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Har moni c functions and fundarent al sol utions of non－I ocal oper at ors gener ating jump Narkov pr ocesses

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# Harmonic functions and fundamental solutions of non-local operators generating jump Markov processes 

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## Contents

Introduction ..... 3
1 Preliminaries ..... 8
1.1 Dirichlet forms and symmetric Hunt processes ..... 8
1.2 Smooth measures and additive functionals ..... 12
2 Perturbation of Dirichlet forms and stability of fundamental solutions ..... 14
2.1 Heat kernel estimates for jump Markov processes ..... 15
2.2 Properties of Kato class measures ..... 21
2.3 Stability of fundamental solutions ..... 28
3 Critical Schrödinger forms and spectral functions ..... 34
3.1 The asymptotic behavior of the resolvent ..... 36
3.2 Growth of Feynman-Kac semigroups ..... 39
3.3 Differentiability of spectral functions ..... 48
4 Continuity of harmonic functions for non-local Markov generators ..... 51
4.1 Preceding results on continuous harmonic functions ..... 52
4.2 Condition (SSJ) and continuous harmonic functions ..... 58
4.3 Condition (QRI) and continuous harmonic functions ..... 60

## Introduction

There is a one to one correspondence between Markov processes and Markov generators. We denote them $\left\{X_{t}\right\}_{t \geq 0}$ and $\mathscr{L}$. When a Markov process $\left\{X_{t}\right\}_{t \geq 0}$ on $\mathbb{R}^{d}$ is given, the behavior of its harmonic function $u(x)$ and the estimates of transition density function $p(t, x, y)$ play a crucial role for knowing the properties of the Markov process. Here, the harmonic function is probabilistically defined as a bounded function such that $\left\{u\left(X_{t}\right)\right\}_{t \geq 0}$ is a martingale and the transition density function is characterized by

$$
\mathbb{E}_{x}\left[f\left(X_{t}\right)\right]=\int_{\mathbb{R}^{d}} p(t, x, y) f(y) d y
$$

The harmonic function $u(x)$ is also characterized as $\mathscr{L} u=0$ and $p(t, x, y)$ is the fundamental solution of equation $\partial u / \partial t=\mathscr{L} u$. A most important Markov process is the standard Brownian motion, whose Markov generator is equal to $(1 / 2) \Delta$. In this case $p(t, x, y)$ is the fundamental solution of the partial differentiable equation $\partial u / \partial t=$ $(1 / 2) \Delta u$. We often call $p(t, x, y)$ heat kernel.

There are a lot of preceding results on the continuity of harmonic functions and estimates of heat kernels. For example, when Markov generator $\mathscr{L}$ is a uniformly elliptic second order differential operator, the corresponding Markov process is a diffusion process and the continuity of harmonic functions is well known. Moreover, this heat kernel $p(t, x, y)$ admits the Gaussian estimate as follows:

$$
C_{1} t^{-\frac{d}{2}} \exp \left(-\frac{C_{2}|x-y|^{2}}{t}\right) \leq p(t, x, y) \leq C_{3} t^{-\frac{d}{2}} \exp \left(-\frac{C_{4}|x-y|^{2}}{t}\right)
$$

Here $C_{i}$ 's are positive constants. The studies on the continuity of harmonic functions and estimate of heat kernels for jump Markov processes have been developed for the last decade. Jump Markov processes have discontinuous sample paths. A most typical example is the rotationally invariant $\alpha$-stable process generated by the fractional Laplacian $-(-\Delta)^{\alpha / 2}$, where $0<\alpha<2$. Unlike diffusion processes, harmonic functions of jump Markov processes are not necessarily continuous (Barlow, Bass and et. al. [3]). Thus, for the continuity of harmonic functions we need to impose some conditions on the generator. This problem is considered by Bass and Levin [5], Bass and Kassmann [4] and Husseini and Kassmann [14]. Moreover, it is proved that heat kernels of jump Markov processes admit two-sided estimates different from those of diffusion processes. Chen and Kumagai [9] treated $\alpha$-stable-like processes, which is a generalization of the rotationally invariant $\alpha$-stable process. In Chen, Kim and Kumagai [7], they gave two-sided estimates for heat kernel of relativistic $\alpha$-stable-like
processes. These processes are symmetric with respect to the Lebesgue measure $m$. To analyze symmetric Markov processes, the Dirichlet form theory is a powerful tool. We introduce the Dirichlet form $(\mathscr{E}, \mathscr{F})$ corresponding to the Markov process $\left\{X_{t}\right\}_{t \geq 0}$ :

$$
\mathscr{E}(u, v)=(-\mathscr{L} u, v),
$$

where $(\cdot, \cdot)$ is the inner product of $L^{2}\left(\mathbb{R}^{d}\right)$ and $\mathscr{F}$ is the domain of the form $\mathscr{E}$, the closure of the domain of $\mathscr{L}$. We define a Schrödinger form $\mathscr{E}^{\mu}$ by

$$
\begin{equation*}
\mathscr{E}^{\mu}(u, v)=\mathscr{E}(u, v)-\int_{\mathbb{R}^{d}} u(x) v(x) \mu(d x) . \tag{0.1}
\end{equation*}
$$

Here $\mu$ is a suitable positive measure. We note that the second term of ( 0.1 ) contributes not killing but creation of the process. We consider the fundamental solution of $\partial u / \partial t=\mathscr{L}^{\mu} u$, where $\mathscr{L}^{\mu}$ is the Schrödinger type operator defined by

$$
\left(-\mathscr{L}^{\mu} u, v\right)=\mathscr{E}^{\mu}(u, v)
$$

We denote by $p^{\mu}(t, x, y)$ the fundamental solution. Since the perturbation in a Schrödinger form is defined by creation, $p^{\mu}(t, x, y)$ is no longer a transition density of a Markov process. However, we can define the corresponding semigroup by using the positive continuous additive functional $A_{t}^{\mu}$. Here $A_{t}^{\mu}$ is determined uniquely from $\mu$ by the Revuz correspondence. Then $p^{\mu}(t, x, y)$ is the integral kernel of the Feynman-Kac semigroup:

$$
\begin{equation*}
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right) f\left(X_{t}\right)\right]=\int_{\mathbb{R}^{d}} p^{\mu}(t, x, y) f(y) d y \tag{0.2}
\end{equation*}
$$

We compare $p^{\mu}(t, x, y)$ with $p(t, x, y)$. If $p^{\mu}(t, x, y)$ has the same type two-sided estimate as $p(t, x, y)$ up to positive constants, we call this phenomenon the stability of fundamental solutions. Intuitively, if the potential $\mu$ is large, $p^{\mu}(t, x, y)$ has a different estimate from $p(t, x, y)$. Thus we need to formulate the smallness of the measure $\mu$. To this end, we use the bottom of spectrum of the time changed process by $A_{t}^{\mu}$. We recall a result for a transient Brownian motion. Assume $\mu$ is in a certain class. Takeda [32] showed that under a certain condition on $\mu$, the stability of fundamental solution holds if and only if

$$
\begin{equation*}
\lambda(\mu):=\inf \left\{\frac{1}{2} \mathbb{D}(u, u) ; u \in H^{1}\left(\mathbb{R}^{d}\right), \int_{\mathbb{R}^{d}} u^{2} d \mu=1\right\}>1 \tag{0.3}
\end{equation*}
$$

where $\mathbb{D}$ is the Dirichlet integral defined by

$$
\mathbb{D}(u, v)=\int_{\mathbb{R}^{d}} \nabla u(x) \cdot \nabla v(x) d x
$$

and $H^{1}\left(\mathbb{R}^{d}\right)$ is the 1 -order Sobolev space. $\lambda(\mu)$ is regarded as the bottom of spectrum of the operator $(1 / 2 \mu) \Delta$, the generator of the time changed process of the Brownian motion by $A_{t}^{\mu}$. The formula (0.3) describes the smallness of the measure $\mu$. Indeed, if
$\mu_{1} \leq \mu_{2}$, then $\lambda\left(\mu_{1}\right) \geq \lambda\left(\mu_{2}\right)$. Moreover, the following each statement is equivalent to (0.3):

$$
\begin{gather*}
\sup _{x \in \mathbb{R}^{d}} \mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right)\right]<\infty,  \tag{0.4}\\
G^{\mu}(x, y):=\int_{0}^{\infty} p^{\mu}(t, x, y) d t<\infty \quad \text { for } x \neq y . \tag{0.5}
\end{gather*}
$$

If the measure $\mu$ satisfies ( 0.4 ) and (0.5), then $\mu$ is said to be gaugeable and subcritical respectively. For the proof of the stability of fundamental solutions, the equivalence of (0.3)-(0.5) is crucial. In this thesis, the measure $\mu$ is said to be subcritical if $\lambda(\mu)>1$. Furthermore,
(1) $\mu$ is said to be critical if $\lambda(\mu)=1$,
(2) $\mu$ is said to be supercritical if $\lambda(\mu)<1$.

In two cases above, we cannot expect that $p^{\mu}(t, x, y)$ have the same two-sided estimate as $p(t, x, y)$. We have not been able to give a two-sided estimate of $p^{\mu}(t, x, y)$. However, there exist some papers which deal with large time asymptotics of the Feynman-Kac semigroups, $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$, if $\mu$ is critical or supercritical. We see from (0.2) that

$$
\begin{equation*}
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=\int_{\mathbb{R}^{d}} p^{\mu}(t, x, y) d y \tag{0.6}
\end{equation*}
$$

Since the subcriticality is equivalent to the gaugeability, the Feynman-Kac semigroup diverges as $t \rightarrow \infty$. Takeda $[28,30]$ showed that

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \frac{1}{t} \log \mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=C(\mu) \tag{0.7}
\end{equation*}
$$

where

$$
C(\mu)=-\inf \left\{\mathscr{E}^{\mu}(u, u) ; \int_{\mathbb{R}^{d}} u^{2} d x=1\right\}
$$

that is, $-C(\mu)$ is the bottom of spectrum of the operator $-\mathscr{L}^{\mu}=-(\mathscr{L}+\mu)$. Hence we see that if $\mu$ is supercritical, then $C(\mu)$ is positive and the expectation (0.6) grows exponentially.

If $\mu$ is critical, $C(\mu)=0$ and the expectation (0.6) seems to have polynomial growth. This conjecture is proved by Simon [25] and Cranston, Kolokoltsov and et. al. [8] in case that $\mu$ is absolutely continuous with respect to the Lebesgue measure. More precisely, when $\mu=V(x) d x$ for $V \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$, the growth of the expectation (0.6) has the following asymptotics:

$$
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right] \sim \begin{cases}C_{1} t^{\frac{1}{2}} & (d=3) \\ C_{2} t / \log t & (d=4) \\ C_{3} t & (d \geq 5)\end{cases}
$$

In this thesis we extend these results to jump Markov processes. The Dirichlet form associated with jump Markov process $\left\{X_{t}\right\}_{t \geq 0}$ is expressed by

$$
\mathscr{E}(u, v)=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}(u(y)-u(x))(v(y)-v(x)) J(x, y) d x d y .
$$

Here $J(x, y)$ is a positive symmetric Borel function called jump intensity measure. The terms ' $\alpha$-stable-like process' and 'relativistic $\alpha$-stable-like process' come from the behavior of jump intensity measure. In this thesis, we introduce three classes of Radon measures: the Kato-class $\mathbf{K}$, Green-tight Kato class $\mathbf{K}_{\infty}$ and conditional Green-tight Kato class $\mathbf{S}_{\infty}$ (For the definitions, see Section 2.2). In particular, the class $\mathbf{K}_{\infty}$ plays a crucial role. For $\mu \in \mathbf{K}_{\infty}$, we consider the Schrödinger form $\mathscr{E}^{\mu}$ defined in (0.1) and denote by $\left\{P_{t}^{\mu}\right\}$ a semigroup generated by $\mathscr{E}^{\mu}$. Using a harmonic function $h(x)$ of Schrödinger operator, the Schrödigner form can be transformed to a Dirichlet form. This method is called Doob's $h$-transform used intensively in Chen and Zhang [12]. We consider the transformed semigroup $\left\{P_{t}^{\mu, h}\right\}_{t \geq 0}: L^{2}\left(h^{2} m\right) \rightarrow L^{2}\left(h^{2} m\right)$,

$$
P_{t}^{\mu, h} f(x)=\mathbb{E}_{x}\left[\frac{h\left(X_{t}\right)}{h\left(X_{0}\right)} \exp \left(A_{t}^{\mu}\right) f\left(X_{t}\right)\right]
$$

The transformed semigroup generates Markov process on $\mathbb{R}^{d}$ with symmetric measure $h^{2} m$. Moreover, we see that the associated Dirichlet form on $L^{2}\left(h^{2} m\right)$ is expressed as

$$
\begin{equation*}
\mathscr{E}^{\mu, h}(u, v)=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}(u(y)-u(x))(v(y)-v(x)) J(x, y) h(x) h(y) d x d y, \tag{0.8}
\end{equation*}
$$

namely, the jump intensity measure $J(x, y)$ is transformed to $J(x, y) h(x) h(y)$.
First we establish a necessary and sufficient condition on $\mu$ for the stability of fundamental solution when the Markov process $\left\{X_{t}\right\}_{t \geq 0}$ is $\alpha$-stable-like or relativistic $\alpha$-stable-like. If $\mu$ is subcritical, the gauge function $h(x)=\mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right)\right]$ satisfies $1 \leq h(x) \leq C_{1}$ for some positive constant. Hence $J(x, y) h(x) h(y)$ is equivalent to $J(x, y)$ and consequently the transition density function of the transformed process is equivalent to the original one. Noting that $p^{\mu}(t, x, y) / h(x) h(y)$ is equal to the transition density function of the transformed process, we can conclude the stability of fundamental solutions. For example, let $\left\{X_{t}\right\}_{t \geq 0}$ be a transient $\alpha$-stable-like process on $\mathbb{R}^{d}$. In this case, the jump intensity measure $J(x, y)$ satisfies

$$
\frac{C_{1}}{|x-y|^{d+\alpha}} \leq J(x, y) \leq \frac{C_{2}}{|x-y|^{d+\alpha}}
$$

for some positive constants $C_{1}$ and $C_{2}$. The transformed Dirichlet form ( 0.8 ) is also $\alpha$-stable-like if and only if the measure $\mu$ satisfies

$$
\begin{equation*}
\inf \left\{\mathscr{E}(u, u) \mid u \in \mathscr{F}, \int_{\mathbb{R}^{d}} u^{2} d \mu=1\right\}>1 . \tag{0.9}
\end{equation*}
$$

From Chen and Kumagai [9], we know the two-sided estimate for the heat kernel of $\alpha$-stable-like process:

$$
C_{1}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) \leq p(t, x, y) \leq C_{2}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right)
$$

Therefore, we can conclude that the stability of fundamental solution is equivalent to the condition (0.9), which is an extension of (0.3).

We next consider large time asymptotics of the expectation $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$, when $\mu$ is critical or supercritical. Since (0.7) is valid for the rotationally invariant $\alpha$-stable process ([33]), the expectation (0.6) grows exponentially if $\mu$ is supercritical.

If $\mu$ is critical, we have a concrete growth order of (0.6) for $d=2$ and $\alpha=1$, which is the same as that of 4 -dimensional Brownian motion. Takeda [34] proved that for the rotationally invariant $\alpha$-stable process with $d / \alpha>2$, the growth of ( 0.6 ) is proportional to the time $t$. In this case the function $h(x)$ belongs to $L^{2}\left(\mathbb{R}^{d}\right)$ and, as a result, the transformed Markov process has the finite invariant measure $h^{2} m$. Applying the ergodic theory, we can obtain the growth order of (0.6). But when $d / \alpha \leq 2$, we cannot use this argument because $h(x)$ is not in $L^{2}\left(\mathbb{R}^{d}\right)$. Hence, we apply an analytical methods due to Simon [19].

If $\mu$ is critical, the construction of $h(x)$ is based on Takeda and Tsuchida [35], in which they proved that $h(x)$ is continuous and satisfies $h(x) \asymp 1 \wedge|x|^{\alpha-d}$. Hence, in order to estimate the heat kernel of the $h$-transformed process, we need to treat the jump intensity like $h(x) h(y) /|x-y|^{d+\alpha}$, which depends not only $|x-y|$ but also $|x|$ and $|y|$. The estimate of heat kernels for these type jump processes is an interesting problem.

We finally consider harmonic functions with respect to the Markov generators. Bass and Levin [5] showed the Hölder continuity of harmonic functions for $\alpha$-stable-like generators. Bass and Kassmann [4] showed the continuity of harmonic functions for more general jump-type Markov generators. They impose two conditions on the jump intensity measure $J(x, y)$ : one is singularity of small jumps, namely how the amount of jumps with size $r$ grows as $r$ tends to 0 . The other is quasi rotationally invariance, namely how the process jumps in any direction to some extent. In this paper we prove the continuity of harmonic functions under conditions weaker than those in Bass and Kassmann [4]. For instance, we consider the case $J(x, y)$ satisfies

$$
\frac{C_{1}}{|x-y|^{d+1}} \log \frac{3}{|x-y|} \leq J(x, y) \leq \frac{C_{2}}{|x-y|^{d+1}} \log \frac{3}{|x-y|} \quad(0<|x-y| \leq 2)
$$

for some positive constants $C_{1}$ and $C_{2}$. This jump intensity measure is a bit different from that of $\alpha$-stable-like processes. For detail, see Example 4.1.

This thesis is organized as follows: In Chapter 1, we prepare the basic material: Dirichlet forms, Hunt processes, smooth measures and additive functionals. In Chapter 2 , we establish a necessary and sufficient condition on $\mu$ for the stability of fundamental solutions when the process is $\alpha$-stable-like or relativistic $\alpha$-stable-like. In Chapter 3 , for the rotationally invariant 1 -stable process on $\mathbb{R}^{2}$, we consider the growth order of Feynman-Kac expectations and differentiability of spectral functions associated with critical Schrödinger forms . In Chapter 4, we study the continuity of harmonic functions for a class of jump-type Markov generators containing $\alpha$-stable-like ones. Our result is an extension of preceding ones in Bass and Kassmann [4], and Husseini and Kassmann [14].

