{2} \\
& T_{3}=\frac{1}{2} \sum_{Q(r) \subset A_{q+1+a}} \sum_{Q\left(r^{\prime}\right)=A_{q+2+a}^{c}} \psi\left(r, r^{\prime}\right) n(r)^{2}  \tag{4.22}\\
& T_{4}=\frac{1}{2} \sum_{Q(r) \subset A_{q+1+a}} \sum_{Q\left(r^{\prime}\right)=A_{q+2+a \backslash A_{q+1}}} \psi\left(r, r^{\prime}\right) n\left(r^{\prime}\right)^{2} \\
& T_{5}=\frac{1}{2} \sum_{Q(r) \leq A_{q+1+a}} \sum_{Q\left(r^{\prime}\right) \subset \Lambda_{q+2+a}^{c}} \psi\left(r, r^{\prime}\right) n\left(r^{\prime}\right)^{2}
\end{align*}
$$

Here $n(r)=n(r,(\tilde{\omega}(\tau)))$ and $n\left(r^{\prime}\right)=n\left(r^{\prime},(\hat{\omega}(\tau))\right)$. Let

$$
\begin{align*}
F & \equiv \sum_{k=0}^{\infty} \psi(k)(k+1)^{v-1}<\infty \\
F(l) & \equiv \sum_{k=1}^{\infty} \psi(k)(k+1)^{v-1}  \tag{4.23}\\
\psi_{k} & \equiv \sum_{k^{\prime}=I_{q+1+a+k}-I_{q+1+a}} \psi\left(k^{\prime}\right)
\end{align*}
$$

Then the following are easy to obtained:

$$
\begin{align*}
& T_{1} \leqslant \frac{1}{2} F \sum_{Q(r) \subset \Lambda_{q+1+a \backslash \Lambda_{q}}} n(r)^{2} \\
& T_{2} \leqslant \frac{1}{2} F\left(l_{q+1}-l_{q}\right) \sum_{Q(r) \subset A_{q}} n(r)^{2} \\
& T_{3} \leqslant \frac{1}{2} F\left(l_{q+2+a}-l_{q+1+a}\right) \sum_{Q(r) \subset \lambda_{q}+1+a} n(r)^{2}  \tag{4.24}\\
& T_{4} \leqslant \frac{1}{2} F \sum_{Q\left(r^{\prime}\right) \subset A_{q+2+a \backslash \Lambda_{q+1}}} n\left(r^{\prime}\right)^{2} \\
& T_{5} \leqslant \frac{1}{2} \sum_{Q(r) \subset A_{q+1+a}} \sum_{k=1}^{\infty} \psi_{k}\left[\sum_{Q\left(r^{\prime}\right) \subset A_{q+2+a+k \backslash A_{q+1+a+k}}} n\left(r^{\prime}\right)^{2}\right]
\end{align*}
$$

By integrating with respect to $\tau$ and by noting that $l_{q+1}-l_{q} \leqslant$ $l_{q+a+2}-l_{q+a+1}$, we obtain that

$$
\begin{align*}
-W((\tilde{\omega}),(\hat{\omega})) \leqslant & F E_{\Lambda_{q+2+a \backslash A_{q}}+F\left(l_{q+1}-l_{q}\right) E_{A_{q+1+a}}} \\
& +\frac{1}{2} \sum_{Q(r) \subset \Lambda_{q+1+a}} \sum_{k=1}^{\infty} \psi_{k}\left[E_{A_{q+2+a+k \backslash A_{q+a+2}}}-E_{A_{q+1+a+k} \backslash A_{q+2+a}}\right] \\
\leqslant & F E_{\Lambda_{q+2+a \backslash A_{q}}+F\left(l_{q+1}-l_{q}\right) E_{A+1+a}} \\
& +\frac{1}{2} \sum_{Q(r) \in \Lambda_{q+1}} \sum_{k=1}^{\infty}\left(\psi_{k}-\psi_{k+1}\right) E_{\Lambda_{q+2+k+a}} \tag{4.25}
\end{align*}
$$

Here for $B \subset \mathbf{R}^{v}$ we have written

$$
\begin{equation*}
E_{B} \equiv \sum_{Q(r) \subset B} \int_{0}^{\beta} n(r, \tau)^{2} d \tau \tag{4.26}
\end{equation*}
$$