## Chapter 1

## Preliminaries

In this chapter we introduce basic material that will be used through this paper. In Section 1, we review basic properties of Dirichlet forms and Markov processes. In Section 2, we mainly treat the smooth measures and additive functionals.

### 1.1 Dirichlet forms and symmetric Hunt processes

Let $(\mathscr{E}, \mathscr{F})$ be a Dirichlet form on $L^{2}\left(\mathbb{R}^{d}\right)$. Here $\mathscr{E}$ is a non-negative definite symmetric bilinear form and $\mathscr{F}$ is an appropriate subspace of $L^{2}\left(\mathbb{R}^{d}\right)$ called domain. Define the norm on $\mathscr{F}$ by

$$
\mathscr{E}_{1}(u, v)=\mathscr{E}(u, v)+\int_{\mathbb{R}^{d}} u(x) v(x) d x
$$

Then $\mathscr{F}$ is a Hilbert space with respect to $\mathscr{E}_{1}^{1 / 2}$-norm. Let $C_{c}\left(\mathbb{R}^{d}\right)$ be the family of continuous functions on $\mathbb{R}^{d}$ with compact supports equipped with the uniform norm. If the domain $\mathscr{F}$ is both dense in $L^{2}\left(\mathbb{R}^{d}\right)$ with respect to the $L^{2}$-norm and in $C_{c}\left(\mathbb{R}^{d}\right)$ with respect to the uniform norm, the Dirichlet form $(\mathscr{E}, \mathscr{F})$ is said to be regular.

It is known that a regular Dirichlet form on $L^{2}\left(\mathbb{R}^{d}\right)$ has a unique representation as follows:

$$
\begin{align*}
& \mathscr{E}(u, v)=\mathscr{E} c \\
&(u, v)+\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}(u(y)-u(x))  \tag{1.1}\\
&(v(y)-v(x)) J(x, y) d x d y+\int_{\mathbb{R}^{d}} u(x) v(x) k(d x) .
\end{align*}
$$

We call (1.1) formulae of Beurling-Deny and LeJan. For further arguments, see $\S 3.2$ of [13]. Furthermore, by the general theory of regular Dirichlet forms, there exists a unique Hunt process on $\mathbb{R}^{d}$ associated with $(\mathscr{E}, \mathscr{F})$ (See Chapter 7 of [13]). Hunt processes on $\mathbb{R}^{d}$ consists of three parts according to the behaviors of sample paths. (1.1) enables us to distinguish these behaviors of sample paths: The first part of (1.1) is called diffusion part which satisfies strong local property: i.e. if $v$ is equal to some constant on the support of $u$, it follows that $\mathscr{E}^{c}(u, v)=0$. Diffusion part describes
the continuous movement of a particle, such as the Brownian motion. The second part of (1.1) is called jump part which satisfies non-local property and describes the discontinuous movement of a particle. The most typical example is the rotationally invariant $\alpha$-stable process for $0<\alpha<2$ (See Example 1.1). The last part of (1.1) is called killing part which describes the disappearance of a particle in the state space.

Next we review definitions and properties of Hunt process. Let $\left(\Omega, \mathscr{M},\left\{X_{t}\right\}_{t \geq 0}, \mathbb{P}\right)$ be a stochastic process with state space $\left(\mathbb{R}^{d}, \mathscr{B}\left(\mathbb{R}^{d}\right)\right)$, where $\mathscr{B}\left(\mathbb{R}^{d}\right)$ stands for the Borel $\sigma$-field on $\mathbb{R}^{d}$, i.e. $(\Omega, \mathscr{M}, \mathbb{P})$ is a probability space and each $X_{t}$ is a measurable map from $\Omega$ to $\mathbb{R}^{d}$. The last condition of the measurability is explicitly indicated by $X_{t} \in \mathscr{M} / \mathscr{B}\left(\mathbb{R}^{d}\right)$. We say that the family $\left\{\mathscr{M}_{t}\right\}_{t>0}$ of sub- $\sigma$-fields of $\mathscr{M}$ is an admissible filtration if $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$ is increasing in $t$ and $X_{t} \in \mathscr{M}_{t} / \mathscr{B}$ for each $t \geq 0$. An admissible filtration $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$ is called right continuous if for any $t \geq 0$,

$$
\mathscr{M}_{t}=\mathscr{M}_{t+}:=\cap_{t^{\prime}>t} \mathscr{M}_{t^{\prime}}
$$

Adjoining an extra point $\Delta$ to a measurable space $\left(\mathbb{R}^{d}, \mathscr{B}\left(\mathbb{R}^{d}\right)\right)$, we set

$$
\overline{\mathbb{R}^{d}}=\mathbb{R}^{d} \cup \Delta, \quad \mathscr{B}\left(\overline{\mathbb{R}^{d}}\right)=\mathscr{B}\left(\mathbb{R}^{d}\right) \cup\left\{B \cup \Delta: B \in \mathscr{B}\left(\mathbb{R}^{d}\right)\right\}
$$

A quadruple $\mathbb{M}=\left(\Omega, \mathscr{M},\left\{X_{t}\right\}_{t \geq 0},\left\{\mathbb{P}_{x}\right\}_{x \in \overline{\mathbb{R}^{d}}}\right)$ is said to be normal Markov process on $\left(\mathbb{R}^{d}, \mathscr{B}\left(\mathbb{R}^{d}\right)\right)$ if the following conditions are satisfied:
(M-1) For each $x \in \overline{\mathbb{R}^{d}},\left(\Omega, \mathscr{M},\left\{X_{t}\right\}_{t \geq 0}, \mathbb{P}_{x}\right)$ is a stochastic process on $\mathbb{R}^{d}$.
(M-2) $\mathbb{P}_{x}\left(X_{t} \in E\right)$ is measurable for $t \geq 0$ and $E \in \mathscr{B}\left(\mathbb{R}^{d}\right)$.
(M-3) There exists an admissible filtration $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$ such that

$$
\mathbb{P}_{x}\left(X_{t+s} \in E \mid \mathscr{M}_{t}\right)=\mathbb{P}_{X_{t}}\left(X_{s} \in E\right), \quad \mathbb{P}_{x} \text {-a.s. }
$$

for any $x \in \mathbb{R}^{d}, t, s \geq 0$ and $E \in \mathscr{B}\left(\mathbb{R}^{d}\right)$.
(M-4) $\mathbb{P}_{\Delta}\left(X_{t}=\Delta\right)=1$ for any $t \geq 0$.
(M-5) $\mathbb{P}_{x}\left(X_{0}=x\right)=1$ for any $x \in \mathbb{R}^{d}$.
Before we define a Hunt process, we introduce stopping time and sub- $\sigma$-field defined by stopping time. A $[0, \infty]$-valued function $\sigma$ on $\Omega$ is called an $\mathscr{M}_{t}$-stopping time if $\{\sigma \leq t\} \in \mathscr{M}_{t}$ for each $t \geq 0$. For a stopping time $\sigma$, we define the sub- $\sigma$-field $\mathscr{M}_{\sigma}$ by

$$
\mathscr{M}_{\sigma}=\left\{\Lambda \in \mathscr{M}: \Lambda \cap\{\sigma \leq t\} \in \mathscr{M}_{t} \text { for all } t \geq 0\right\} .
$$

Markov process $\mathbb{M}$ is said to be Hunt process, if the following additional conditions are satisfied:
(M-6) (i) $X_{\infty}(\omega)=\Delta$ for any $\omega \in \Omega$.
(ii) $X_{t}(\omega)=\Delta$ for $t \geq \zeta(\omega)$, where $\zeta(\omega)$ is the life time defined by $\zeta(\omega)=$ $\inf \left\{t \geq 0: X_{t}(\omega)=\Delta\right\}$
(iii) for each $t \in[0, \infty]$, there exists a map $\theta_{t}$ from $\Omega$ to $\Omega$ such that $X_{s} \circ \theta_{t}=X_{t+s}$ for $s \geq 0$.
(iv) For each $\omega \in \Omega$, the sample path $t \rightarrow X_{t}(\omega)$ is right continuous on $[0, \infty)$ and has the left limit on $(0, \infty)$.
(M-7) (i) The admissible filtration $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$ is right continuous and for any probability measure $v$ on $\overline{\mathbb{R}^{d}}, E \in \mathscr{B}\left(\overline{\mathbb{R}^{d}}\right)$ and $s \geq 0$,

$$
\mathbb{P}_{v}\left(X_{\sigma+s} \in E \mid \mathscr{M}_{\sigma}\right)=\mathbb{P}_{X_{\sigma}}\left(X_{s} \in E\right), \quad \mathbb{P}_{v} \text {-a.s. }
$$

where $\sigma$ is $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$-stopping time.
(ii) Let $v$ be an arbitrary probability measure on $\overline{\mathbb{R}^{d}}$. For any $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$-stopping time $\sigma_{n}$ increasing to $\sigma$, it holds that

$$
\mathbb{P}_{v}\left(\lim _{n \rightarrow \infty} X_{\sigma_{n}}=X_{\sigma}, \sigma<\infty\right)=\mathbb{P}_{v}(\sigma<\infty)
$$

The properties in (M-7) are called the strong Markov property and the quasi left continuity respectively.

In the sequel, we abbreviate the Hunt process by $\left\{X_{t}\right\}_{t \geq 0}$ for the convenience. Let $\left\{P_{t}\right\}_{t>0}$ be the semigroup generated by the Hunt process $\left\{X_{t}\right\}_{t \geq 0}$ :

$$
P_{t} f(x)=\mathbb{E}_{x}\left[f\left(X_{t}\right)\right]=\int_{\mathbb{R}^{d}} p(t, x, y) f(y) d y,
$$

where $p(t, x, y)$ is the transition probability density of $\left\{X_{t}\right\}_{t \geq 0}$. If we use the unique operator $\mathscr{L}$ satisfying $\mathscr{E}(u, v)=-(\mathscr{L} u, v), p(t, x, y)$ is the fundamental solution of the equation $\partial u / \partial t=\mathscr{L} u$. For $\beta>0$, we define the resolvent kernel of $\beta$-order as follows:

$$
G_{\beta}(x, y)=\int_{0}^{\infty} e^{-\beta t} p(t, x, y) d t
$$

If we admit the case the right hand side is equal to infinity, this formula is valid for $\beta=0$. We denote by $G(x, y)$ the resolvent kernel of 0 -order, which we call the green kernel. The Hunt process $\left\{X_{t}\right\}_{t \geq 0}$ or the associated Dirichlet form $(\mathscr{E}, \mathscr{F})$ is said to be transient if $G(x, y)<\infty$ for $x \neq y$. In order to characterize transience property, we often use the extended Dirichlet space $\mathscr{F}_{e}$. $\mathscr{F}_{e}$ is defined as a family of measurable functions on $\mathbb{R}^{d}$ satisfying the following two conditions:

- $|u|<\infty \quad m$-a.e.
- There exists an $\mathscr{E}$-Cauchy sequence $\left\{u_{n}\right\}_{n \in \mathbb{N}}$ of functions in $\mathscr{F}$ such that $\lim _{n \rightarrow \infty} u_{n}=u \quad m$-a.e.

Obviously $\mathscr{F}$ is a subspace of $\mathscr{F}_{e}$, and we see that $\mathscr{F}_{e} \cap L^{2}\left(\mathbb{R}^{d}\right)=\mathscr{F}$ from Theorem 1.5.2 of [13]. The following proposition is taken from Theorem 1.5.3 of [13], which describes necessary and sufficient conditions for the transience of Dirichlet form.

Proposition 1.1. $\mathscr{F}_{e}$ is the extended Dirichlet space of transient Dirichlet form $\mathscr{E}$ if and only if the following conditions are satisfied.
(i) $\mathscr{F}_{e}$ is a real Hilbert space with inner product $\mathscr{E}$.
(ii) There exists an m-integrable bounded function $g$ strictly positive m-a.e. such that $\mathscr{F}_{e} \subset L^{1}(g \cdot m)$ and

$$
\int_{\mathbb{R}^{d}}|u(x)| g(x) m(d x) \leq \sqrt{\mathscr{E}(u, u)}, \quad \forall u \in \mathscr{F}_{e}
$$

(iii) $\mathscr{F}_{e} \cap L^{2}\left(\mathbb{R}^{d}\right)$ is dense both in $L^{2}\left(\mathbb{R}^{d}\right)$ and in $\left(\mathscr{F}_{e}, \mathscr{E}\right)$.
(iv) For any $u \in \mathscr{F}_{e}$ and its normal contraction $v$, it follows that $v \in \mathscr{F}_{e}$ and $\mathscr{E}(v, v) \leq \mathscr{E}(u, u)$. Here we say $v$ is a normal contraction of $u$, if $|v(x)| \leq|u(x)|$ for all $x \in \mathbb{R}^{d}$ and $|v(x)-v(y)| \leq|u(x)-u(y)|$ for all $x, y \in \mathbb{R}^{d}$.

In this thesis, we consider the Hunt process which consists of only discontinuous sample paths. These Hunt process is called jump processes. (1.1) implies that the jump process $\left\{X_{t}\right\}_{t \geq 0}$ is characterized by a regular Dirichlet form $(\mathscr{E}, \mathscr{F})$ on $L^{2}\left(\mathbb{R}^{d}\right)$ as follows:

$$
\mathscr{E}(u, u)=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}(u(y)-u(x))^{2} J(x, y) d x d y, \quad \mathscr{F}=\overline{\left\{C_{c}\left(\mathbb{R}^{d}\right)\right\}^{\mathscr{E}_{1}^{1 / 2}}},
$$

where $J(x, y)$ is a symmetric Borel function. $J(x, y)$ describes the frequency of jump and is called jump intensity measure. We close this section by introducing three important examples of jump processes.

Example 1.1. If $J(x, y)=C_{\alpha} /|x-y|^{d+\alpha}$ for $0<\alpha<2$ and positive constant $C_{\alpha}$, $\left\{X_{t}\right\}_{t \geq 0}$ is called the rotationally invariant $\alpha$-stable process. The characteristic function satisfies

$$
\mathbb{E}_{0}\left[\exp \left(i u \cdot X_{t}\right)\right]=\exp \left(-t|u|^{\alpha}\right) .
$$

Example 1.2. We call the associated process $\left\{X_{t}\right\}_{t \geq 0} \alpha$-stable-like, if $J(x, y)$ satisfies

$$
\frac{C_{1}}{|x-y|^{d+\alpha}} \leq J(x, y) \leq \frac{C_{2}}{|x-y|^{d+\alpha}}
$$

for some positive constants $C_{1}$ and $C_{2}$.
Example 1.3. $\left\{X_{t}\right\}_{t \geq 0}$ is called relativistic $\alpha$-stable-like if $J(x, y)$ satisfies

$$
\frac{C_{1}}{|x-y|^{d+\alpha}} \exp \left(-m_{0}|x-y|\right) \leq J(x, y) \leq \frac{C_{2}}{|x-y|^{d+\alpha}} \exp \left(-m_{0}|x-y|\right)
$$

for some positive constants $C_{1}, C_{2}$ and $m_{0}$.

### 1.2 Smooth measures and additive functionals

In this section we define smooth measures and positive continuous additive functionals. In order to define smooth measures, we first review the definition of capacity. Denote by $\mathscr{O}$ the family of all open subset of $\mathbb{R}^{d}$. For $A \in \mathscr{O}$, we define

$$
\begin{aligned}
\mathscr{L}_{A} & =\{u \in \mathscr{F} ; u \geq 1 \\
\operatorname{Cap}(A) & = \begin{cases}\inf _{u \in \mathscr{L}_{A}} \mathscr{E}_{1}(u, u), & \mathscr{L}_{A} \neq \emptyset \\
\infty, & \mathscr{L}_{A}=\emptyset .\end{cases}
\end{aligned}
$$

For any set $A \subset \mathbb{R}^{d}$, we set

$$
\operatorname{Cap}(A)=\inf _{B \in \mathscr{O}, A \subset B} \operatorname{Cap}(B) .
$$

We call this the capacity of $A$. The set $A$ is said to be exceptional if $\operatorname{Cap}(A)=0$. We use the term quasi everywhere (in abbreviation, q.e.) in order to mention 'except for an exceptional set'.

A positive Radon measure $\mu$ on $\mathbb{R}^{d}$ is said to be of finite energy integral if

$$
\int_{\mathbb{R}^{d}}|v(x)| \mu(d x) \leq C \mathscr{E}_{1}(v, v) \quad\left(v \in \mathscr{F} \cap C_{c}\left(\mathbb{R}^{d}\right)\right)
$$

for some positive constant $C$. We denote by $\mathbf{S}_{0}$ the family of positive Radon measures of finite energy integral. Since $\mathscr{F}$ is a Hilbert space with respect to $\mathscr{E}_{1}$-norm, Riesz representation theorem implies that there exists a unique $U_{1} \mu \in \mathscr{F}$ such that

$$
\mathscr{E}_{1}\left(U_{1} \mu, v\right)=\int_{\mathbb{R}^{d}} v(x) \mu(d x) .
$$

We call the function $U_{1} \mu$ 1-potential. Moreover, we define a subset $\mathbf{S}_{00}$ of $\mathbf{S}_{0}$ by

$$
\mathbf{S}_{00}=\left\{\mu \in \mathbf{S}_{0} ; \mu\left(\mathbb{R}^{d}\right)<\infty, \quad\left\|U_{1} \mu\right\|_{\infty}<\infty\right\}
$$

where $\|\cdot\|_{\infty}$ stands for the norm of $L^{\infty}\left(\mathbb{R}^{d}\right)$.
Using these material, we define smooth measures and smooth measures in the strict sense.

Definition 1.1. (i) A positive Borel measure $\mu$ is said to be smooth $(\mu \in S)$ if $\mu$ charges no set of zero capacity and there exists an increasing sequence $\left\{F_{n}\right\}_{n \in \mathbb{N}}$ of closed sets such that $\mu\left(F_{n}\right)<\infty$ for all $n \in \mathbb{N}$ and $\lim _{n \rightarrow \infty} \operatorname{Cap}\left(K \backslash F_{n}\right)=0$ for any compact set $K$.
(ii) A positive Borel measure $\mu$ is said to be smooth in the strict sense ( $\mu \in S_{1}$ ) if there exists a sequence $\left\{E_{n}\right\}_{n \geq 1}$ of Borel sets increasing to $\mathbb{R}^{d}$ such that $1_{E_{n}} \cdot \mu \in S_{00}$ for each $n$ and $\mathbb{P}_{x}\left(\lim _{n \rightarrow \infty} T_{\mathbb{R}^{d} \backslash E_{n}}=\infty\right)=1$ for any $x \in \mathbb{R}^{d}$, where

$$
T_{\mathbb{R}^{d} \backslash E_{n}}:=\inf \left\{t>0 ; X_{t} \in \mathbb{R}^{d} \backslash E_{n}\right\} .
$$

Next we introduce additive functionals. The $\mathbb{R}$-valued stochastic process $\left\{A_{t}\right\}_{t \geq 0}$ is called positive continuous additive functional (PCAF in abbreviation) if it satisfies the following conditions:
(A-1) $A_{t}$ is $\mathscr{F}_{t}$-measurable, $\left\{\mathscr{F}_{t}\right\}$ being the minimum completed admissible filtration of the Hunt process.
(A-2) There exist a set $\Lambda \in \mathscr{F}_{\infty}$ and an exceptional set $N \subset \mathbb{R}^{d}$ such that $\mathbb{P}_{x}(\Lambda)=1$ for any $x \in \mathbb{R}^{d} \backslash N, \theta_{t} \Lambda \subset \Lambda$ for any $t>0$, and moreover for each $\omega \in \Lambda, A .(\omega)$ is right continuous and has the left limit on $[0, \zeta(\omega)), A_{0}(\omega)=0,\left|A_{t}(\omega)\right|<\infty$ for $t<\zeta(\omega), A_{t}(\omega)=A_{\zeta(\omega)}(\omega)$ for $t \geq \zeta(\omega)$ and $A_{t+s}(\omega)=A_{s}(\omega)+A_{t}\left(\theta_{s} \omega\right)$ for any $t, s \geq 0$.

If we can choose the empty set as an exceptional set in (A-2), $\left\{A_{t}\right\}$ is called PCAF in the strict sense. It is known from Theorem 5.1.3 (Theorem 5.1.7) of [13] that there is one to one correspondence between the set of smooth measures (in the strict sense) and the set of PCAFs (in the strict sense). This relation is called the Revuz correspondence: for all positive bounded Borel measurable function $f$ and $\gamma$-excessive function $h$,

$$
\int_{\mathbb{R}^{d}} f(x) h(x) \mu(d x)=\lim _{t \rightarrow 0} \frac{1}{t} \mathbb{E}_{h \cdot m}\left[\int_{0}^{t} f\left(X_{s}\right) d A_{s}^{\mu}\right]
$$

In particular, if $\mu$ is absolutely continuous with respect to the Lebesgue measure $m$ and consequently $\mu=V(x) d x$, it holds that $A_{t}^{\mu}=\int_{0}^{t} V\left(X_{s}\right) d s$.

## Chapter 2

## Perturbation of Dirichlet forms and stability of fundamental solutions

In this chapter we assume that the Hunt process $\left\{X_{t}\right\}_{t \geq 0}$ or the associated Dirichlet form $(\mathscr{E}, \mathscr{F})$ is $\alpha$-stable-like or relativistic $\alpha$-stable-like. Denote by $\left\{P_{t}\right\}$ the associated semigroup and let $p(t, x, y)$ be the transition density function of $\left\{X_{t}\right\}_{t \geq 0}$. It is well known that $\left\{P_{t}\right\}$ admits the integral kernel $p(t, x, y)$.

We consider the perturbation of Dirichlet form defined by

$$
\mathscr{E}^{\mu}(u, u):=\mathscr{E}(u, u)-\int_{\mathbb{R}^{d}} u^{2} d \mu=-(\mathscr{L} u, u)-\int_{\mathbb{R}^{d}} u^{2} d \mu .
$$

Here $\mu$ is a positive measure in the Kato class satisfying Green tightness (in abbreviation $\mu \in \mathbf{K}_{\infty}$ ). Let $\left\{P_{t}^{\mu}\right\}$ be the associated semigroup. This semigroup also admits the integral kernel $p^{\mu}(t, x, y)$ defined on $(0, \infty) \times \mathbb{R}^{d} \times \mathbb{R}^{d}$. We compare $p^{\mu}(t, x, y)$ with $p(t, x, y)$. Chen, Kim and Kumagai [7] and Chen and Kumagai [9] proved that both upper estimates and lower estimates of $p(t, x, y)$ are the same function up to positive constants. If $p^{\mu}(t, x, y)$ has the same estimate as $p(t, x, y)$ up to positive constants, we call this phenomenon stability of fundamental solution. Suppose that the measure $\mu \in \mathbf{K}_{\infty}$ is of 0-order finite energy integral, namely, $\mu$ and the Green kernel $G(x, y)$ satisfies

$$
\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} G(x, y) \mu(d x) \mu(d y)<\infty .
$$

Our goal is to prove that the subcriticality of $\mu$ is the necessary and sufficient condition for the stability of fundamental solutions.

For the construction of the necessary condition, we need to check some classes of measures. Takeda [29] showed that there are some conditions equivalent to the subcriticality of the measure $\mu$. However, this equivalence is valid for the class of conditionally Green-tight measures $\mathbf{S}_{\infty}$, which is in general a subclass of $\mathbf{K}_{\infty}$. In order to apply this argument to the class $\mathbf{K}_{\infty}$, we first prove $\mathbf{K}_{\infty}=\mathbf{S}_{\infty}$ using 3G-inequality.

For the construction of the sufficient condition, we apply Doob's $h$-transformation. Here $h(x)$ is a harmonic function of the perturbed operator $\mathscr{L}^{\mu}:=\mathscr{L}+\mu$. If $h(x)=$ $\exp (u(x))$ for $u \in \mathscr{F}_{e}$, we can apply the argument of Chen and Zhang [12], and thus construct a transformed semigroup $\left\{P_{t}^{\mu, h}\right\}$ on $L^{2}\left(h^{2} d m\right)$. We can describe the associated Dirichlet form $\left(\mathscr{E}^{\mu, h}, \mathscr{D}\left(\mathscr{E}^{\mu, h}\right)\right)$, which is equivalent to the original $(\mathscr{E}, \mathscr{F})$ and then conclude the stability of fundamental solutions. In order to apply this argument, we assume that $\mu$ is of 0 -order finite energy integral.

This chapter is organized as follows: in Section 1 we review two properties for $\alpha$ -stable-like processes and relativistic $\alpha$-stable-like processes: one is the conservativeness of processes given by Masamune and Uemura [21], and the other is the two-sided heat kernel estimates given by Chen, Kim and Kumagai [7] and Chen and Kumagai [9]. We also give the two-sided estimates for Green kernels. In Section 2 we will give the definition of some classes of smooth measures: the Kato class $\mathbf{K}$, the Green-tight Kato class $\mathbf{K}_{\infty}$ and the conditional Green-tight Kato class $\mathbf{S}_{\infty}$. In Section 3 we prove the main result following the arguments of Chen and Zhang [12] and Takeda [31].

### 2.1 Heat kernel estimates for jump Markov processes

We first define the Dirichlet forms associated with $\alpha$-stable-like process and relativistic $\alpha$-stable-like process. Let $(\mathscr{E}, \mathscr{F})$ be a jump type regular Dirichlet form on $L^{2}\left(\mathbb{R}^{d}\right)$ as follows:

$$
\begin{equation*}
\mathscr{E}(u, u)=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}(u(y)-u(x))^{2} J(x, y) d x d y, \quad \mathscr{F}=\overline{\left\{C_{c}\left(\mathbb{R}^{d}\right)\right\}^{\mathscr{E}_{1}^{1 / 2}}} \tag{2.1}
\end{equation*}
$$

where $C_{c}\left(\mathbb{R}^{d}\right)$ is the family of all continuous functions on $\mathbb{R}^{d}$ with compact supports, $\mathscr{E}_{1}(u, u)=\mathscr{E}(u, u)+\int_{\mathbb{R}^{d}} u^{2}(x) d x$ and $J(x, y)$ is a symmetric Borel function called jump measure. Here we consider the two cases:

$$
\begin{align*}
& \frac{C_{1}}{|x-y|^{d+\alpha}} \leq J(x, y) \leq \frac{C_{2}}{|x-y|^{d+\alpha}} \quad(0<\alpha<2),  \tag{2.2}\\
& \frac{C_{1} \exp \left(-m_{0}|x-y|\right)}{|x-y|^{d+\alpha}} \leq J(x, y) \leq \frac{C_{2} \exp \left(-m_{0}|x-y|\right)}{|x-y|^{d+\alpha}} \quad\left(0<\alpha<2, \quad m_{0}>0\right) \text {. } \tag{2.3}
\end{align*}
$$

In the sequel $C_{i}$ 's are unimportant positive constants varying line to line. Denote by $\left\{X_{t}\right\}_{t \geq 0}$ the associated Hunt process. If $J(x, y)$ satisfies (2.2), $\left\{X_{t}\right\}_{t \geq 0}$ is called $\alpha$ -stable-like. If $J(x, y)$ satisfies (2.3), $\left\{X_{t}\right\}_{t \geq 0}$ is called relativistic $\alpha$-stable-like.

The Hunt process $\left\{X_{t}\right\}_{t \geq 0}$ is said to be conservative if

$$
\mathbb{P}_{x}(\zeta=\infty)=1 \quad \text { q.e. } x \in \mathbb{R}^{d}
$$

where $\zeta$ is the life time of $\left\{X_{t}\right\}_{t \geq 0}$ defined by

$$
\zeta:=\inf \left\{t \geq 0, X_{t}=\Delta\right\}
$$

We first show the conservativeness of $\left\{X_{t}\right\}_{t \geq 0}$. Masamune and Uemura [21] established some sufficient conditions for conservativeness of jump Markov processes on locally compact metric spaces. If we rewrite their theorem in the framework of jump Markov processes on $\mathbb{R}^{d}$, we have the following assertion.

Theorem 2.1. Let $(\mathscr{E}, \mathscr{F})$ be a jump regular Dirichlet form defined by (2.1). The associated Hunt process $\left\{X_{t}\right\}_{t \geq 0}$ is conservative if the following two conditions are satisfied:
(i) $\sup _{x \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d} \backslash\{x\}}\left(1 \wedge|x-y|^{2}\right) J(x, y) d y<\infty$.
(ii) For any $a>0$, it holds that $e^{-a|x|} \in L^{1}\left(\mathbb{R}^{d}\right)$.

If $J(x, y)$ satisfies (2.2) or (2.3), we see that $\left\{X_{t}\right\}_{t \geq 0}$ is conservative from this theorem. Indeed, in the both cases, the upper bound of $J(x, y)$ implies

$$
\begin{gathered}
\int_{\mathbb{R}^{d} \backslash\{x\}}\left(1 \wedge|x-y|^{2}\right) J(x, y) d y \leq \int_{\mathbb{R}^{d} \backslash\{x\}}\left(1 \wedge|x-y|^{2}\right) \cdot \frac{C_{1}}{|x-y|^{d+\alpha}} d y \\
=\int_{|x-y| \leq 1} \frac{C_{1}}{|x-y|^{d+\alpha-2}} d y+\int_{|x-y| \geq 1} \frac{C_{1}}{|x-y|^{d+\alpha}} d y .
\end{gathered}
$$

Using the spherical coordinates, the last line of the above formula is equal to

$$
C_{2}\left(\int_{0}^{1} r^{1-\alpha} d r+\int_{1}^{\infty} r^{-1-\alpha} d r\right)
$$

Since $0<\alpha<2$, this value is bounded by some positive constant not depending on $x \in \mathbb{R}^{d}$. Hence we obtain (i). Moreover, we see that

$$
\int_{\mathbb{R}^{d}} \exp (-a|x|) d x=C_{3} \int_{0}^{\infty} r^{d-1} \exp (-a r) d r<\infty
$$

Thus we obtain (ii).
Next we review the two-sided heat kernel estimates. Chen and Kumagai [9] showed heat kernel estimates for $\alpha$-stable-like processes.

Theorem 2.2. Let $\left\{X_{t}\right\}_{t \geq 0}$ be an $\alpha$-stable-like process on $\mathbb{R}^{d}$. Then there exist positive constants $C_{1}$ and $C_{2}$ such that for all $t>0$ and $x, y \in \mathbb{R}^{d}$,

$$
\begin{equation*}
C_{1}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) \leq p(t, x, y) \leq C_{2}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) . \tag{2.4}
\end{equation*}
$$

Furthermore, Chen, Kim and Kumagai [7] gave heat kernel estimates for relativistic $\alpha$-stable-like processes.

Theorem 2.3. Let $\left\{X_{t}\right\}_{t \geq 0}$ be an $\alpha$-stable-like process on $\mathbb{R}^{d}$. Then the heat kernel $p(t, x, y)$ has different two-sided estimates according to the sizes of t and $|x-y|$.
(i) If $0 \leq t \leq 1$ and $|x-y| \leq 1, p(t, x, y)$ satisfies the same two-sided estimates as $\alpha$-stable-like processes:

$$
\begin{equation*}
C_{1}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) \leq p(t, x, y) \leq C_{2}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) \tag{2.5}
\end{equation*}
$$

(ii) If $0 \leq t \leq 1$ and $|x-y| \geq 1, p(t, x, y)$ satisfies

$$
\begin{equation*}
C_{1} t \frac{\exp \left(-C_{2}|x-y|\right)}{|x-y|^{d+\alpha}} \leq p(t, x, y) \leq C_{3} t \frac{\exp \left(-C_{4}|x-y|\right)}{|x-y|^{d+\alpha}} \tag{2.6}
\end{equation*}
$$

(iii) If $1 \leq t \leq|x-y|, p(t, x, y)$ satisfies

$$
\begin{equation*}
C_{1} t^{-\frac{d}{2}} \exp \left(-C_{2}|x-y|\right) \leq p(t, x, y) \leq C_{3} t^{-\frac{d}{2}} \exp \left(-C_{4}|x-y|\right) \tag{2.7}
\end{equation*}
$$

(iv) If $1 \vee|x-y| \leq t, p(t, x, y)$ satisfies

$$
\begin{equation*}
C_{1} t^{-\frac{d}{2}} \exp \left(-\frac{C_{2}|x-y|^{2}}{t}\right) \leq p(t, x, y) \leq C_{3} t^{-\frac{d}{2}} \exp \left(-\frac{C_{4}|x-y|^{2}}{t}\right) \tag{2.8}
\end{equation*}
$$

In the sequel, we assume that $\left\{X_{t}\right\}_{t \geq 0}$ is transient and consequently $\left\{X_{t}\right\}_{t \geq 0}$ admits the finite Green kernel. We obtain the two-sided estimates for Green kernels from Theorems 2.2 and 2.3.

Proposition 2.1. Let $\left\{X_{t}\right\}_{t \geq 0}$ be the transient Hunt process generated by a Dirichlet form ( $\mathscr{E}, \mathscr{F}$ ) satisfying (2.1).
(i) Suppose $J(x, y)$ satisfies (2.2). Then the Green kernel $G(x, y)$ satisfies

$$
\begin{equation*}
\frac{C_{1}}{|x-y|^{d-\alpha}} \leq G(x, y) \leq \frac{C_{2}}{|x-y|^{d-\alpha}} \tag{2.9}
\end{equation*}
$$

(ii) Suppose J $(x, y)$ satisfies (2.3). Then the Green kernel $G(x, y)$ satisfies

$$
\begin{align*}
C_{1}\left(\frac{1}{|x-y|^{d-\alpha}} \vee \frac{1}{|x-y|^{d-2}}\right) & \leq G(x, y) \\
& \leq C_{2}\left(\frac{1}{|x-y|^{d-\alpha}} \vee \frac{1}{|x-y|^{d-2}}\right) \tag{2.10}
\end{align*}
$$

Proof. (i) The transience of $\left\{X_{t}\right\}_{t \geq 0}$ implies $\alpha<d$. Since we obtain the heat kernel estimates of $\alpha$-stable-like processes in Theorem 2.2, it follows that

$$
C_{1} \int_{0}^{\infty}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) d t \leq G(x, y) \leq C_{2} \int_{0}^{\infty}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) d t
$$

We first compare $t^{-\frac{d}{\alpha}}$ with $t /|x-y|^{d+\alpha}$. Since it holds that

$$
t^{-\frac{d}{\alpha}} \leq \frac{t}{|x-y|^{d+\alpha}} \Leftrightarrow|x-y|^{d+\alpha} \leq t^{1+\frac{d}{\alpha}} \Leftrightarrow|x-y|^{\alpha} \leq t
$$

we obtain

$$
t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}= \begin{cases}t /|x-y|^{d+\alpha} & \left(t \leq|x-y|^{\alpha}\right) \\ t^{-\frac{d}{\alpha}} & \left(t \geq|x-y|^{\alpha}\right)\end{cases}
$$

This estimate implies that

$$
\begin{aligned}
\int_{0}^{\infty}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) d t & =\int_{0}^{|x-y|^{\alpha}} \frac{t}{|x-y|^{d+\alpha}} d t+\int_{|x-y|^{\alpha}}^{\infty} t^{-\frac{d}{\alpha}} d t \\
& =\frac{1}{2|x-y|^{d-\alpha}}+\left(\frac{d}{\alpha}-1\right)^{-1} \frac{1}{|x-y|^{d-\alpha}}=\frac{C_{3}}{|x-y|^{d-\alpha}} .
\end{aligned}
$$