We assert that, if $q_{0}$ is sufficiently large $\left(q_{0} \leqslant q\right)$ and if $l=1+\delta$ with sufficiently small $\delta>0$, then the following bounds hold:

$$
\begin{array}{r}
{\left[\left|\Lambda_{q+2+a}\right|^{l}-\left|\Lambda_{q}\right|^{l}\right] /\left|\Lambda_{q+1}\right|^{l}<A / 48 F} \\
{\left[\left|\Lambda_{q+1+a}\right|^{l} F\left(l_{q+1}-l_{q}\right)\right] /\left|\Lambda_{q+1}\right|^{l}<A / 48}  \tag{4.27}\\
\left|\Lambda_{q+1+a}\right|\left|\Lambda_{q+1}\right|^{-l} \sum_{k=1}^{\infty}\left(\psi_{k}-\psi_{k+1}\right)\left|\Lambda_{q+2+a+k}\right|^{l}<A / 24
\end{array}
$$

for $\alpha$ small enough. We prove our assertions. By the definition of $\Lambda_{q}$ in (4.1.9) the left expression in the first inequality is bounded by $\exp [\alpha v(a+1)]-1$, and so the first inequality holds for sufficiently small $\alpha$. We note that $l_{q+1}-l_{q}>\alpha l_{q}$ and so $F\left(l_{q+1}-l_{q}\right) \rightarrow 0$ as $q \rightarrow \infty$. This proves the second inequality. To prove the last inequality in (4.27) we note that

$$
\begin{aligned}
& \left|A_{q+2+a+k}\right|^{l} /\left[2\left(l_{q+2+a+k}-l_{q+2+a}\right)+3\right]^{v l} \\
& \quad \leqslant\left[l_{q+a+2+k} /\left(l_{q+1+a+k}-l_{q+a+1}\right)\right]^{v l} \\
& \quad \leqslant\left[e^{\alpha} /\left(1-e^{-\alpha k}\right)\right]^{v l} \\
& \quad \leqslant\left[e^{\alpha} /\left(1-e^{-\alpha}\right)\right]^{v l}
\end{aligned}
$$

Thus the left side of the third inequality in (4.29) is bounded by

$$
\begin{aligned}
& {\left[e^{\alpha} /\left(1-e^{-\alpha}\right)\right]^{v l} e^{\alpha v a} \sum_{k=1}^{\infty}\left(\psi_{k}-\psi_{k+1}\right)\left[2\left(l_{q+2+a+k}-l_{q+a+1}\right)\right]^{v l}} \\
& \quad \leqslant\left[e^{\alpha} /\left(1-e^{-\alpha}\right)\right]^{v l} e^{\alpha v a} \sum_{s=I_{q+2+a-l_{q+1+a}}}[\psi(s)-\psi(s+1)][2(s+3)]^{v l} \\
& \quad \rightarrow 0 \text { as } q \rightarrow \infty
\end{aligned}
$$

by the lower regurality and the fact that $l_{q+2+a}-l_{q+1+a} \rightarrow \infty$ as $q \rightarrow \infty$. This proves the third inequality in (4.27).

We note that for $q_{0} \leqslant q<q^{\prime}$ that on $\mathscr{E}_{q}$

$$
\begin{align*}
E_{A_{q^{\prime}}} & <\left|A_{q^{\prime}}\right|^{l} \\
E_{A_{q^{\prime}} \backslash A_{q}} & <E_{A_{q^{\prime}}}-E_{A_{q}}  \tag{4.28}\\
& \leqslant\left(E_{A_{q^{\prime}}}-E_{\Lambda_{q}}\right)+\left[F_{q^{\prime}}(p)-F_{q}(p)\right] \\
& <\left|\Lambda_{q^{\prime}}\right|^{l}-\left|\Lambda_{q}\right|^{l}
\end{align*}
$$

where for $B \subset \mathbf{R}^{\nu}$

$$
F_{B}(p) \equiv \sup _{\tau \in[0, \beta]} \sum_{Q(r) \subset B} n(r, \tau)^{p}
$$

Thus the lemma follows from (4.25), (4.27), and (4.28). This completes the proof of Lemma 4.6.

## 5. PRESSURE AND EQUIVALENCE OF THE BOUNDARY CONDITIONS: PROOF OF THEOREM 2.2 .4

We will prove Theorem 2.2.4 in the following manner. From (2.2.12) it follows that

$$
\begin{equation*}
P_{\Lambda}=P_{\infty, A} \leqslant P_{\sigma, A} \leqslant P_{0, A} \tag{5.1}
\end{equation*}
$$

Thus, if one can prove that the limit

$$
\begin{equation*}
P=\lim _{A \uparrow \mathbf{R}^{v}} P_{A} \tag{5.2}
\end{equation*}
$$

exists as $\Lambda$ tends to $\mathbf{R}^{v}$ in the sense of van Hove, and that

$$
\begin{equation*}
P_{0, \Lambda}-P_{\Lambda} \rightarrow 0 \quad \text { as } \quad \Lambda \rightarrow \mathbf{R}^{v} \tag{5.3}
\end{equation*}
$$

one proves the existence of the infinite volume limit pressure and its independence of the boundary conditions.