Hence, we conclude (2.9).
(ii) The transience of the process implies $d \geq 3$. First we assume that $|x-y| \leq 1$. We have (2.5) for $0 \leq t \leq 1$ and (2.8) for $t \geq 1$. Thus it follows that

$$
\begin{align*}
& \int_{0}^{\infty} p(t, x, y) d t=\int_{0}^{1} p(t, x, y) d t+\int_{1}^{\infty} p(t, x, y) d t \\
& \leq \int_{0}^{1} C_{1}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) d t+\int_{1}^{\infty} C_{2} t^{-\frac{d}{2}} \exp \left(-\frac{C_{3}|x-y|^{2}}{t}\right) d t \tag{2.11}
\end{align*}
$$

Here the first term of (2.11) satisfies

$$
\begin{align*}
& \int_{0}^{1} C_{1}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) d t \\
\leq & \int_{0}^{\infty} C_{1}\left(t^{-\frac{d}{\alpha}} \wedge \frac{t}{|x-y|^{d+\alpha}}\right) d t \leq \frac{C_{4}}{|x-y|^{d-\alpha}} . \tag{2.12}
\end{align*}
$$

Moreover, the second term of (2.11) satisfies

$$
\begin{align*}
\int_{1}^{\infty} C_{2} t^{-\frac{d}{2}} \exp ( & \left.-\frac{C_{3}|x-y|^{2}}{t}\right) d t \leq \int_{|x-y|^{2}}^{\infty} C_{2} t^{-\frac{d}{2}} \exp \left(-\frac{C_{3}|x-y|^{2}}{t}\right) d t \\
& \leq \int_{|x-y|^{2}}^{\infty} C_{2} t^{-\frac{d}{2}} d t \leq \frac{C_{5}}{|x-y|^{d-2}} \leq \frac{C_{5}}{|x-y|^{d-\alpha}} \tag{2.13}
\end{align*}
$$

Combining (2.12) and (2.13), we obtain the upper bound:

$$
\begin{equation*}
\int_{0}^{\infty} p(t, x, y) d t \leq \frac{C_{6}}{|x-y|^{d-\alpha}} \tag{2.14}
\end{equation*}
$$

For the lower bound, we have

$$
\begin{align*}
\int_{0}^{\infty} p(t, x, y) d t & \geq \int_{0}^{|x-y|^{\alpha}} p(t, x, y) d t \\
& \geq \int_{0}^{|x-y|^{\alpha}} C_{7} \frac{t}{|x-y|^{d+\alpha}} d t=\frac{C_{8}}{|x-y|^{d-\alpha}} \tag{2.15}
\end{align*}
$$

Combining (2.14) and (2.15), we have for $|x-y| \leq 1$

$$
\begin{equation*}
\frac{C_{1}}{|x-y|^{d-\alpha}} \leq G(x, y) \leq \frac{C_{2}}{|x-y|^{d-\alpha}} . \tag{2.16}
\end{equation*}
$$

We next assume that $|x-y| \geq 1$. We obtained (2.6) for $0 \leq t \leq 1$. However, noting that $|x-y| \leq \exp (|x-y|)$, we can make the left hand side of (2.6) smaller:

$$
\begin{align*}
C_{1} t \frac{\exp \left(-C_{2}|x-y|\right)}{|x-y|^{d+\alpha}} & \geq C_{1} t \exp \left(-\left(C_{2}+d+\alpha\right)|x-y|\right) \\
& \geq C_{3} t \exp \left(-C_{4}|x-y|\right) \tag{2.17}
\end{align*}
$$

Similarly, noting that $|x-y| \geq 1$, we can make the right hand side of (2.6) larger:

$$
\begin{equation*}
C_{5} t \frac{\exp \left(-C_{6}|x-y|\right)}{|x-y|^{d+\alpha}} \leq C_{5} t \exp \left(-C_{6}|x-y|\right) \tag{2.18}
\end{equation*}
$$

Hence (2.17) and (2.18) imply

$$
C_{1} t \exp \left(-C_{2}|x-y|\right) \leq p(t, x, y) \leq C_{3} t \exp \left(-C_{4}|x-y|\right)
$$

Moreover, we obtained (2.7) for $1 \leq t \leq|x-y|$ and (2.8) for $t \geq|x-y|$. Thus it follows that

$$
\begin{align*}
\int_{0}^{\infty} p(t, x, y) d t & \leq \int_{0}^{1} C_{3} t \exp \left(-C_{4}|x-y|\right) d t+\int_{1}^{|x-y|} C_{5} t^{-\frac{d}{2}} \exp \left(-C_{6}|x-y|\right) d t \\
& +\int_{|x-y|}^{\infty} C_{7} t^{-\frac{d}{2}} \exp \left(-\frac{C_{8}|x-y|^{2}}{t}\right) d t \tag{2.19}
\end{align*}
$$

The first term of (2.19) satisfies

$$
\begin{equation*}
\int_{0}^{1} C_{3} t \exp \left(-C_{4}|x-y|\right) d t \leq C_{9} \exp \left(-C_{10}|x-y|\right) \leq \frac{C_{11}}{|x-y|^{d-2}} \tag{2.20}
\end{equation*}
$$

The second term of (2.19) satisfies

$$
\begin{align*}
\int_{1}^{|x-y|} C_{5} t^{-\frac{d}{2}} \exp \left(-C_{6}|x-y|\right) d t & \leq \int_{1}^{\infty} C_{5} t^{-\frac{d}{2}} \exp \left(-C_{6}|x-y|\right) d t \\
& \leq C_{12} \exp \left(-C_{13}|x-y|\right) \leq \frac{C_{14}}{|x-y|^{d-2}} \tag{2.21}
\end{align*}
$$

For the third term of (2.19), we substitute $s=|x-y|^{2} / t$. It follows that

$$
\begin{align*}
\int_{|x-y|}^{\infty} C_{7} t^{-\frac{d}{2}} \exp ( & \left.-\frac{C_{8}|x-y|^{2}}{t}\right) d t=\int_{0}^{|x-y|} C_{7}|x-y|^{2-d} s^{\frac{d}{2}-2} \exp \left(-C_{8} s\right) d s \\
& \leq \frac{C_{7}}{|x-y|^{d-2}} \int_{0}^{\infty} s^{\frac{d}{2}-2} \exp \left(-C_{8} s\right) d s \leq \frac{C_{15}}{|x-y|^{d-2}} \tag{2.22}
\end{align*}
$$

Thus, combining (2.20)-(2.22), we have the upper bound

$$
\begin{equation*}
p(t, x, y) \leq \frac{C_{16}}{|x-y|^{d-2}} \tag{2.23}
\end{equation*}
$$

We can also obtain the lower bound:

$$
\begin{align*}
\int_{0}^{\infty} p(t, x, y) d t & \geq \int_{|x-y|}^{\infty} p(t, x, y) d t \\
& \geq \int_{|x-y|}^{\infty} C_{17} t^{-\frac{d}{2}} \exp \left(-\frac{C_{18}|x-y|^{2}}{t}\right) d t \\
& =\int_{0}^{|x-y|} C_{17}|x-y|^{2-d} s^{\frac{d}{2}-2} \exp \left(-C_{18} s\right) d s \tag{2.24}
\end{align*}
$$

Noting that $|x-y| \geq 1$, we can make the right hand side of (2.24) smaller:

$$
\begin{align*}
& \int_{0}^{|x-y|} C_{17}|x-y|^{2-d} s^{\frac{d}{2}-2} \exp \left(-C_{18} s\right) d s \\
& \geq \int_{0}^{1} C_{17}|x-y|^{2-d} s^{\frac{d}{2}-2} \exp \left(-C_{18} s\right) d s \geq \frac{C_{19}}{|x-y|^{d-2}} . \tag{2.25}
\end{align*}
$$

We thus have for $|x-y| \geq 1$

$$
\begin{equation*}
\frac{C_{1}}{|x-y|^{d-2}} \leq G(x, y) \leq \frac{C_{2}}{|x-y|^{d-2}} \tag{2.26}
\end{equation*}
$$

from (2.23) and (2.25). On account of (2.16) and (2.26), we obtain (2.10).

We close this section introducing a certain function space. Let $\mathscr{G}$ be a function space as follows:
$\mathscr{G}=\left\{g ; g\right.$ is a positive decreasing function on $(0, \infty), \exists C_{1}, C_{2}$ s.t. $\left.C_{1} \leq \frac{g(2 r)}{g(r)} \leq C_{2}\right\}$,
where $C_{1}$ and $C_{2}$ are independent of $r>0$. The following corollary plays a crucial role in the next section.

Corollary 2.1. Let $\left\{X_{t}\right\}_{t \geq 0}$ be either an $\alpha$-stable-like process or a relativistic $\alpha$-stablelike process. Then there exists a function $g \in \mathscr{G}$ such that for all $x, y \in \mathbb{R}^{d}$ with $x \neq y$,

$$
C_{1} g(|x-y|) \leq G(x, y) \leq C_{2} g(|x-y|) .
$$

Proof. By Proposition 2.1, we have only to prove that both $g_{1}(r)=r^{\alpha-d}$ and $g_{2}(r)=$ $r^{\alpha-d} \vee r^{2-d}$ belong to the function space $\mathscr{G}$. Since $g_{1}(2 r) / g_{1}(r)=2^{\alpha-d}$, we conclude $g_{1} \in \mathscr{G}$. Similarly, we obtain

$$
\frac{g_{2}(2 r)}{g_{2}(r)}= \begin{cases}2^{\alpha-d} & (r \leq 1 / 2) \\ 2^{2-d} \cdot r^{2-\alpha} & (1 / 2 \leq r \leq 1) \\ 2^{2-d} & (r \geq 1)\end{cases}
$$

and consequently $2^{\alpha-d} \leq g_{2}(2 r) / g_{2}(r) \leq 2^{2-d}$. Thus we conclude $g_{2} \in \mathscr{G}$.

### 2.2 Properties of Kato class measures

A set $B \subset \overline{\mathbb{R}^{d}}$ is called nearly Borel measurable if for each probability measure $v$ on $\overline{\mathbb{R}^{d}}$, there exist Borel sets $B_{1}, B_{2} \in \mathscr{B}\left(\overline{\mathbb{R}^{d}}\right)$ such that $B_{1} \subset B \subset B_{2}$ and $\mathbb{P}_{v}\left(X_{t} \in B_{2} \backslash B_{1}, \exists t \geq\right.$ $0)=0$. We introduce some kinds of 'small sets' for a Markov process $\left\{X_{t}\right\}_{t \geq 0}$.
Definition 2.1. (i) $A$ set $A \subset \mathbb{R}^{d}$ is said to be polar if there exists a nearly Borel measurable set $B$ such that $A \subset B$ and $\mathbb{P}_{x}\left(\sigma_{B}=\infty\right)=1$ for all $x \in \mathbb{R}^{d}$. Here $\sigma_{B}=\inf \left\{t>0 \mid X_{t} \in B\right\}$.
(ii) The subset $A \subset \mathbb{R}^{d}$ is said to be m-polar if there exists a nearly Borel measurable set B such that $A \subset B$ and $\mathbb{P}_{m}\left(\sigma_{B}=\infty\right)=1$. Here, $\mathbb{P}_{m}(\Lambda)=\int_{\mathbb{R}^{d}} \mathbb{P}_{x}(\Lambda) m(d x)$.
Moreover, the set $A$ is $m$-polar if and only if $A$ is an exceptional set defined in Section 1.2. Recall that $\left\{X_{t}\right\}_{t \geq 0}$ admits an absolute continuous heat kernel with respect to the Lebesgue measure. It is known that the absolute continuity of heat kernel implies that the polar set or the exceptional set is empty set. Thus, when we consider a smooth measure $\mu$ or the associated PCAF $A_{t}^{\mu}$, we can strengthen them to the strict ones. We denote by $\mathbf{S}_{1}$ the family of smooth measure in the strict sense. Now we define some subclasses of $\mathbf{S}_{1}$.
Definition 2.2. A smooth measure in the strict sense $\mu$ is said to be in the Kato class ( $\mu \in \boldsymbol{K}$ in notation), if it holds that

$$
\begin{equation*}
\lim _{\beta \rightarrow \infty} \sup _{x \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d}} G_{\beta}(x, y) \mu(d y)=\lim _{\beta \rightarrow \infty} \sup _{x \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d}} \int_{0}^{\infty} e^{-\beta t} p(t, x, y) d t \mu(d y)=0 \tag{2.27}
\end{equation*}
$$

The following definition on the Green-tight smooth measures of Kato class is taken from Takeda [29].

Definition 2.3. A measure $\mu \in \boldsymbol{K}$ is said to be Green-tight ( $\mu \in \boldsymbol{K}_{\infty}$ in notation), if for any $\varepsilon>0$, there exist a compact set $K_{\varepsilon}$ and a positive constant $\delta_{\varepsilon}$ such that

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{d}} \int_{K_{\varepsilon}^{c}} G(x, y) \mu(d y) \leq \varepsilon \tag{2.28}
\end{equation*}
$$

and for any $B \subset K_{\varepsilon}$ with $\mu(B)<\delta_{\varepsilon}$, it holds that

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{d}} \int_{B} G(x, y) \mu(d y) \leq \varepsilon . \tag{2.29}
\end{equation*}
$$

The following definition on the conditionally Green-tight smooth measures of Kato class is also taken from Takeda [29].

Definition 2.4. A measure $\mu \in \boldsymbol{K}$ is said to be in the class $\boldsymbol{S}_{\infty}$, if for any $\varepsilon>0$, there exists a compact set $K_{\varepsilon}$ and a positive constant $\delta_{\varepsilon}$ such that

$$
\sup _{x, z \in \mathbb{R}^{d}} \int_{K_{\varepsilon}^{c}} \frac{G(x, y) G(y, z)}{G(x, z)} \mu(d y) \leq \varepsilon
$$

and for any $B \subset K_{\mathcal{\varepsilon}}$ with $\mu(B)<\delta_{\varepsilon}$, it holds that

$$
\sup _{x, z \in \mathbb{R}^{d}} \int_{B} \frac{G(x, y) G(y, z)}{G(x, z)} \mu(d y) \leq \varepsilon .
$$

There are different definitions for $\mathbf{K}, \mathbf{K}_{\infty}$ and $\mathbf{S}_{\infty}$ in $[6,36]$. We first make sure that these definitions are equivalent each other.

Proposition 2.2. The following assertions are equivalent each other.
(i) $\mu \in \boldsymbol{K}$. i.e. $\mu$ satisfies (2.27).
(ii) $\lim _{t \rightarrow 0} \sup _{x \in \mathbb{R}^{d}} \mathbb{E}_{x}\left[A_{t}^{\mu}\right]=0$.
(iii) $\lim _{a \rightarrow 0} \sup _{x \in \mathbb{R}^{d}} \int_{|x-y| \leq a} G(x, y) \mu(d y)=0$.

Proof. As a definition of K, [6] used (ii) and [36] used (iii) instead of (2.27). Note that

$$
\begin{equation*}
\mathbb{E}_{x}\left[A_{t}^{\mu}\right]=\int_{\mathbb{R}^{d}} \int_{0}^{t} p(s, x, y) d s \mu(d y) \tag{2.30}
\end{equation*}
$$

by [1]. Thus the formula (ii) is rewritten as follows:

$$
\begin{equation*}
\lim _{t \rightarrow 0} \sup _{x \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d}} \int_{0}^{t} p(s, x, y) d s \mu(d y)=0 . \tag{2.31}
\end{equation*}
$$

We see from Lemma 3.1 of [20] that (2.27) and (2.31) are equivalent, which implies the equivalence between (i) and (ii).

As for the equivalence between (ii) and (iii), Kuwae and Takahashi proved in Theorem 3.2 of [20] for more general Markov processes, but we give another proof here by checking some conditions in Zhao [38]. Let

$$
\begin{aligned}
\tau_{B(x, r)} & :=\inf \left\{t>0 ; X_{t} \notin B(x, r)\right\}, \\
T_{B(x, r)} & :=\inf \left\{t>0 ; X_{t} \in B(x, r)\right\} .
\end{aligned}
$$

$\tau_{B(x, r)}$ and $T_{B(x, r)}$ are called first exiting time and first hitting time of $B(x, r)$ respectively. On account of Theorem 1 of [38], it is sufficient to prove the following three formulae:

$$
\begin{align*}
\alpha_{0} & :=\sup _{t>0} \inf _{r>0} \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, r)}>t\right)<1,  \tag{2.32}\\
\beta_{0} & :=\sup _{r>0} \inf _{t>0} \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, r)}<t\right)<1,  \tag{2.33}\\
\lambda_{0} & :=\sup _{u>0} \inf _{r>0} \sup _{|x-y| \geq u} \mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right)<1 . \tag{2.34}
\end{align*}
$$

Define the function $\phi(r)$ as follows:

$$
\phi(r)= \begin{cases}r^{\alpha} & \left(\left\{X_{t}\right\}_{t \geq 0} \text { is } \alpha \text {-stable-like }\right) \\ r^{\alpha} e^{m_{0} r} & \left(\left\{X_{t}\right\}_{t \geq 0} \text { is relativistic } \alpha \text {-stable-like }\right)\end{cases}
$$

We first show that there exists a positive constant $C_{1}$ such that for all $x \in \mathbb{R}^{d}$ and $0<$ $r<1 / 4$,

$$
\begin{equation*}
\mathbb{E}_{x}\left[\tau_{B(x, r)}\right] \leq C_{1} \phi(r) . \tag{2.35}
\end{equation*}
$$

Indeed, using the Lèvy system formula and the lower bound of jump measure, we have

$$
\begin{align*}
1 & \geq \mathbb{P}_{x}\left(X_{\tau_{B(x, r)}} \notin \overline{B(x, 2 r)}\right)=\mathbb{E}_{x}\left[\int_{0}^{\tau_{B(x, r)}} \int_{\overline{B(x, 2 r)^{c}}} J\left(X_{s}, u\right) d u d s\right] \\
& \geq \mathbb{E}_{x}\left[\int_{0}^{\tau_{B(x, r)}} \int_{\overline{B(x, 2 r)^{c}}} \frac{C_{2}}{\left|X_{s}-u\right|^{d} \phi\left(\left|X_{s}-u\right|\right)} d u d s\right] . \tag{2.36}
\end{align*}
$$

Note that $X_{s} \in B(x, r)$ for $0 \leq s<\tau_{B(x, r)}$ and for fixed $w \in B(x, r),\{u:|w-u| \geq$ $3 r\} \subset \overline{B(x, 2 r)}^{c}$. Applying this to the right hand side of (2.36) and using the spherical coordinates, we obtain

$$
\begin{align*}
& \mathbb{E}_{x}\left[\int_{0}^{\tau_{B(x, r)}} \int_{\overline{B(x, 2 r)}} c\right.\left.\frac{C_{2}}{\left|X_{s}-u\right|^{d} \phi\left(\left|X_{s}-u\right|\right)} d u d s\right] \\
& \geq C_{3} \int_{3 r}^{\infty} \frac{1}{\rho \phi(\rho)} d \rho \mathbb{E}_{x}\left[\tau_{B(x, r)}\right]  \tag{2.37}\\
& \geq C_{4} \int_{3 r}^{1} \rho^{-\alpha-1} d \rho \mathbb{E}_{x}\left[\tau_{B(x, r)}\right]=C_{5}\left((3 r)^{-\alpha}-1\right) \mathbb{E}_{x}\left[\tau_{B(x, r)}\right]
\end{align*}
$$

Since we assumed that $0<r<1 / 4$, it follows that

$$
(3 r)^{-\alpha}-1=(3 r)^{-\alpha}\left(1-(3 r)^{\alpha}\right) \geq C_{6} r^{-\alpha} .
$$

Hence the right hand side of (2.37) is estimated as follows:

$$
\begin{equation*}
C_{5}\left((3 r)^{-\alpha}-1\right) \mathbb{E}_{x}\left[\tau_{B(x, r)}\right] \geq C_{6} r^{-\alpha} \mathbb{E}_{x}\left[\tau_{B(x, r)}\right] \geq \frac{C_{6}}{\phi(r)} \mathbb{E}_{x}\left[\tau_{B(x, r)}\right] \tag{2.38}
\end{equation*}
$$

Thus (2.36)-(2.38) imply $1 \geq C_{6} / \phi(r) \cdot \mathbb{E}_{x}\left[\tau_{B(x, r)}\right]$ and we have (2.35).This is an extension of Theorem 5.1 in [9]. It is clear that (2.35) implies

$$
\mathbb{P}_{x}\left(\tau_{B(x, r)}>t\right) \leq \frac{C_{8} \phi(r)}{t}
$$

and we obtain (2.32) with $\alpha_{0}=0$.
Applying Proposition 4.9 of Chen and Kumagai [10], we see that for arbitrary $\varepsilon>0$ there exists $\gamma_{\varepsilon}>0$ such that for $0<r<1$ and $x \in \mathbb{R}^{d}$

$$
\mathbb{P}_{x}\left(\tau_{B(x, r)}<\gamma_{\varepsilon} \phi(r)\right) \leq \varepsilon .
$$

Thus, for $0<r<1$ we obtain

$$
\inf _{t>0} \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, r)}<t\right) \leq \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, r)}<\gamma_{\varepsilon} \phi(r)\right) \leq \varepsilon .
$$

For $r \geq 1$, we obtain

$$
\begin{aligned}
\inf _{t>0} \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, r)}<t\right) & \leq \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, r)}<\gamma_{\varepsilon} \phi(1 / 2)\right) \\
& \leq \sup _{x \in \mathbb{R}^{d}} \mathbb{P}_{x}\left(\tau_{B(x, 1 / 2)}<\gamma_{\varepsilon} \phi(1 / 2)\right) \leq \varepsilon .
\end{aligned}
$$

Hence, we obtain (2.33) with $\beta_{0}=0$.
We can prove (2.34) in the same way as is used in Lemma 5 of [38], which deals with Lévy processes. Corollary 2.1 implies that there exist positive constants $C_{1}, C_{2}$ and $g \in \mathscr{G}$ such that

$$
C_{1} g(|x-y|) \leq G(x, y) \leq C_{2} g(|x-y|) .
$$

Fix $u>0$ and $b \geq 1$. If $|x-y| \geq u, 0<r \leq u /(2 b+1)$ and $|z-x| \leq r$, it follows that $|y-z| \geq 2 b r$ and consequently

$$
C_{2} g(2 b r) \geq C_{2} g(|y-z|) \geq G(y, z) .
$$

Thus, we have

$$
\begin{align*}
C_{3} g(2 b r) r^{d} & \geq \int_{B(x, r)} G(y, z) d z \geq \mathbb{E}_{y}\left[\int_{T_{B(x, r)}}^{\infty} 1_{B(x, r)}\left(X_{t}\right) d t\right] \\
& =\mathbb{E}_{y}\left[\mathbb{E}_{X_{T_{B(x, r)}}}\left[\int_{0}^{\infty} 1_{B(x, r)}\left(X_{t}\right) d t\right]: T_{B(x, r)}<\infty\right], \tag{2.39}
\end{align*}
$$

where we used the strong Markov property in the last equality. Note that $X_{T_{B(x, r)}} \in$ $\overline{B(x, r)}$ and

$$
\mathbb{E}_{w}\left[\int_{0}^{\infty} 1_{B(x, r)}\left(X_{t}\right) d t\right]=\int_{B(x, r)} G(w, z) d z .
$$

Hence the right hand side of (2.39) satisfies

$$
\begin{align*}
& \mathbb{E}_{y}\left[\mathbb{E}_{X_{T_{B(x, r)}}}\left[\int_{0}^{\infty} 1_{B(x, r)}\left(X_{t}\right) d t\right]: T_{B(x, r)}<\infty\right] \\
& \geq \mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) \cdot \inf _{|w-x| \leq r} \int_{B(x, r)} G(w, z) d z \\
& \geq \mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) \cdot \inf _{|w-x| \leq r} \int_{B(x, r)} C_{1} g(|w-z|) d z . \tag{2.40}
\end{align*}
$$

Moreover, $|w-x| \leq r$ and $z \in B(x, r)$ imply $|w-z| \leq 2 r$, and using the monotone decreasing property of $g$, we obtain

$$
\begin{equation*}
\mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) \cdot \inf _{|w-x| \leq r} \int_{B(x, r)} C_{1} g(|w-z|) d z \geq C_{4} g(2 r) r^{d} \mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) \tag{2.41}
\end{equation*}
$$

From (2.39)-(2.41), we see that

$$
\mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) \leq \frac{C_{3} g(2 b r)}{C_{4} g(2 r)}
$$

for $|x-y| \geq u, 0<r \leq u /(2 b+1)$ and $u>0$. We thus see that

$$
\begin{aligned}
\inf _{r>0} \sup _{|x-y| \geq u} \mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) & \leq \inf _{0<r \leq u /(2 b+1)} \sup _{|x-y| \geq u} \mathbb{P}_{y}\left(T_{B(x, r)}<\infty\right) \\
& \leq \inf _{0<r \leq u /(2 b+1)} \sup _{|x-y| \geq u} \frac{C_{3} g(2 b r)}{C_{4} g(2 r)} \\
& =\inf _{0<r \leq u /(2 b+1)} \frac{C_{3} g(2 b r)}{C_{4} g(2 r)} \leq \limsup _{r \rightarrow 0} \frac{C_{3} g(2 b r)}{C_{4} g(2 r)}
\end{aligned}
$$

For a sufficiently large $b$, the right hand side of the above formula is smaller than 1 . Hence, we obtain (2.34).

In the sequel, we assume $\mu \in \mathbf{K}$. The following proposition says that the two definitions on the Green-tight smooth measures of Kato class from [6] and [29] coincide with each other.

Proposition 2.3. For $\mu \in \boldsymbol{K}$, the following assertions are equivalent each other.
(i) $\mu \in \boldsymbol{K}_{\infty}$. i.e. $\mu$ satisfies (2.28)-(2.29).
(ii) For any $\varepsilon>0$ there exist a set $F_{\varepsilon}$ of $\mu$-finite measure and a positive constant $\tilde{\delta}_{\varepsilon}$ such that

$$
\sup _{x \in \mathbb{R}^{d}} \int_{F_{\varepsilon}^{c}} G(x, y) \mu(d y) \leq \varepsilon
$$

and for any $B \subset F_{\varepsilon}$ with $\mu(B)<\tilde{\delta}_{\varepsilon}$,

$$
\sup _{x \in \mathbb{R}^{d}} \int_{B} G(x, y) \mu(d y) \leq \varepsilon .
$$

(iii) It holds that

$$
\lim _{r \rightarrow \infty} \sup _{x \in \mathbb{R}^{d}} \int_{|y|>r} G(x, y) \mu(d y)=0 .
$$

Proof. We assume (i) and let $K_{\varepsilon}$ and $\delta_{\varepsilon}$ be a compact set and a positive constant in (2.28)-(2.29) respectively. Since $\mu \in \mathbf{K}$ and $g(|x-y|) \rightarrow \infty$ as $|x-y| \rightarrow 0$, (iii) of Proposition 2.2 implies that $\mu(B(x, a)) \rightarrow 0$ as $a \rightarrow 0$ uniformly in $x \in \mathbb{R}^{d}$. In particular, there exist positive constants $a_{0}$ and $C_{0}$ such that $\mu\left(B\left(x_{0}, a_{0}\right)\right) \leq C_{0}$ for all $x \in \mathbb{R}^{d}$. Since $K_{\varepsilon}$ is compact, $K_{\varepsilon}$ is covered by finite subset of $\left\{B\left(x, a_{0}\right)\right\}_{x \in \mathbb{R}^{d}}$ and thus we obtain $\mu\left(K_{\varepsilon}\right)<\infty$. Hence (ii) follows for $F_{\varepsilon}=K_{\varepsilon}$ and $\tilde{\delta}_{\varepsilon}=\delta_{\varepsilon}$.

We next show (ii) implies (i). We follow the proof of Theorem 2.1 (3) in [6]. Let $F_{\varepsilon}$ and $\tilde{\delta}_{\varepsilon}$ be a set of $\mu$-finite measure and a positive constant satisfying (ii). Since

$$
\mu\left(\overline{B(0, R)}^{c} \cap F_{\varepsilon}\right) \rightarrow 0 \quad(R \rightarrow \infty)
$$

there exists a positive constant $R_{\varepsilon}$ such that

$$
\left.\sup _{x \in \mathbb{R}^{d}} \int_{\overline{B\left(0, R_{\varepsilon}\right)}}{ }^{c} \cap F_{\varepsilon}\right] .
$$

We thus obtain

$$
\begin{aligned}
& \sup _{x \in \mathbb{R}^{d}} \int_{\overline{B\left(0, R_{\varepsilon}\right)}} G(x, y) \mu(d y) \\
& \leq \sup _{x \in \mathbb{R}^{d}}\left(\int_{F_{\varepsilon}^{c}} G(x, y) \mu(d y)+\int_{\overline{B\left(0, R_{\varepsilon}\right)}{ }^{c} \cap F_{\varepsilon}} G(x, y) \mu(d y)\right) \leq 2 \varepsilon .
\end{aligned}
$$

It holds that for $B \subset \overline{B\left(0, R_{\varepsilon}\right)}$ with $\mu(B)<\tilde{\delta}_{\varepsilon}$

$$
\sup _{x \in \mathbb{R}^{d}} \int_{B} G(x, y) \mu(d y) \leq \sup _{x \in \mathbb{R}^{d}}\left(\int_{B \cap F_{\varepsilon}} G(x, y) \mu(d y)+\int_{F_{\varepsilon}^{c}} G(x, y) \mu(d y)\right) \leq 2 \varepsilon .
$$

Hence, (i) follows for $K_{\varepsilon}=\overline{B\left(0, R_{\frac{\varepsilon}{2}}\right)}$ and $\delta_{\varepsilon}=\tilde{\delta}_{\frac{\varepsilon}{2}}$.
We easily see that (iii) follows from (i) by choosing the empty set as $B$ and $\overline{B(0, r)}$ as $K$ respectively.

If (iii) is valid, it follows that for arbitrary $\varepsilon>0$ there exists a sufficient large $r_{\varepsilon}>0$ such that

$$
\sup _{x \in \mathbb{R}^{d}} \int_{|y|>r_{\varepsilon}} G(x, y) \mu(d y) \leq \frac{\varepsilon}{2} .
$$

Set $K_{\varepsilon}=\overline{B\left(0, r_{\varepsilon}\right)}$. Since $\mu \in \mathbf{K}$, (iii) of Proposition 2.2 implies that there exists a sufficient small positive constant $a_{\varepsilon}$ such that

$$
\sup _{x \in \mathbb{R}^{d}} \int_{|x-y| \leq a_{\varepsilon}} G(x, y) \mu(d y) \leq \frac{\varepsilon}{2} .
$$

Thus, it holds that for a measurable set $A \subset K_{\varepsilon}$

$$
\begin{align*}
\int_{A} G(x, y) \mu(d y) & \leq \int_{B\left(x, a_{\varepsilon}\right)} G(x, y) \mu(d y)+\int_{A \cap B\left(x, a_{\varepsilon}\right)^{c}} G(x, y) \mu(d y) \\
& \leq \frac{\varepsilon}{2}+C_{1} g\left(a_{\varepsilon}\right) \mu(A) . \tag{2.42}
\end{align*}
$$

If we choose a sufficient small positive constant $\delta_{\varepsilon}$, the second term of (2.42) can be smaller than $\varepsilon / 2$ for any set $A$ with $\mu(A)<\delta_{\varepsilon}$. Hence we have (i).

Denote by $\hat{\mathbf{K}}_{\infty}$ the class of measures satisfying the conditions in (ii) of Proposition 2.3. This is the definition of Green-tight smooth measures of Kato class taken from [6]. Next we compare the definition of $\mathbf{S}_{\infty}$ with that in [6]. The following definition on the conditionally Green-tight smooth measures of Kato class is taken from Chen [6].

Definition 2.5. The measure $\mu$ belongs to the class $\hat{\boldsymbol{S}}_{\infty}$ if for any $\varepsilon>0$ there exist a set $F_{\varepsilon}$ of $\mu$-finite measure and $\tilde{\delta}_{\varepsilon}>0$ such that

$$
\sup _{x, z \in \mathbb{R}^{d} \times \mathbb{R}^{d} \backslash \Delta} \int_{F_{\varepsilon}^{c}} \frac{G(x, y) G(y, z)}{G(x, z)} \mu(d y) \leq \varepsilon,
$$

$\Delta=\left\{(x, x) ; x \in \mathbb{R}^{d}\right\}$ and for any $B \subset F_{\varepsilon}$ with $\mu(B)<\tilde{\delta}_{\varepsilon}$

$$
\sup _{x, z \in \mathbb{R}^{d} \times \mathbb{R}^{d} \backslash \Delta} \int_{B} \frac{G(x, y) G(y, z)}{G(x, z)} \mu(d y) \leq \varepsilon
$$

The following proposition can be proved in the same way as is used in the proof of (i) $\Rightarrow$ (ii) in Proposition 2.3.

Proposition 2.4. It holds that $\boldsymbol{S}_{\infty} \subset \hat{\boldsymbol{S}}_{\infty}$.
Proof. Let $\mu$ is in the class $\mathbf{S}_{\infty}$. Since $\mu \in \mathbf{K}$ from Definition 2.4, we can refrain the argument of the beginning of the proof for Proposition 2.3. In particular, we can choose $K_{\varepsilon}$ and $\delta_{\varepsilon}$ in Definition 2.4 as $F_{\varepsilon}$ and $\tilde{\delta}_{\varepsilon}$ respectively.

The next theorem is proved by means of $3 G$-theorem.
Theorem 2.4. It holds that $\boldsymbol{K}_{\infty}=\boldsymbol{S}_{\infty}=\hat{\boldsymbol{S}}_{\infty}$.
Proof. Let $g$ be a function in Corollary 2.1. Note that either $|x-y| \geq|x-z| / 2$ or $|y-z| \geq|x-z| / 2$ holds. If $|x-y| \geq|x-z| / 2$, it follows that

$$
\begin{align*}
\frac{G(x, y) G(y, z)}{G(x, z)} & \leq \frac{C_{2} g(|x-y|)}{C_{1} g(|x-z|)} G(y, z) \\
& \leq \frac{C_{2} g(|x-z| / 2)}{C_{1} g(|x-z|)} G(y, z) \leq C_{3} G(y, z)=C_{3} G(z, y) \tag{2.43}
\end{align*}
$$

by Proposition 2.1. If $|y-z| \geq|x-z| / 2$ we can similarly obtain

$$
\begin{equation*}
\frac{G(x, y) G(y, z)}{G(x, z)} \leq C_{4} G(x, y) . \tag{2.44}
\end{equation*}
$$

Combining (2.43) and (2.44) we see that

$$
\frac{G(x, y) G(y, z)}{G(x, z)} \leq C_{5}(G(x, y)+G(z, y))
$$

and thus $\mathbf{K}_{\infty} \subset \mathbf{S}_{\infty}$. Moreover, we see $\hat{\mathbf{S}}_{\infty} \subset \hat{\mathbf{K}}_{\infty}$ from Corollary 3.1 of [11] and p. 4663 of [6]. It is also proved that $\hat{\mathbf{K}}_{\infty} \subset \mathbf{K}$ by Proposition 2.3 of [6]. Thus we obtain $\hat{\mathbf{S}}_{\infty} \subset \mathbf{K}_{\infty}$. Combining with Proposition 2.4, we obtain $\mathbf{K}_{\infty} \subset \mathbf{S}_{\infty} \subset \hat{\mathbf{S}}_{\infty} \subset \mathbf{K}_{\infty}$ and this is the desired assertion.

We close this section introducing an important property of the measure in $\mathbf{K}_{\infty}$.