We begin to prove Theorem 2.2 .4 by showing the existence of the pressure with Dirichlet boundary conditions.

Proposition 5.1. Let $A$ tend to $\mathbf{R}^{v}$ in the sense of van Hove. Under the assumptions in Theorem 2.2.4 the limit

$$
P=\lim _{A \uparrow \mathbf{R}^{v}} P_{A}
$$

exists.
Proof. Let us consider the following rectangular regions:

$$
\begin{align*}
A & =\left(-l_{1}, l_{1}\right) \times S \\
A_{1} & =\left(-l_{1},-\frac{1}{2} b\right) \times S \\
A_{2} & =\left(\frac{1}{2} b, l_{1}\right) \times S  \tag{5.4}\\
S & =\prod_{i=2}^{v}\left(-l_{i}, l_{i}\right)
\end{align*}
$$

where $b>0$ is the constant given in Theorem 2.2.4. We will show that

$$
\begin{equation*}
\exp \left[-m_{0}|S|\right] \Xi_{A_{1}} \Xi_{A_{2}} \leqslant \Xi_{A} \tag{5.5}
\end{equation*}
$$

for some constant $m_{0} \geqslant 0$. The proposition then follows for translationinvariant interactions by a standard argument in statistical mechanics. ${ }^{(1,13)}$

We prove (5.5). Since $A_{1} \cup A_{2} \subset A$, it follows that

$$
\begin{equation*}
\Xi_{\Lambda} \geqslant \operatorname{Tr}_{\mathrm{F}^{(s)}\left(\Lambda_{1} \cup \Lambda_{2}\right)}\left\{\exp \left[-\beta\left(H_{\Lambda_{1}}+H_{\Lambda_{2}}+W\left(\Lambda_{1}, \Lambda_{2}\right)\right)\right]\right\} \tag{5.6}
\end{equation*}
$$

where $W\left(\Lambda_{1}, \Lambda_{2}\right)$ is the operator corresponding to the interaction between the particles in $\Lambda_{1}$ and the particles in $\Lambda_{2}$. The decay property in Theorem 2.2.4 yields

$$
\begin{equation*}
W\left(\Lambda_{1}, A_{2}\right) \leqslant \frac{1}{2} \sum_{\substack{Q\left(r^{\prime}\right)=\Lambda_{1} \\ Q\left(r^{\prime}\right) \in A_{2}}} \tilde{\varphi}\left(\left|r-r^{\prime}\right|\right)\left[N_{Q(r)}^{2}+N_{Q\left(r^{\prime}\right)}^{2}\right] \tag{5.7}
\end{equation*}
$$

We use the Peierls-Bogoliubov-type inequality [see (2.16) of Ref. 9]

$$
\begin{equation*}
\operatorname{Tr}\left(e^{A+B}\right) \geqslant \operatorname{Tr}\left(e^{A}\right) \exp \left\{\operatorname{Tr}\left(B e^{A}\right) / \operatorname{Tr}\left(e^{A}\right)\right\} \tag{5.8}
\end{equation*}
$$

for trace class operators $e^{A}$ and $e^{A+B}$. Applying (5.8) and (5.7) to (5.6) and using the bounds in (2.2.10) we obtain

$$
\Xi_{A} \geqslant \Xi_{\Lambda_{1}} \Xi_{\Lambda_{2}} \exp \left\{-c \sum_{\substack{Q(r) \subset \Lambda_{1} \\ \varrho\left(r^{\prime}\right) \subset \Lambda_{2}}} \tilde{\varphi}\left(\left|r-r^{\prime}\right|\right)\right\}
$$

The bounds in (5.5) now follows from the decay property in (2.2.14) and the above bounds. This completes the proof of the proposition.

We next prove (5.3) and so complete the proof of Theorem 2.2.4.
Proposition 5.2. Let the interaction satisfy Assumption A. Then

$$
P_{0.4}-P_{A} \rightarrow 0
$$

as $\Lambda$ tends to $\mathbf{R}^{v}$ in the sense of van Hove.
The main idea in proving the above proposition is careful applications of the methods in Section 3 and Section 4 with convex combinations of Wiener measures corresponding to different boundary conditions.