Corollary 2.2. For $\mu \in \boldsymbol{K}_{\infty}$,

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{d}} \mathbb{E}_{x}\left[A_{\infty}^{\mu}\right]<\infty . \tag{2.45}
\end{equation*}
$$

Proof. The left hand side of (2.45) is rewritten as follows:

$$
\mathbb{E}_{x}\left[A_{\infty}^{\mu}\right]=\int_{\mathbb{R}^{d}} G(x, y) \mu(d y)
$$

We see that the right hand side of the above formula is equal to

$$
\begin{equation*}
\int_{|x-y| \leq a_{\varepsilon}} G(x, y) \mu(d y)+\int_{K_{\varepsilon}^{c}} G(x, y) \mu(d y)+\int_{K_{\varepsilon} \cap|x-y|>a_{\varepsilon}} G(x, y) \mu(d y), \tag{2.46}
\end{equation*}
$$

where $a_{\varepsilon}$ is a positive constant such that the first term of (2.46) is smaller than $\varepsilon$ uniformly in $x \in \mathbb{R}^{d}$ and $K_{\varepsilon}$ is a compact set taken from Definition 2.3. Thus the first and the second terms of (2.46) are uniformly bounded. Since $K_{\mathcal{\varepsilon}}$ is of finite $\mu$-measure from Proposition 2.3 and $G(x, y)$ is bounded on $|x-y| \geq a_{\varepsilon}$, the third term of (2.46) is also uniformly bounded. Hence, we have the desired result.

Remark 2.1. In [6], the property (2.45) is proved for the wider class of measures $\boldsymbol{K}_{1}$. $\mu$ is said to be in the class $\boldsymbol{K}_{1}$ if there is a Borel set $F$ of finite $\mu$-measure and a constant $\tilde{\delta}$ such that

$$
\sup _{B \subset F: \mu(B)<\tilde{\delta}} \sup _{x \in \mathbb{R}^{d}} \int_{F^{c} \cup B} G(x, y) \mu(d y)<1
$$

If $\mu \in \boldsymbol{K}_{\infty}$, we can make the left hand side arbitrarily small.

### 2.3 Stability of fundamental solutions

In this section we assume $\mu \in \mathbf{K}_{\infty}$ is positive and satisfies

$$
\begin{equation*}
\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} G(x, y) \mu(d x) \mu(d y)<\infty \tag{2.47}
\end{equation*}
$$

$\mu$ is said to be of 0 -order finite energy integral if $\mu$ satisfies (2.47). Consider the Schrödinger form $\left(\mathscr{E}^{\mu}, \mathscr{F}\right)$ defined by

$$
\begin{equation*}
\mathscr{E}^{\mu}(u, u):=\mathscr{E}(u, u)-\int_{\mathbb{R}^{d}} u^{2} d \mu \tag{2.48}
\end{equation*}
$$

Denote by $\left\{P_{t}^{\mu}\right\}_{t \geq 0}$ the corresponding semigroup. It is known that $P_{t}^{\mu}$ is written by

$$
P_{t}^{\mu} f(x)=\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right) f\left(X_{t}\right)\right]
$$

where $A_{t}^{\mu}$ is a positive continuous additive functional in the Revuz correspondence with $\mu$. Following [1, 2], we see that $\left\{P_{t}^{\mu}\right\}_{t \geq 0}$ admits the integral kernel $p^{\mu}(t, x, y)$ defined on $(0, \infty) \times \mathbb{R}^{d} \times \mathbb{R}^{d}$. Our goal is to prove the following theorem:

Theorem 2.5. Let $(\mathscr{E}, \mathscr{F})$ be a transient Dirichlet form associated with $\alpha$-stable-like process or relativistic $\alpha$-stable-like process. Suppose $\mu$ is a Green-tight measure of 0 -order finite energy integral. Then the stability of fundamental solution is valid if and only if $\mu$ satisfies

$$
\begin{equation*}
\inf \left\{\mathscr{E}(u, u) ; u \in \mathscr{F}, \int_{\mathbb{R}^{d}} u^{2} d \mu=1\right\}>1 \tag{2.49}
\end{equation*}
$$

For the proof of Theorem 2.5, the following proposition plays a crucial role.
Proposition 2.5. For $\mu \in \boldsymbol{K}_{\infty}$, the following assertions are equivalent.
(i) $G^{\mu}(x, y):=\int_{0}^{\infty} p^{\mu}(t, x, y) d t<\infty \quad$ for $x \neq y$;
(ii) $\inf \left\{\mathscr{E}(u, u) ; u \in \mathscr{F}, \int_{\mathbb{R}^{d}} u^{2} d \mu=1\right\}>1$;
(iii) $\sup \mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right)\right]<\infty$.

$$
x \in \mathbb{R}^{d}
$$

Proof. This proposition is proved for $\mu \in \mathbf{S}_{\infty}$ in Theorem 2.4 and Theorem 3.9 of [29]. Since we see that $\mathbf{K}_{\infty}=\mathbf{S}_{\infty}$ from Theorem 2.4, we obtain the desired result.

We can prove the 'only if' part of Theorem 2.5 from equivalence between (i) and (ii) in Proposition 2.5. Suppose the stability of fundamental solution holds. Using the Green kernel estimates similarly obtained as in Proposition 2.1, we see that $G^{\mu}(x, y)<$ $\infty$. This is equivalent to (2.49).
Before proving the 'if' part, we introduce the definition of gaugeability.
Definition 2.6. The Green-tight measure $\mu$ is called gaugeable if $A_{t}^{\mu}$ satisfies (iii) of Proposition 2.5.

Now we prove the 'if' part of Theorem 2.5. The equation (2.49) and Proposition 2.5 imply that

$$
1 \leq h(x):=\mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right)\right] \leq C_{1}<\infty .
$$

In order to apply Theorem 3.4 in [12], we need to show that there exists $u \in \mathscr{F}_{e}$ such that $h(x)=\exp (u(x))$, where $\mathscr{F}_{e}$ is the extended Dirichlet space, namely the closure of $\mathscr{F}$ with respect to the $\mathscr{E}^{1 / 2}$-norm.

Lemma 2.1. For $\mu \in \boldsymbol{K}_{\infty}$, define

$$
G \mu(x)=\int_{\mathbb{R}^{d}} G(x, y) \mu(d y) .
$$

Under assumption (2.47), $G \mu \in \mathscr{F}_{e}$.

Proof. This lemma is an extension of Lemma 3.1 in [31]. Let $\mu \in \mathbf{K}_{\infty}$. Then it holds that

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} u^{2} d \mu \leq\|G \mu\|_{\infty} \mathscr{E}(u, u) \tag{2.50}
\end{equation*}
$$

for $u \in \mathscr{F}_{e}$. This is the modification of Theorem 3.1 in Stollmann and Voigt [27]. Define $\mu_{K}(\cdot)=\mu(K \cap \cdot)$ for a set $K$ of $\mu$-finite measure. Applying (2.50), we have

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} \psi d \mu_{K} & \leq(\mu(K))^{1 / 2}\left(\int_{\mathbb{R}^{d}} \psi^{2} d \mu_{K}\right)^{1 / 2} \\
& \leq(\mu(K))^{1 / 2}\left\|G \mu_{K}\right\|_{\infty}^{1 / 2} \mathscr{E}(\psi, \psi)^{1 / 2}
\end{aligned}
$$

By (2.45), we see $\left\|G \mu_{K}\right\|_{\infty} \leq\|G \mu\|_{\infty}=\sup _{x \in \mathbb{R}^{d}} \mathbb{E}_{x}\left[A_{\infty}^{\mu}\right]<\infty$ and consequently $\mu_{K}$ is of 0 -order finite energy integral in the sense of [13]. We thus have

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} \psi d \mu_{K} & \leq \mathscr{E}\left(G \mu_{K}, G \mu_{K}\right)^{1 / 2} \mathscr{E}(\psi, \psi)^{1 / 2} \\
& \leq\left(\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} G(x, y) d \mu_{K}(x) d \mu_{K}(y)\right)^{1 / 2} \mathscr{E}(\psi, \psi)^{1 / 2} .
\end{aligned}
$$

We see that $\mu$ is also of 0 -order finite energy integral and $G \mu \in \mathscr{F}_{e}$ by letting $K$ to $\mathbb{R}^{d}$.

Lemma 2.2. Suppose $\mu \in \boldsymbol{K}_{\infty}$ is gaugeable. Then it holds that

$$
h(x)=1+G(h \mu)(x) .
$$

Proof. This is an extension of Lemma 3.2 of [31]. Let $\left\{\mathscr{M}_{t}\right\}_{t \geq 0}$ be the filtration equipped with the Hunt process $\left\{X_{t}\right\}_{t \geq 0}$. Set $\bar{M}_{t}=\mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right) \mid \mathscr{M}_{t}\right]$. By the Markov property, we have

$$
h\left(X_{t}\right)=\mathbb{E}_{X_{t}}\left[\exp \left(A_{\infty}^{\mu}\right)\right]=\mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu} \circ \theta_{t}\right) \mid \mathscr{M}_{t}\right]
$$

$\mathscr{M}_{t}$-measurability of $A_{t}^{\mu}$ and the property of additive functionals imply

$$
\begin{aligned}
\mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu} \circ \theta_{t}\right) \mid \mathscr{M}_{t}\right] & =\exp \left(-A_{t}^{\mu}\right) \mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}+A_{\infty}^{\mu} \circ \theta_{t}\right) \mid \mathscr{M}_{t}\right] \\
& =\exp \left(-A_{t}^{\mu}\right) \mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right) \mid \mathscr{M}_{t}\right]=\exp \left(-A_{t}^{\mu}\right) \bar{M}_{t}
\end{aligned}
$$

where $\theta_{t}$ is the shift operator satisfying $X_{t+s}=X_{s} \circ \theta_{t}$ for all $s \geq 0$. Hence, we have

$$
\begin{equation*}
\mathbb{E}_{x}\left[\int_{0}^{t} h\left(X_{s}\right) d A_{s}^{\mu}\right]=\mathbb{E}_{x}\left[\int_{0}^{t} \exp \left(-A_{s}^{\mu}\right) \bar{M}_{s} d A_{s}^{\mu}\right] \tag{2.51}
\end{equation*}
$$

Noting that

$$
\exp \left(-A_{s}^{\mu}\right) d A_{s}^{\mu}=-d\left(\exp \left(-A_{s}^{\mu}\right)\right)
$$

and using the integral by parts, the right hand side of (2.51) is equal to

$$
\mathbb{E}_{x}\left[\bar{M}_{0}\right]-\mathbb{E}_{x}\left[\exp \left(-A_{t}^{\mu}\right) \bar{M}_{t}\right]+\mathbb{E}_{x}\left[\int_{0}^{t} \exp \left(-A_{s}^{\mu}\right) d \bar{M}_{s}\right]=h(x)-\mathbb{E}_{x}\left[h\left(X_{t}\right)\right]
$$

where we use the martingale property of $\left\{\bar{M}_{t}\right\}_{t \geq 0}$ in the last equality. Thus, we have

$$
\begin{equation*}
\mathbb{E}_{x}\left[\int_{0}^{t} h\left(X_{s}\right) d A_{s}^{\mu}\right]=h(x)-\mathbb{E}_{x}\left[h\left(X_{t}\right)\right] \tag{2.52}
\end{equation*}
$$

and note that

$$
\lim _{t \rightarrow \infty} h\left(X_{t}\right)=\lim _{t \rightarrow \infty} \exp \left(-A_{t}^{\mu}\right) \bar{M}_{t}=\exp \left(-A_{\infty}^{\mu}\right) \exp \left(A_{\infty}^{\mu}\right)=1 .
$$

We then have the desired result by letting $t \rightarrow \infty$ in (2.52).
Recall that $\mathscr{F}_{e}$ is a Hilbert space with inner product $\mathscr{E}^{1 / 2}$ if $\mathscr{E}$ is transient. Since $h$ is a positive bounded function and $G \mu \in \mathscr{F}_{e}, G(h \mu) \in \mathscr{F}_{e}$. Moreover, define

$$
u(x):=\log h(x)=\log (1+G(h \mu)(x)) .
$$

Since $u(x)$ is a normal contraction of $G(h \mu)(x), u \in \mathscr{F}_{e}$ and hence we conclude that $h(x)=\exp (u(x))$ for $u \in \mathscr{F}_{e}$.

Next, consider Fukushima's decomposition of $G(h \mu)$ :

$$
\begin{equation*}
G(h \mu)\left(X_{t}\right)-G(h \mu)\left(X_{0}\right)=M_{t}^{[G(h \mu)]}+N_{t}^{[G(h \mu)]}, \tag{2.53}
\end{equation*}
$$

where $M_{t}^{[G(h \mu)]}$ is a martingale additive functional of finite energy and $N_{t}^{[G(h \mu)]}$ is a continuous additive functional of zero energy. Since the left hand side of (2.53) equals $h\left(X_{t}\right)-h\left(X_{0}\right)$ by Lemma 2.2, $M_{t}^{[G(h \mu)]}$ equals $M_{t}^{[h]}$. Moreover we see from Lemma 5.4.1 of [13] that

$$
N_{t}^{[G(h \mu)]}=-\int_{0}^{t} h\left(X_{s}\right) d A_{s}^{\mu} .
$$

We thus have

$$
h\left(X_{t}\right)-h\left(X_{0}\right)=M_{t}^{[h]}-\int_{0}^{t} h\left(X_{s}\right) d A_{s}^{\mu} .
$$

Furthermore, we define a martingale by

$$
M_{t}=\int_{0}^{t} \frac{1}{h\left(X_{s-}\right)} d M_{s}^{[h]}
$$

and denote by $L_{t}$ the unique solution of Doleans-Dade equation: $L_{t}=1+\int_{0}^{t} L_{s-} d M_{s}$. From Theorem 9.39 of [15], we know that $L_{t}$ is expressed as

$$
L_{t}=\exp \left(M_{t}-\frac{1}{2}\left\langle M^{c}\right\rangle_{t}\right) \prod_{0<s \leq t}\left(1+\Delta M_{s}\right) \exp \left(-\Delta M_{s}\right)
$$

where $M^{c}$ is the continuous part of martingale $M$ and $\left\langle M^{c}\right\rangle$ is the quadratic variation of $M^{c}$. Noting that

$$
\Delta M_{s}=M_{s}-M_{s-}=\frac{1}{h\left(X_{s-}\right)}\left(h\left(X_{s}\right)-h\left(X_{s-}\right)\right)=\frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)}-1,
$$

we obtain

$$
L_{t}=\exp \left(M_{t}-\frac{1}{2}\left\langle M^{c}\right\rangle_{t}\right) \prod_{0<s \leq t} \frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)} \exp \left(1-\frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)}\right)
$$

In order to calculate $L_{t}$, we apply Itô formula to the semimartingale $h\left(X_{t}\right)$ and the function $\log x$. Thus, we have

$$
\begin{align*}
& \log h\left(X_{t}\right)=\log h\left(X_{0}\right)+\int_{0}^{t} \frac{1}{h\left(X_{s-}\right)}\left(d M_{s}^{[h]}-h\left(X_{s}\right) d A_{s}^{\mu}\right) \\
& +\sum_{0<s \leq t}\left(\log h\left(X_{s}\right)-\log h\left(X_{s-}\right)-\frac{1}{h\left(X_{s-}\right)} \Delta M_{s}^{[h]}\right)+\frac{1}{2} \int_{0}^{t} \frac{-1}{h\left(X_{s-}\right)^{2}} d\left\langle h(X)^{c}\right\rangle_{s} . \tag{2.54}
\end{align*}
$$

The right hand side of (2.54) is equal to

$$
\begin{aligned}
\log h\left(X_{0}\right) & +M_{t}-\int_{0}^{t} \frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)} d A_{s}^{\mu} \\
& +\sum_{0<s \leq t}\left\{\log \frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)}+\left(1-\frac{h\left(X_{S}\right)}{h\left(X_{s-}\right)}\right)\right\}-\frac{1}{2} \int_{0}^{t} \frac{1}{h\left(X_{s-}\right)^{2}} d\left\langle M^{[h], c}\right\rangle_{s}
\end{aligned}
$$

Noting that the set $\left\{s ; X_{s} \neq X_{s-}\right\}$ is at most countable, the above formula is equal to

$$
\log h\left(X_{0}\right)+M_{t}-A_{t}^{\mu}+\sum_{0<s \leq t}\left\{\log \frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)}+\left(1-\frac{h\left(X_{s}\right)}{h\left(X_{s-}\right)}\right)\right\}-\frac{1}{2}\left\langle M^{c}\right\rangle_{t}
$$

Hence we obtain

$$
L_{t}=\frac{h\left(X_{t}\right)}{h\left(X_{0}\right)} \exp \left(A_{t}^{\mu}\right)
$$

We consider the transformed semigroup $\left\{P_{t}^{\mu, h}\right\}_{t \geq 0}$ by $L_{t}$,

$$
P_{t}^{\mu, h} f(x)=\mathbb{E}_{x}\left[L_{t} f\left(X_{t}\right)\right]=\frac{1}{h(x)} \mathbb{E}_{x}\left[h\left(X_{t}\right) \exp \left(A_{t}^{\mu}\right) f\left(X_{t}\right)\right] .
$$

We then know from Theorem 3.4 of [12] that the Dirichlet form generated by $\left\{P_{t}^{\mu, h}\right\}_{t \geq 0}$ is identified.

Proposition 2.6. There exists a Dirichlet form $\left(\mathscr{E}^{\mu, h}, \mathscr{D}\left(\mathscr{E}^{\mu, h}\right)\right)$ on $L^{2}\left(h^{2} d x\right)$ corresponding to the semigroup $\left\{P_{t}^{\mu, h}\right\}$ and it has the representation as follows:

$$
\begin{align*}
\mathscr{E}^{\mu, h}(u, v) & =\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}(u(y)-u(x))(v(y)-v(x)) J(x, y) h(x) h(y) d x d y,  \tag{2.55}\\
\mathscr{D}\left(\mathscr{E}^{\mu, h}\right) & =\mathscr{F} . \tag{2.56}
\end{align*}
$$

Using this proposition, we can prove the 'if' part of Theorem 2.5. Note that $P_{t}^{\mu, h}$ admits the integral kernel $h(x)^{-1} p^{\mu}(t, x, y) h(y)^{-1}$ with respect to the measure $h^{2}(y) d y$. Since $1 \leq h(x) \leq C_{1}$ for some positive constant $C_{1}$, we see that the form defined by $(2.55)-(2.56)$ is the regular Dirichlet form on $L^{2}\left(\mathbb{R}^{d}\right)$ with jump measure $J_{1}(x, y):=$ $J(x, y) h(x) h(y)$. Moreover, there exists positive constants $C_{2}$ and $C_{3}$ such that

$$
\frac{C_{2}}{|x-y|^{d} \phi(|x-y|)} \leq J_{1}(x, y) \leq \frac{C_{3}}{|x-y|^{d} \phi(|x-y|)}
$$

where $\phi(r)=r^{\alpha}\left(\phi(r)=r^{\alpha} \exp \left(m_{0} r\right)\right)$ if $\left\{X_{t}\right\}_{t \geq 0}$ is $\alpha$-stable-like (relativistic $\alpha$-stablelike). Hence, we see that $h(x)^{-1} p^{\mu}(t, x, y) h(y)^{-1}$ has the same two sided estimates as those given in Theorems 2.2 and 2.3 for $p(t, x, y)$. Since $1 \leq h(x) \leq C_{1}$, so does $p^{\mu}(t, x, y)$.

The following example is constructed by Takeda and Uemura [36].
Example 2.1. Let $\sigma_{r}$ be a surface measure of sphere $\partial B_{r}=\{|x|=r\}$. Since the symmetric $\alpha$-stable process hits the sphere $\partial B_{r}$ if $1<\alpha \leq 2$ (see e.g., E44.19 in [23]), the measure $\sigma_{r}$ is in $K_{\infty}$. The measure $\sigma_{r}$ is then gaugeable if and only if

$$
\inf \left\{\mathscr{E}(u, u): \int_{\{|x|=r\}} u^{2} d \sigma_{r}=1\right\}>1
$$

Since the measure $\sigma_{r}$ is spherically symmetric, the infimum is attained by the function

$$
u(x)=c \mathbb{P}_{x}\left(\sigma_{\partial B_{r}}<\infty\right), x \in \mathbb{R}^{d},
$$

where $c=1 / \sqrt{\sigma_{r}\left(\partial B_{r}\right)}$. Let $\operatorname{Cap}^{(\alpha)}(\cdot)$ be the 0 -order capacity with respect to the symmetric $\alpha$-stable process. Then the infimum above equals to

$$
\frac{\mathrm{Cap}^{(\alpha)}\left(\partial B_{r}\right)}{\sigma_{r}\left(\partial B_{r}\right)},
$$

because

$$
\mathscr{E}^{(\alpha)}\left(\mathbb{P} \cdot\left(\sigma_{\partial B_{r}}<\infty\right), \mathbb{P} \cdot\left(\sigma_{\partial B_{r}}<\infty\right)\right)=\operatorname{Cap}^{(\alpha)}\left(\partial B_{r}\right)
$$

It is known (see e.g., Corollary 2.2 in [22]) that

$$
\operatorname{Cap}^{(\alpha)}\left(\partial B_{r}\right)=\frac{2 \pi^{(d+1) / 2} \Gamma\left(\frac{\alpha+d}{2}-1\right) \Gamma\left(\frac{\alpha}{2}\right)}{\Gamma\left(\frac{d}{2}\right) \Gamma\left(\frac{\alpha-1}{2}\right) \Gamma\left(\frac{d-\alpha}{2}\right)} r^{d-\alpha} .
$$

Therefore, the measure $\sigma_{r}$ is gaugeable if and only if

$$
\left\{\frac{\sqrt{\pi} \Gamma\left(\frac{d+\alpha}{2}-1\right) \Gamma\left(\frac{\alpha}{2}\right)}{\Gamma\left(\frac{\alpha-1}{2}\right) \Gamma\left(\frac{d-\alpha}{2}\right)}\right\}^{\frac{1}{\alpha-1}}>r
$$

We close this chapter introducing the extension of the stability by Kim and Kuwae.
Remark 2.2. The assumption that $\mu$ is of 0-order finite energy integral is unnecessary. Kim and Kuwae [16] extends the argument of h-transformation to more general class of function. Moreover, Kim and Kuwae [17] proved the stability of fundamental solutions for perturbations containing non-local parts.

## Chapter 3

## Critical Schrödinger forms and spectral functions

Let $\left\{X_{t}\right\}_{t \geq 0}$ be a transient jump Markov process on $\mathbb{R}^{d}$ and $\mu$ be a green-tight measure in the Kato class. In the previous chapter, we established the necessary and sufficient condition on the measure $\mu$ for the fundamental solution to satisfy the stability i.e. $p^{\mu}(t, x, y)$ has the same two sided estimates as $p(t, x, y)$ up to positive constants. This condition is said to be subcriticality of $\mu$, which describes the smallness of measure $\mu$. If $\mu$ is not subcritical, there are two cases to be considered: critical case and supercritical case. In these two cases, we expect that $p^{\mu}(t, x, y)$ has different behavior from $p(t, x, y)$. We have not determined the estimates of $p^{\mu}(t, x, y)$ however, there are some papers mentioning the estimates of the integral $\int_{\mathbb{R}^{d}} p^{\mu}(t, x, y) d y$. Note that this integral is equal to the expectation $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$. Furthermore, we see from the previous section that $\mu$ is subcritical if and only if $\sup _{x \in \mathbb{R}^{d}} \mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right)\right]<\infty$. Hence, the expectation $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ diverges as $t \rightarrow \infty$ if $\mu$ is critical or supercritical. To know the growth of $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$, we first consider the limit defined by

$$
\lim _{t \rightarrow \infty} \frac{1}{t} \log \mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]
$$

Takeda [30, 33] and Tsuchida [37] proved this limit is equal to the spectral bound as follows:

$$
C(\mu):=-\inf \left\{\mathscr{E}^{\mu}(u, u) u \in \mathscr{F}, \int_{\mathbb{R}^{d}} u^{2}(x) d x=1\right\} .
$$

[30] treats the standard Brownian motion, [33] treats the rotationally invariant $\alpha$-stable process and [37] treats the relativistic $\alpha$-stable process. It is known that $C(\mu)>0$ if $\mu$ is supercritical, and hence $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ grows exponentially. If $\mu$ is critical, $C(\mu)=0$ and the growth is slower than the exponential one. Simon [25] and Cranston, Kolokoltsov and et. al. [8] considered the same problem when $\left\{X_{t}\right\}_{t \geq 0}$ is the transient Brownian motion and $\mu=V \cdot m$ for non-negative $V \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$. The growth order of the expecta-
tion $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ depends on the dimension $d$ and satisfies

$$
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right] \sim \begin{cases}C_{1} h(x) \cdot t^{1 / 2} & (d=3) \\ C_{2} h(x) \cdot t / \log t & (d=4) \\ C_{3} h(x) \cdot t & (d \geq 5)\end{cases}
$$

where $h(x)$ is the harmonic function of $(-\Delta-V)$. In the sequel, we assume that $\mu$ is critical and $\left\{X_{t}\right\}_{t \geq 0}$ is the rotationally invariant $\alpha$-stable process. Takeda [34] proved that $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ is proportional to $t$ if $d / \alpha>2$. This is an analogy of the $d$-dimensional Brownian motion for $d \geq 5$. The outline of the proof is based on the probabilistic theory and as follows:

First we rewrite the expectation as follows:

$$
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=1+\mathbb{E}_{x}\left[\int_{0}^{t} \exp \left(A_{s}^{\mu}\right) d A_{s}^{\mu}\right] .
$$

Next we consider the transformation of Schrödinger semigroup by the harmonic function $h(x)$. In the critical case, we can also construct the harmonic function following the argument of Takeda and Tsuchida [35]. Using the $h$-transformed semigroup defined by

$$
\mathbb{E}_{x}^{h}\left[f\left(X_{t}\right)\right]=\mathbb{E}_{x}\left[\frac{h\left(X_{t}\right)}{h\left(X_{0}\right)} \exp \left(A_{t}^{\mu}\right) f\left(X_{t}\right)\right],
$$

we conclude that

$$
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=1+\mathbb{E}_{x}\left[\int_{0}^{t} \exp \left(A_{s}^{\mu}\right) d A_{s}^{\mu}\right]=1+h(x) \mathbb{E}_{x}^{h}\left[\int_{0}^{t} \frac{d A_{s}^{\mu}}{h\left(X_{s}\right)}\right] .
$$

Moreover, [35] proved that the harmonic function $h(x)$ satisfies

$$
C_{1}\left(1 \wedge|x|^{\alpha-d}\right) \leq h(x) \leq C_{2}\left(1 \wedge|x|^{\alpha-d}\right) .
$$

The function $h(x)$ is in $L^{2}\left(\mathbb{R}^{d}\right)$ for $d / \alpha>2$. Thus, the transformed process is an ergodic process with the finite invariant measure $h^{2} \cdot m$ and consequently, we obtain

$$
\frac{1}{t} \mathbb{E}_{x}^{h}\left[\int_{0}^{t} \frac{d A_{s}^{\mu}}{h\left(X_{s}\right)}\right] \rightarrow \int_{\mathbb{R}^{d}} h(x) \mu(d x) \quad(t \rightarrow \infty)
$$

Hence we determine the growth of $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$.
However, we cannot apply this method for the case $d / \alpha \leq 2$ because $h(x)$ is not in $L^{2}\left(\mathbb{R}^{d}\right)$. Thus, we use other analytical method based on the argument of Simon [25]. In this chapter, we consider the case $\left\{X_{t}\right\}_{t \geq 0}$ is the rotationally invariant 1 -stable process on $\mathbb{R}^{2}$ and $\mu$ is absolutely continuous with respect to the Lebesgue measure. This chapter is organized as follows: In Section 1, we calculate the transition density function of $\left\{X_{t}\right\}_{t \geq 0}$ and give the asymptotic expansion of the resolvent $\left\{G_{\beta}\right\}$ when $\beta$ tends to 0 . In Section 2 we consider the asymptotic behavior of $\left\{G_{\beta}^{\mu}\right\}$, namely the resolvent of Schrödinger form. Furthermore we apply the Tauberian theorem, which describes the
relation between the asymptotic behavior of resolvent and that of semigroup. In Section 3 we consider the behavior of the spectral function. For fixed measure $\mu$, the spectral function $C(\lambda)$ is defined as $C(\lambda)=C(\lambda \mu)$, where $C(\lambda \mu)$ stands for the spectral bound of $\lambda \mu$. Takeda and Tsuchida [35] established the criterion of differentiability of $C(\lambda)$. Using the asymptotic expansion obtained in Section 2, we can determine the precise behavior of $C(\lambda)$.

### 3.1 The asymptotic behavior of the resolvent

We first calculate concrete transition density function and resolvent of the Markov process. Let $\left\{X_{t}\right\}_{t \geq 0}$ be the rotationally invariant $\alpha$-stable process. The characteristic function of $\left\{X_{t}\right\}_{t \geq 0}$ satisfies

$$
\begin{equation*}
\mathbb{E}_{0}\left[\exp \left(i u \cdot X_{t}\right)\right]=\exp \left(-t|u|^{\alpha}\right) . \tag{3.1}
\end{equation*}
$$

Using the transition density function, (3.1) is rewritten as follows:

$$
\int_{\mathbb{R}^{d}} \exp (i u \cdot(y-x)) p(t, x, y) d y=\exp \left(-t|u|^{\alpha}\right)
$$

Furthermore, we obtain the explicit transition density function by the inverse Fourier transform:

$$
\begin{equation*}
p(t, x, y)=\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \exp \left(-t|u|^{\alpha}-i u \cdot(y-x)\right) d u . \tag{3.2}
\end{equation*}
$$

Proposition 3.1. Let $p(t, x, y)$ be the heat kernel of 1 -stable process on $\mathbb{R}^{2}$. Then it follows that

$$
\begin{equation*}
p(t, x, y)=\frac{t}{2 \pi\left(t^{2}+|x-y|^{2}\right)^{\frac{3}{2}}} . \tag{3.3}
\end{equation*}
$$

Proof. We have only to substitute $d=2$ and $\alpha=1$ in (3.2). Using the polar coordinates, we obtain

$$
\begin{aligned}
p(t, x, y) & =\frac{1}{(2 \pi)^{2}} \int_{\mathbb{R}^{2}} \exp (-t|u|+i(x-y) \cdot u) d u \\
& =\frac{1}{(2 \pi)^{2}} \int_{0}^{2 \pi} \int_{0}^{\infty} r \exp (-t r+i r|x-y| \cos \theta) d r d \theta \\
& =\frac{1}{(2 \pi)^{2}} \int_{0}^{2 \pi} \frac{d \theta}{(t-i|x-y| \cos \theta)^{2}}=\frac{1}{2 \pi^{2}} \int_{0}^{\pi} \frac{d \theta}{(t-i|x-y| \cos \theta)^{2}} .
\end{aligned}
$$

Moreover we can calculate the integral of the right hand side as follows:

$$
\int_{0}^{\pi} \frac{d \theta}{(t-i|x-y| \cos \theta)^{2}}=\int_{0}^{\frac{\pi}{2}}\left(\frac{1}{(t+i|x-y| \cos \theta)^{2}}+\frac{1}{(t-i|x-y| \cos \theta)^{2}}\right) d \theta
$$

$$
\begin{align*}
& =2 \int_{0}^{\frac{\pi}{2}} \frac{t^{2}-|x-y|^{2} \cos ^{2} \theta}{\left(t^{2}+|x-y|^{2} \cos ^{2} \theta\right)^{2}} d \theta \\
& =2 \int_{0}^{\frac{\pi}{2}} \frac{(C-B \cos 2 \theta)}{(A+B \cos 2 \theta)^{2}} d \theta=\int_{0}^{\pi} \frac{(C-B \cos \theta)}{(A+B \cos \theta)^{2}} d \theta \tag{3.4}
\end{align*}
$$

where $A=t^{2}+|x-y|^{2} / 2, B=|x-y|^{2} / 2$ and $C=t^{2}-|x-y|^{2} / 2$. Here we consider the integrals as follows:

$$
\int_{0}^{\pi} \frac{1}{(A+B \cos \theta)^{2}} d \theta, \quad \int_{0}^{\frac{\pi}{2}} \frac{\cos \theta}{(A+B \cos \theta)^{2}} d \theta
$$

The former is calculated as follows:

$$
\begin{aligned}
\int_{0}^{\pi} \frac{1}{(A+B \cos \theta)^{2}} d \theta & =\int_{0}^{\infty}\left(A+B \cdot \frac{1-t^{2}}{1+t^{2}}\right)^{-2} \cdot \frac{2}{1+t^{2}} d t \\
& =\int_{0}^{\infty} \frac{2\left(1+t^{2}\right)}{\left((A+B)+(A-B) t^{2}\right)^{2}} d t \\
& =\int_{0}^{\infty}\left(\frac{2 /(A-B)}{(A+B)+(A-B) t^{2}}-\frac{4 B /(A-B)}{\left((A+B)+(A-B) t^{2}\right)^{2}}\right) d t
\end{aligned}
$$

Here we used the substitution of $\tan (\theta / 2)=t$ in the first equality. Since both $A+B$ and $A-B$ are positive, the right hand side of the above formula is equal to

$$
\begin{equation*}
\frac{2}{(A-B)^{2}} \cdot\left(\frac{A-B}{A+B}\right)^{\frac{1}{2}} \cdot \frac{\pi}{2}-\frac{4 B}{(A-B)^{3}} \cdot\left(\frac{A-B}{A+B}\right)^{\frac{3}{2}} \cdot \frac{\pi}{4}=\pi \cdot \frac{A}{(A-B)^{\frac{3}{2}}(A+B)^{\frac{3}{2}}} . \tag{3.5}
\end{equation*}
$$

The latter is calculated as follows:

$$
\begin{align*}
\int_{0}^{\pi} \frac{\cos \theta}{(A+B \cos \theta)^{2}} d \theta & =\int_{0}^{\infty}\left(A+B \cdot \frac{1-t^{2}}{1+t^{2}}\right)^{-2} \cdot \frac{1-t^{2}}{1+t^{2}} \cdot \frac{2}{1+t^{2}} d t \\
& =\int_{0}^{\infty} \frac{2\left(1-t^{2}\right)}{\left((A+B)+(A-B) t^{2}\right)^{2}} d t \\
& =\int_{0}^{\infty}\left(\frac{-2 /(A-B)}{(A+B)+(A-B) t^{2}}+\frac{4 A /(A-B)}{\left((A+B)+(A-B) t^{2}\right)^{2}}\right) d t \\
& =-\frac{2}{(A-B)^{2}} \cdot\left(\frac{A-B}{A+B}\right)^{\frac{1}{2}} \cdot \frac{\pi}{2}+\frac{4 A}{(A-B)^{3}} \cdot\left(\frac{A-B}{A+B}\right)^{\frac{3}{2}} \cdot \frac{\pi}{4} \\
& =\pi \cdot \frac{-B}{(A-B)^{\frac{3}{2}}(A+B)^{\frac{3}{2}}} \tag{3.6}
\end{align*}
$$

(3.5) and (3.6) imply that the right hand side of (3.4) is equal to

$$
\pi \cdot \frac{A C+B^{2}}{(A-B)^{\frac{3}{2}}(A+B)^{\frac{3}{2}}}=\pi \cdot \frac{t^{4}}{t^{3}\left(t^{2}+|x-y|^{2}\right)^{\frac{3}{2}}}=\pi \cdot \frac{t}{\left(t^{2}+|x-y|^{2}\right)^{\frac{3}{2}}} .
$$

Thus, we obtain (3.3).