Before proving the proposition we need some preparations. Let $P_{\sigma, 1}^{\beta}(x, y ; d \omega)$ be the conditional Wiener measure corresponding to the transition function

$$
P_{\sigma, A}^{t}(x, y)=\exp \left[-\frac{1}{2} t \Delta_{\sigma, A}\right](x, y)
$$

the kernel of $\exp \left[-\frac{1}{2} t \Delta_{\sigma, A}\right]$. For a notational convenience we write that for $s \in[0,1]$

$$
\begin{equation*}
P_{A, s}(x, y ; d \omega)=s P_{0, A}^{\beta}(x, y ; d \omega)+(1-s) P_{\infty, A}^{\beta}(x, y ; d \omega) \tag{5.9}
\end{equation*}
$$

That is, $P_{A, s}(x, y ; d \omega)$ is the convex combination of the Wiener measures corresponding to Neumann and Dirichlet boundary conditions. We denote that
$\Xi_{A}(s)=\sum_{n=0}^{\infty} \frac{z^{n}}{n!} \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} \int_{\Omega^{n}} P_{A, s}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) e^{-U\left((\omega)_{n}\right)}$
and for any function $A\left((x)_{n}\right)$ invariant under $S_{n}$

$$
\begin{align*}
\rho_{\Lambda, s}(A)= & \Xi_{\Lambda}(s)^{-1} \sum_{n=0} \frac{z^{n}}{n!} \sum_{\pi \in S_{n}} \int_{A^{n}} d(x)_{n} A\left((x)_{n}\right) \int_{\Omega^{n}} P_{A, s}\left((x)_{n}, \pi(x)_{n} ; d(\omega)_{n}\right) \\
& \times e^{-U\left((\omega)_{n}\right)} \tag{5.11}
\end{align*}
$$

From the definitions it follows that

$$
\begin{gather*}
\Xi_{\Lambda}(1)=\Xi_{0, A}, \quad \Xi_{\Lambda}(0)=\Xi_{\infty, A}=\Xi_{A}  \tag{5.12}\\
\rho_{A, 1}(A)=\rho_{0, A}(A) \quad \text { and } \quad \rho_{A, 0}(A)=\rho_{A}(A)
\end{gather*}
$$

Let $\mathbf{K}_{\Delta, l, s}$ and $\overline{\mathbf{K}}_{\Delta, l, s}$ be the operators corresponding to the kernels

$$
\begin{equation*}
\mathbf{K}_{\Delta, l, s}(x, y)=\int_{\mathscr{E}(\Delta, l)} P_{A, s}(x, y ; d \omega) \tag{5.13}
\end{equation*}
$$

and

$$
\begin{align*}
\overline{\mathbf{K}}_{\Delta, l, s}(x, y) & =P_{A(A, l+1)}\left(s E^{(\beta / 4) \Delta_{0, A}}+(1-s) e^{(\beta / 4) \Delta_{A}}\right) P_{A(\Delta, l+1)}(x, y) \\
& =\int_{\Omega} \chi_{A(A, l+1)}(\omega) P_{A, s}^{\beta / 2}(x, y ; d \omega) \tag{5.14}
\end{align*}
$$

respectively [see (3.2.23) and (3.2.24)]. Following the method used in the proof of Proposition III.2.2 of Ref. 10 it is easy to check that Lemma 3.2.3 holds for $\mathbf{K}_{\Delta, l, s}$ and $\overline{\mathbf{K}}_{\Delta, l, s}$. Furthermore all arguments and results in Section 4 hold even if one replaces $P_{A}^{\beta}(x, y ; d \omega)$ by $P_{A, s}^{\beta}(x, y ; d \omega)$.

We next use the fundamental theorem of calculus to get

$$
\begin{align*}
P_{0, \Lambda}-P_{\Lambda} & =\frac{1}{|\Lambda|} \int_{0}^{1} d\left(\log \Xi_{A}(s)\right)  \tag{5.15}\\
& =\frac{1}{|\Lambda|} \int_{0}^{1}\left[\frac{1}{\Xi_{A}(s)}\right] \frac{d}{d s}\left(\Xi_{\Lambda}(s)\right) d s
\end{align*}
$$

From (5.10) it follows that

$$
\begin{equation*}
\frac{d}{d s} \Xi_{A}(s)=\sum_{Q(r): Q(r) \cap A \neq 0} \Xi_{A}^{\prime}(s, r) \tag{5.16}
\end{equation*}
$$

where

$$
\begin{align*}
\Xi_{A}^{\prime}(s, r)= & \sum_{n=0}^{\infty} \frac{z^{n}}{(n-1)!} \sum_{\pi \in S_{n}} \int_{Q(r) \cap A} d x_{1} \int_{A^{n-1}} d(x)_{n-1} \\
& \times \int_{\Omega}\left[\left.P_{A, s}\left(x_{1}, \pi\left(x_{1}\right) ; d \omega_{1}\right)\right|_{0} ^{1}\right] \\
& \times \int_{\Omega^{n-1}} P_{A, s}\left((x)_{n-1}, \pi(x)_{n-1} ; d(\omega)_{n-1}\right) \\
& \times e^{-U\left((\omega)_{n}\right)} \tag{5.17}
\end{align*}
$$

We will show that

$$
\begin{equation*}
\Xi_{A}^{\prime}(s, r) \leqslant\left\{e^{D} \exp \left[-\frac{1}{90 \beta} d(Q(r), \partial \Lambda)^{2}\right]\right\} \Xi_{A}(s) \tag{5.18}
\end{equation*}
$$

for some constant $D>0$, where $d(Q(r), \partial \Lambda)$ is the distance between $Q(r)$ and $\partial A$. Then from (5.15), (5.16), and (5.18) it follows that

$$
\begin{aligned}
\left|P_{0, \Lambda}-P_{\Lambda}\right| & \leqslant c|\partial \Lambda| /|\Lambda| \\
& \rightarrow 0 \quad \text { as } \Lambda \rightarrow \mathbf{R}^{v}
\end{aligned}
$$

In the remainder of this section we prove the bounds in (5.18).
Let $\mathbf{K}_{\Delta, l}^{\prime}$ be the operators of which the kernels are given by

$$
\begin{equation*}
\mathbf{K}_{\Delta, l}^{\prime}(x, y)=\int_{\mathscr{E}(\Delta, l)}\left[\left.P_{A, s}^{\beta}(x, y ; d \omega)\right|_{0} ^{1}\right. \tag{5.19}
\end{equation*}
$$

We then have the following result:
Lemma 5.3. For $l=-1,0,1,2, \ldots$ there are constants $c_{1}>0$ and $c_{2}>0$ such that
(a) $\operatorname{Tr}_{L^{2}(A)}\left(\mathbf{K}_{\Delta, l, s}\right) \leqslant c_{1}|\Lambda(\Delta, l+1)| \exp \left(-\frac{1}{16 \beta} l^{2}\right)$
(b) $\left|\operatorname{Tr}_{L^{2}(A)}\left(\mathbf{K}_{\Delta, l}^{\prime}\right)\right| \leqslant c_{2} \exp \left(\frac{-1}{60 \beta} l^{2}\right) \exp \left[-\frac{1}{60 \beta} d(\Delta, \partial \Lambda)^{2}\right]$

Proof. As we stated, the part (a) follows from the method similar to that used in the proof of Lemma 3.2.3 (see the proof of Proposition III.2.2 of Ref. 10). To prove the part (b) we note that

$$
\mathbf{K}_{\Delta, \lambda}^{\prime}(x, x)=\int_{\varepsilon(A, l)}\left[P_{0, A}(x, x ; d \omega)-P_{A}(x, x ; d \omega)\right]
$$

If $l \leqslant \frac{2}{3} d(\Delta, \partial A)$, then

$$
\begin{aligned}
\left|\mathbf{K}_{\Delta, l}^{\prime}(x, x)\right| \leqslant & \left|\chi_{A(\Delta, l+1)}(x) \int\left[P_{0, A}(x, x ; d \omega)-P_{A}(x, x ; d \omega)\right]\right| \\
\leqslant & \left|\chi_{A(A, l+1)}(x)\left[\int P(x, x ; d \omega)-\int P_{0, A}(x, x ; d \omega)\right]\right| \\
& +\left|\chi_{A(A, l+1)}(x)\left[\int P(x, x ; d \omega)-\int P_{A}(x, x ; d \omega)\right]\right| \\
& \leqslant c \exp \left\{-\frac{1}{4 \beta}\left[\frac{1}{3} d(A, \partial \lambda)\right]^{2}\right\} \chi_{A(A, l+1)}(x)
\end{aligned}
$$

Here we have use Theorem 6.3.8 of Ref. 1 to get the last inequality. Hence

$$
\left|\operatorname{Tr}\left(\mathbf{K}_{\Delta, l}^{\prime}\right)\right| \leqslant c \exp \left[-\frac{1}{37 \beta} d(\Delta, \partial \Lambda)^{2}\right]
$$

On the other hand, if $l>\frac{2}{3} d(\Lambda, \partial A)$, then

$$
\begin{align*}
\left|\operatorname{Tr}\left(\mathbf{K}_{\Delta, l}^{\prime}\right)\right| \leqslant & \int d^{v} x \int_{\delta(A, l)} P_{0, A}(x, x ; d \omega) \\
& +\int d^{v} x \int_{\mathscr{E}(\Lambda, l)} P_{A}(x, x ; d \omega) \\
\leqslant & c|A(\Delta, l+1)| \exp \left(-\frac{1}{16 \beta} l^{2}\right)  \tag{5.20}\\
\leqslant & c^{\prime} \exp \left[-\frac{1}{37 \beta} d(\Delta, \partial \Lambda)^{2}\right]
\end{align*}
$$

Here we have use the part (a) of the lemma to get the second inequality. Thus we conclude that

$$
\begin{equation*}
\left|\operatorname{Tr}\left(\mathbf{K}_{\Delta, l}^{\prime}\right)\right| \leqslant c \exp \left[-\frac{1}{37 \beta} d(\Delta, \hat{\partial} \Lambda)^{2}\right] \tag{5.21}
\end{equation*}
$$

By the method used in (5.20) we also have

$$
\begin{equation*}
\left|\operatorname{Tr}\left(\mathbf{K}_{4, l}^{\prime}\right)\right| \leqslant c \exp \left[-\frac{1}{18 \beta} l^{2}\right] \tag{5.22}
\end{equation*}
$$

Combining (5.21) and (5.22) we prove the part (b).
We now begin to prove Proposition 5.2.
Proof of Proposition 5.2. As we remarked before it is sufficient to show the bounds in (5.18). The basic idea is a modification of the method used in Section 3.2. Since the detailed derivation of the bounds in (5.18) would be very long and since the main idea is same as that in Section 3.2, we will give a sketch of the proof by giving necessary modifications and replacements. We leave the details to the reader.

For a given $Q(r), Q(r) \cap A \neq \phi$, we use the following abused notation:

$$
\begin{equation*}
\Lambda_{q}=\prod_{i=1}^{\nu}\left[-l_{q}-\frac{1}{2}+r^{i}, l_{q}+\frac{1}{2}+r^{i}\right] \tag{5.23}
\end{equation*}
$$

where $l_{q}$ is given in (3.1.9). $\Lambda_{q}$ is the closed cube centered at $r$. In the rest of this section $\Lambda_{q}$ is the cube defined in (5.23), not in (3.1.9). For a given integer we use the decompositions in (3.1.11) and (3.1.19) for $\Omega^{n-1}$ :

$$
\Omega^{n-1}=\mathscr{E}_{0} \cup\left(\bigcup_{q \geqslant q_{0}} \bigcup_{k} \mathscr{E}_{q, k-1}\right)
$$

Following each step used to derive the expressions in (3.1.20) and (3.1.26) we obtain

$$
\begin{equation*}
\Xi_{A}^{\prime}(s, r)=\widetilde{G}_{0}+\sum_{n=0}^{\infty} \sum_{q \geqslant q 0} \sum_{k} \widetilde{G}^{(n)}(q, k) \tag{5.24}
\end{equation*}
$$

where $\widetilde{G}_{0}$ is the expression obtained by repacing $\Omega^{n-1}$ in (5.17) by $\mathscr{E}_{0}$ and

$$
\begin{align*}
\tilde{G}^{(n)}(q, k)= & \frac{z^{n}}{(k-1)!(n-k)!} \sum_{\pi \in S_{n}} \int_{Q(r) \cap A} d x_{1} \int_{A^{n-1}} d(x)_{n-1} \\
& \times \int_{\Omega}\left[\left.P_{\Lambda, s}\left(x_{1}, \pi\left(x_{1}\right) ; d \omega_{1}\right)\right|_{0} ^{1}\right] \\
& \times \int_{\tilde{\delta}_{q, k-1}} P_{\Delta, s}\left((x)_{n-1}, \pi(x)_{n-1} ; d(\omega)_{n-1}\right) e^{-U\left((\omega)_{n}\right)} \tag{5.25}
\end{align*}
$$

where $\overline{\mathscr{E}}_{q, k-1}=\left\{(\omega)_{n-1} \in \mathscr{E}_{q}: \omega_{i}(\tau) \in A_{q+1}\right.$ for some $\left.\tau \in[0, \beta], i=2, \ldots, k\right\}$.

We assert that for $q_{0}$ sufficiently large one can choose $1<l<p<$ $\gamma(v, 2)$ such that the following bounds hold: For any $b>0$ there exist positive constants $c_{1}, c_{2}$ and $\varepsilon$ such that

$$
\begin{equation*}
U\left(\omega_{1}\right)+b V\left(\omega_{1}\right)+W\left(\omega_{1},(\omega)_{n-1}\right) \geqslant c_{1} \tag{5.26}
\end{equation*}
$$

on $\mathscr{E}_{0}$, and

$$
\begin{gather*}
U\left((\omega)_{k}\right)+b V\left((\omega)_{k}\right)+W\left((\omega)_{k},(\omega)_{n-k}\right) \\
\geqslant c_{3}\left\{\left|\Lambda_{q}\right|^{1+[(l-1)(1+\varepsilon) / p]}+k^{1+\varepsilon}\right\} \tag{5.27}
\end{gather*}
$$

on $\overline{\mathscr{E}}_{q, k-1}$. We prove our assertion. We note that $E_{q}<\left|A_{q}\right|^{2}$ on $\mathscr{E}_{0}$ for all $q \geqslant q_{0}$. In order to show (5.26) we consider two cases: If $\omega_{1}(\tau) \in A_{40}$ for all $\tau \in[0, \beta]$, then (5.26) follows from the lower regularity condition [Assumption A(b)]. If $\omega_{1}\left(\tau_{0}\right) \notin \Lambda_{q_{0}}$ for some $\tau_{0} \in[0, \beta]$, then (5.26) follows from the methods similar to those used in the proofs of Lemma 4.7 and (4.15). Next, the bound in (5.27) can be obtained by the method used in the proof of Theorem 3.1.1 (see Section 4). We leave the detailed proofs of (5.26) and (5.27) to the reader.

As in (3.2.9) we define

$$
\begin{align*}
\tilde{F}_{0}^{(n)}(Q(r), l)= & \sum_{\pi \in S_{n}} \int_{Q(r)} d x_{1} \int_{A^{n-1}} d(x)_{n-1} \tilde{\mathbf{K}}_{0}\left((x)_{n}, \pi(x)_{n}\right) \\
& \times \tilde{F}^{(n)}\left(q, k ;\left(Q(r), l_{1}\right),\left(A_{2}, l_{2}\right), \ldots,\left(A_{k}, l_{k}\right)\right) \\
= & \sum_{\pi \in S_{n}} \int_{Q(r)} d x_{1} \int_{A^{n-1}} d(x)_{n-1} \tilde{\mathbf{K}}\left((x)_{n}, \pi(x)_{n}\right) \tag{5.28}
\end{align*}
$$

where

$$
\begin{align*}
& \tilde{\mathbf{K}}_{0}\left((x)_{n}, \pi(x)_{n}\right)=\int_{\delta(Q(r), l)}\left[\left.P_{A, s}\left(x_{1}, \pi(x)_{1}, d \omega_{1}\right)\right|_{0} ^{1}\right] \\
& \times \int_{\bar{\delta}_{q, k-1}} P_{A, s}\left((x)_{n-1}, \pi(x)_{n-1} ; d(\omega)_{n-1}\right) e^{-U\left((\omega)_{n-1}\right)} \\
& \widetilde{\mathbf{K}}\left((x)_{n}, \pi(x)_{n}\right)=\int_{\mathcal{B}\left(Q(r), l_{1}\right)}\left[\left.P_{\Lambda, s}\left(x_{1}, \pi(x)_{1} ; d \omega_{1}\right)\right|_{0} ^{1}\right] \\
& \times \int_{\bar{G}\left(q, k-1 ;\left(\Lambda_{2}, l 2\right) \ldots,\left(d_{k}, l k\right)\right.} P_{A, s}\left((x)_{n-1}, \pi(x)_{n-1} ; d(\omega)_{n-1}\right) \\
& \times e^{-U\left((\omega)_{n-k}\right)} \tag{5.29}
\end{align*}
$$

Let $\mathbf{E}_{n, s}$ be the operator of which its kernel is given by

$$
\begin{equation*}
\mathbf{E}_{n, s}\left((x)_{n},(y)_{n}\right)=\int_{\Omega^{n}} P_{A, s}\left((x)_{n},(y)_{n} ; d(\omega)_{n}\right) e^{-U\left((\omega)_{n}\right)} \tag{5.30}
\end{equation*}
$$

Notice that

$$
\begin{equation*}
\Xi_{A}(s)=\sum_{n=0}^{\infty} \operatorname{Tr}_{\mathscr{H}_{n}^{(s)}(A)}\left(\mathbf{E}_{n, s}\right) \tag{5.31}
\end{equation*}
$$

Using the bounds in (5.26) and (5.27) and following the steps to derive (3.2.8) we obtain the following bounds:

$$
\begin{align*}
\widetilde{G}_{0} \leqslant & c_{1} \sum_{n=0}^{\infty} \frac{z^{n}}{(n-1)!} \sum_{l=-1}^{\infty} \exp \left[4 v b(l+2)^{2}\right] \widetilde{F}_{0}(Q(r), l) \\
\widetilde{G}^{(n)}(q, k) \leqslant & c_{2}\left\{\frac{z^{n}}{(k-1)!(n-k)!} \sum_{\Delta_{2} \subset \Lambda_{q+1}} \cdots \sum_{A_{k} \in A_{q+1}} \sum_{l_{1}=-1}^{\infty} \cdots \sum_{l_{k}=-1}^{\infty}\right. \\
& \left.\times \exp \left[4 v b \sum_{i=1}^{k}\left(l_{i}+2\right)^{2}\right] \tilde{F}^{(n)}\left(q, k ;\left(Q(r), l_{1}\right), \ldots,\left(\Delta_{k}, l_{k}\right)\right)\right\} \\
& \times \exp \left[-c\left(\left|\Lambda_{q}\right|^{1+[(l-1)(1+\varepsilon) / p]}+k^{1+\varepsilon)]}\right.\right. \tag{5.32}
\end{align*}
$$

What we have done up to now is the derivation of the forms in (5.24) and (5.32) so that the methods used in Section 3.2 can be applied. To prove the bound in (5.18) we use the following replacements in Section 3.2:

$$
\left.\begin{array}{l}
\mathbf{K}_{A_{1}, l_{1}}\left(x_{1}, y_{1}\right) \\
\mathbf{K}_{\Delta_{i} l_{i}}\left(x_{i}, y_{i}\right) \\
i=2, \ldots, k \\
\exp \left[-\beta H_{\Lambda}^{(n)}\right]
\end{array}\right\} \quad \text { by } \quad\left\{\begin{array}{l}
\mathbf{K}_{Q(r), l_{1}}^{\prime}\left(x_{1}, y_{1}\right) \\
\mathbf{K}_{A_{i}, l_{i}, s}\left(x_{i}, y_{i}\right) \\
i=2, \ldots, k \\
\mathbf{E}_{n, s}
\end{array}\right.
$$

We then use Lemma 5.3 instead of Lemma 3.2.3. Following each step used in the second part of (3.2.9) in Section 3.2, one may get the bounds of the following forms from (5.32) and (5.33):

$$
\begin{aligned}
\widetilde{G}_{0} \leqslant & c_{1} \exp \left[-\frac{1}{90 \beta} d(Q(r), \partial \Lambda)^{2}\right] \sum_{n=1}^{\infty} \operatorname{Tr}_{\mathscr{H}_{n}^{(s)}(\Lambda)}\left(\mathbf{E}_{n, s}\right) \\
G^{(n)}(q, k) \leqslant & c_{2} \exp \left[-c(q+k)^{1+\varepsilon}\right] \exp \left[\frac{-1}{90 \beta} d(Q(r), \partial \Lambda)^{2}\right] \\
& \times \operatorname{Tr}_{\mathscr{H}_{n-k}^{(s)}(\Lambda)}\left(\mathbf{E}_{n-k, s}\right)
\end{aligned}
$$

The bound in (5.18) follows from (5.24), (5.31) and the above inequalities. This completes the proof of Proposition 5.2.

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## Appendix. Corrections of Printing Errata in Ref. 10

$n \downarrow$ means $n$th line from the top; $n \uparrow$ means $n$th line from the bottom:

| Page |  | Replace | With |
| :---: | :---: | :---: | :---: |
| 11 | 7」 | $\exp \left(-\frac{1}{16 \beta}\right)$ | $\exp \left(-\frac{l^{2}}{16 \beta}\right)$ |
| 25 | $14 \downarrow$ | $-\beta H_{A}^{\left(n^{\prime}\right)}$ | $\exp \left[-\beta H_{A}^{\left(n^{\prime}\right)}\right]$ |
| 29 | 12 $\downarrow$ | $\sum_{\substack{r \in A_{4 q}^{c} \\ r \in A_{2 q}}}$ | $\sum_{\substack{r \in A \\ r \in A_{2 q}^{2}}}$ |
| 29 | $8 \uparrow$ | $\bar{n}\left(r^{\prime}, \tau\right) n(r, \tau) n(r, \tau)$ | $\bar{n}\left(r^{\prime}, \tau\right) n(r, \tau)$ |
| 30 | 15 $\downarrow$ | $\sum_{r \in A_{4 q}^{c}}$ | $\sum_{r \in A}$ |
| 32 | $4 \downarrow$ | $\int_{d}^{v}(x)$ | $\int d^{v} x$ |

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