We next consider the behavior of the $\beta$-order resolvent when $\beta$ tends to 0 .
Lemma 3.1. The $\beta$-order resolvent kernel $G_{\beta}(x, y)$ has the asymptotic expansion as follows.

$$
\begin{equation*}
G_{\beta}(x, y)=G_{0}(x, y)+\frac{1}{2 \pi} \beta \log \beta+\frac{\gamma-\log 2+\log |x-y|}{2 \pi} \beta+\mathscr{O}\left(\beta^{2}\right) . \tag{3.7}
\end{equation*}
$$

Proof. (3.3) implies that

$$
G_{\beta}(x, y)=\int_{0}^{\infty} p(t, x, y) e^{-\beta t} d t=\frac{1}{2 \pi} \int_{0}^{\infty} \frac{t e^{-\beta t}}{\left(t^{2}+|x-y|^{2}\right)^{\frac{3}{2}}} d t
$$

Substituting $t=|x-y| \sinh z$, we have

$$
\begin{align*}
G_{0}(x, y) & =\frac{1}{2 \pi} \int_{0}^{\infty} \frac{t}{\left(t^{2}+|x-y|^{2}\right)^{\frac{3}{2}}} d t=\frac{1}{2 \pi} \int_{0}^{\infty} \frac{|x-y| \sinh z}{|x-y|^{3} \cosh ^{3} z} \cdot|x-y| \cosh z d z \\
& =\frac{1}{2 \pi|x-y|} \int_{0}^{\infty} \frac{\sinh z}{\cosh ^{2} z} d z=\frac{1}{2 \pi|x-y|} \tag{3.8}
\end{align*}
$$

and

$$
\begin{align*}
& G_{\beta}(x, y)=\frac{1}{2 \pi} \int_{0}^{\infty} \frac{t e^{-\beta t}}{\left(t^{2}+|x-y|^{2}\right)^{\frac{3}{2}}} d t=\frac{1}{2 \pi|x-y|} \int_{0}^{\infty} \frac{\sinh z}{\cosh ^{2} z} e^{-\beta|x-y| \sinh z} d z \\
& =\frac{1}{2 \pi|x-y|}\left\{\left[\frac{-1}{\cosh z} e^{-\beta|x-y| \sinh z}\right]_{z=0}^{\infty}-\int_{0}^{\infty} \beta|x-y| e^{-\beta|x-y| \sinh z} d z\right\} \\
& =\frac{1}{2 \pi|x-y|}-\frac{\beta}{2 \pi} \int_{0}^{\infty} e^{-\beta|x-y| \sinh z} d z=G_{0}(x, y)-\frac{\beta}{2 \pi} \int_{0}^{\infty} e^{-\beta|x-y| \sinh z} d z . \tag{3.9}
\end{align*}
$$

Now we evaluate the second term of (3.9). We see that

$$
\begin{aligned}
\int_{0}^{\infty} e^{-\beta|x-y| \sinh z} d z & =\int_{0}^{\infty} \exp \left(-\frac{\beta|x-y|}{2} e^{z}\right) \cdot \exp \left(\frac{\beta|x-y| e^{-z}}{2}\right) d z \\
& =\int_{0}^{\infty} \exp \left(-\frac{\beta|x-y|}{2} e^{z}\right) \cdot \sum_{n=0}^{\infty} \frac{1}{n!}\left(\frac{\beta|x-y| e^{-z}}{2}\right)^{n} d z
\end{aligned}
$$

Using the Fubini's theorem, the right hand side of the above formula is equal to

$$
\sum_{n=0}^{\infty} \int_{0}^{\infty} \frac{(\beta|x-y|)^{n}}{2^{n} \cdot n!} \exp \left(-n z-\frac{\beta|x-y|}{2} e^{z}\right) d z=: \sum_{n=0}^{\infty} I_{n},
$$

First we consider the integral $I_{0}$. Substituting $w=\beta|x-y| e^{z} / 2$, we see that

$$
\begin{equation*}
I_{0}=\int_{\frac{\beta|x-y|}{2}}^{\infty} \frac{e^{-w}}{w} d w=-\log \left(\frac{\beta|x-y|}{2}\right)-\gamma-\sum_{n=1}^{\infty} \frac{(-\beta|x-y|)^{n}}{2^{n} \cdot n \cdot n!}, \tag{3.10}
\end{equation*}
$$

where $\gamma$ is Euler's Gamma and we use the residue integral as follows:

$$
\int_{x}^{\infty} \frac{e^{-w}}{w} d w=-\log x-\gamma-\sum_{n=1}^{\infty} \frac{(-x)^{n}}{n \cdot n!}
$$

Furthermore, for $n \geq 1, I_{n}$ has the upper bound as follows:

$$
\begin{equation*}
I_{n} \leq \frac{(\beta|x-y|)^{n}}{2^{n} \cdot n!} \int_{0}^{\infty} e^{-n z} d z=\frac{(\beta|x-y|)^{n}}{2^{n} \cdot n \cdot n!} \tag{3.11}
\end{equation*}
$$

Thus, for fixed $x, y \in \mathbb{R}^{d}$, we obtain $\sum_{n=1}^{\infty} I_{n}=\mathscr{O}(\beta)$ and

$$
\begin{equation*}
\int_{0}^{\infty} e^{-\beta|x-y| \sinh z} d z=-\log \left(\frac{\beta|x-y|}{2}\right)-\gamma+\mathscr{O}(\beta) \quad(\beta \rightarrow 0) \tag{3.12}
\end{equation*}
$$

Hence we obtain (3.7).

### 3.2 Growth of Feynman-Kac semigroups

Let $(\mathscr{E}, \mathscr{F})$ be the symmetric jump Dirichlet form associated with $\left\{X_{t}\right\}_{t \geq 0}$, i.e.

$$
\mathscr{E}(u, u)=\int_{\mathbb{R}^{2} \times \mathbb{R}^{2}}(u(y)-u(x))^{2} \frac{\kappa_{0}}{|x-y|^{3}} d x d y, \quad \mathscr{F}=H^{\frac{1}{2}}\left(\mathbb{R}^{2}\right)
$$

where $\kappa_{0}$ is an appropriate positive constant and $H^{\frac{1}{2}}\left(\mathbb{R}^{2}\right)$ is the Sobolev space with order $1 / 2$. Let $H$ is a non-local operator satisfying

$$
\mathscr{E}(u, v)=(H u, v),
$$

where the right hand side is the inner product of $L^{2}\left(\mathbb{R}^{2}\right)$.
Let $\mu$ be a Green-tight measure in Definition 2.3 and define the Schrödinger form by

$$
\mathscr{E}^{\mu}(u, u)=\mathscr{E}(u, u)-\int_{\mathbb{R}^{d}} u^{2} d \mu .
$$

We divide the class of Green-tight measures into three subclasses according to the smallness of measures.

Definition 3.1. For $\mu \in \boldsymbol{K}_{\infty}$, define

$$
\lambda(\mu):=\inf \left\{\mathscr{E}(u, u) ; u \in \mathscr{F}_{e}, \int_{\mathbb{R}^{d}} u^{2} d \mu=1\right\}
$$

(1) $\mu$ is said to be subcritical if $\lambda(\mu)>1$.
(2) $\mu$ is said to be critical if $\lambda(\mu)=1$.
(3) $\mu$ is said to be supercritical if $\lambda(\mu)<1$.

We proved that the stability of fundamental solution is equivalent to the subcriticality of $\mu$ in Theorem 2.5. Thus we see that the behavior of $p^{\mu}(t, x, y)$ is different from that of $p(t, x, y)$ when $\mu$ is critical or supercritical. We have not determined the behavior of $p^{\mu}(t, x, y)$. Instead, we consider the behavior of

$$
\begin{equation*}
\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=\int_{\mathbb{R}^{d}} p^{\mu}(t, x, y) d y \tag{3.13}
\end{equation*}
$$

Note that the subcriticality of $\mu$ is equivalent to the gaugeability of $A_{t}^{\mu}$, namely,

$$
\sup _{x \in \mathbb{R}^{d}} \mathbb{E}_{x}\left[\exp \left(A_{\infty}^{\mu}\right)\right]<\infty
$$

Thus, the value $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ diverges as $t$ tends to $\infty$ in the other cases. First we define the spectral bound $C(\mu)$ by

$$
\begin{equation*}
C(\mu)=-\inf \left\{\mathscr{E}^{\mu}(u, u) ; u \in \mathscr{F}, \int_{\mathbb{R}^{d}} u^{2}(x) d x=1\right\} . \tag{3.14}
\end{equation*}
$$

The following theorem is given in Theorem 5.2 of [33].
Theorem 3.1. For Green-tight measure $\mu$, it follows that

$$
\lim _{t \rightarrow \infty} \frac{1}{t} \log \mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=C(\mu)
$$

Moreover we can characterize supercriticality of $\mu$ by using the spectral bound $C(\mu)$. The following proposition is proved in Lemma 2.2 of [35].

Proposition 3.2. The measure $\mu$ is supercritical if and only if $C(\mu)>0$.
Combining Theorem 3.1 and Proposition 3.2, we can obtain the following corollary.
Corollary 3.1. Suppose the Green-tight measure $\mu$ is supercritical. $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ has exponential growth as $t \rightarrow \infty$.

If $\mu$ is critical, $C(\mu)=0$ and thus $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ does not have the exponential growth. We first introduce two items: one is the ground state, equivalently the harmonic function with respect to Schrödinger operator $H-\mu$ and the other is the result for Brownian motion.

The following theorem is proved in Theorem 3.4 of Takeda and Tsuchida [35].
Theorem 3.2. For $\mu \in \boldsymbol{K}_{\infty}$, the extended Dirichlet space $\mathscr{F}_{e}$ is compactly embedded into $L^{2}(\mu)$

Thus we see that there exists a $h_{0} \in \mathscr{F}_{e}$ such that

$$
\mathscr{E}\left(h_{0}, h_{0}\right)=\lambda(\mu), \quad \int_{\mathbb{R}^{d}} h_{0}^{2} d \mu=1
$$

We call $h_{0}$ ground state of the Schrödinger form $\mathscr{E}^{\mu}$. Moreover, they showed the precise behavior of $h_{0}$ as follows:

Lemma 3.2. The function $h_{0}(x)$ is continuous and there exist positive constants $C_{1}$ and $C_{2}$ such that

$$
C_{1}\left(1 \wedge \frac{1}{|x|}\right) \leq h_{0}(x) \leq C_{2}\left(1 \wedge \frac{1}{|x|}\right)
$$

Proof. We know from Lemma 4.7 and Proposition 4.14 of [35] that the ground state $h_{0}(x)$ satisfies

$$
C_{1}\left(1 \wedge \frac{1}{|x|^{d-\alpha}}\right) \leq h_{0}(x) \leq C_{2}\left(1 \wedge \frac{1}{|x|^{d-\alpha}}\right)
$$

for the rotationally invariant $\alpha$-stable process on $\mathbb{R}^{d}$. Since $d=2$ and $\alpha=1$, the assertion follows.

Our goal is to prove the following theorem:
Theorem 3.3. Suppose $\left\{X_{t}\right\}_{t \geq 0}$ is the 1 -stable process on $\mathbb{R}^{2}$ and $\mu$ is critical. Assume $\mu=V \cdot m$ for positive $V \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$. Then there exists a positive constant $C_{1}$ such that

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \frac{\log t}{t} \mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]=C_{1} h_{0}(x), \tag{3.15}
\end{equation*}
$$

where $h_{0}(x)$ is the ground state of the Schrödinger operator $H-V$.
In the sequel, we consider the growth of $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ as $t \rightarrow \infty$ when $V$ is critical. Denote the semigroup and the resolvent associated with $\mathscr{E}^{\mu}$ by $\left\{P_{t}^{\mu}\right\}$ and $\left\{G_{\beta}^{\mu}\right\}$ respectively. Recall that

$$
P_{t}^{\mu} f(x)=\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right) f\left(X_{t}\right)\right]
$$

and consequently $P_{t}^{\mu} 1(x)=\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$.
Proposition 3.3. The semigroup $\left\{P_{t}^{\mu}\right\}$ satisfies the following formula:

$$
\begin{equation*}
P_{t}^{\mu} 1=1+\int_{0}^{t} P_{s}^{\mu} V d s \tag{3.16}
\end{equation*}
$$

Proof. Noting that the absolute continuity of the measure $\mu$, we have

$$
\exp \left(A_{t}^{\mu}\right)-1=\int_{0}^{t} \exp \left(A_{s}^{\mu}\right) d A_{s}^{\mu}=\int_{0}^{t} \exp \left(A_{s}^{\mu}\right) V\left(X_{s}\right) d s
$$

Thus, it follows that

$$
\begin{aligned}
P_{t}^{\mu} f(x) & =\mathbb{E}_{x}\left[\left(1+\int_{0}^{t} \exp \left(A_{s}^{\mu}\right) V\left(X_{s}\right) d s\right) f\left(X_{t}\right)\right] \\
& =P_{t} f(x)+\mathbb{E}_{x}\left[\int_{0}^{t} \exp \left(A_{s}^{\mu}\right) V\left(X_{s}\right) f\left(X_{t}\right) d s\right]
\end{aligned}
$$

$$
=P_{t} f(x)+\int_{0}^{t} \mathbb{E}_{x}\left[\exp \left(A_{s}^{\mu}\right) V\left(X_{s}\right) f\left(X_{t}\right)\right] d s
$$

Using the Markov property, we can rewrite the above formula as follows:

$$
\begin{aligned}
& P_{t} f(x)+\int_{0}^{t} \mathbb{E}_{x}\left[\mathbb{E}_{x}\left[\exp \left(A_{s}^{\mu}\right) V\left(X_{s}\right) f\left(X_{t}\right) \mid \mathscr{F}_{s}\right]\right] \\
= & P_{t} f(x)+\int_{0}^{t} \mathbb{E}_{x}\left[\exp \left(A_{s}^{\mu}\right) V\left(X_{s}\right) \mathbb{E}_{x}\left[f\left(X_{t}\right) \mid \mathscr{F}_{s}\right]\right] \\
= & P_{t} f(x)+\int_{0}^{t} P_{s}^{\mu}\left(V P_{t-s} f\right)(x) d s .
\end{aligned}
$$

Hence, we obtain

$$
P_{t}^{\mu} f(x)=P_{t} f(x)+\int_{0}^{t} P_{s}^{\mu}\left(V P_{t-s} f\right)(x) d s
$$

If we substitute $f$ by the constant 1 , the conservativeness of $\left\{X_{t}\right\}_{t \geq 0}$ implies (3.16).
In order to evaluate $P_{s}^{\mu} V$, we consider its Laplace transform, $G_{\beta}^{\mu} V$. The following proposition is proved by the resolvent equation.

Proposition 3.4. It follows that

$$
\begin{equation*}
G_{\beta}^{\mu} V(x)=\left(1-L_{\beta}\right)^{-1}\left(G_{\beta} V\right)(x) \tag{3.17}
\end{equation*}
$$

where $G_{\beta} V$ is defined by

$$
G_{\beta} V(x)=\int_{\mathbb{R}^{2}} G_{\beta}(x, y) V(y) d y
$$

and $L_{\beta}$ is an operator defined by $L_{\beta} f(x)=G_{\beta}(V f)(x)$.
Proof. Since $G_{\beta}=(H+\beta)^{-1}$ and $G_{\beta}^{\mu}=(H-V+\beta)^{-1}$, we obtain

$$
(H-V+\beta)^{-1}-(H+\beta)^{-1}=(H+\beta)^{-1} V(H-V+\beta)^{-1}
$$

by the resolvent equation. If we solve this equation with respect to $(H-V+\beta)^{-1}$, we have

$$
\left(1-(H+\beta)^{-1} V\right)(H-V+\beta)^{-1}=(H+\beta)^{-1}
$$

and consequently

$$
(H-V+\beta)^{-1}=\left(1-(H+\beta)^{-1} V\right)^{-1}(H+\beta)^{-1}
$$

If we operate the function $V$ on both sides, we can have the desired result.
For the calculation of the resolvent $G_{\beta}^{\mu}$, we have to consider the behavior of the function $G_{\beta} V$ and operator $L_{\beta}$. We begin with function $G_{\beta} V$.

Lemma 3.3. For $\beta \geq 0$ and $V \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$, $G_{\beta} V$ belongs to the class $L^{\infty}\left(\mathbb{R}^{2}\right)$ and converges uniformly to $G_{0} V$ as $\beta \rightarrow 0$.

Proof. Since the integral kernels $\left\{G_{\beta}(x, y)\right\}_{\beta \geq 0}$ is increasing as $\beta \rightarrow 0$, it suffices to prove $G_{0} V \in L^{\infty}\left(\mathbb{R}^{2}\right)$. We know from (3.8) that

$$
G_{0}(x, y)=\frac{C_{1}}{|x-y|}
$$

Set $R>0$ satisfying $\operatorname{supp}(V) \subset B(0, R)$, where $B(0, R)$ is a ball in $\mathbb{R}^{2}$ with center $O$ and radius $R$. If $|x|>2 R$, it follows that

$$
\begin{equation*}
G_{0} V(x) \leq \frac{C_{1}}{|x|-R} \int_{\mathbb{R}^{2}} V(y) d y \leq \frac{C_{2}}{|x|-R} \tag{3.18}
\end{equation*}
$$

As for $|x| \leq 2 R$, we obtain

$$
\begin{equation*}
G_{0} V(x) \leq \int_{|x-y| \leq 3 R} \frac{V(y)}{|x-y|} \leq\|V\|_{\infty} \int_{|x-y| \leq 3 R} \frac{d y}{|x-y|} \leq C_{3} \tag{3.19}
\end{equation*}
$$

Thus we can conclude $G_{0} V \in L^{\infty}\left(\mathbb{R}^{d}\right)$ and so is $G_{\beta} V$.
Next we prove $\lim _{\beta \rightarrow 0}\left\|G_{0} V-G_{\beta} V\right\|_{\infty}=0$. We may assume that $|x|$ is bounded. Indeed, for any $\varepsilon>0$, (3.18) implies that there exists $R_{\varepsilon}>0$ such that $G_{0} V(x)<\varepsilon$ for $|x|>R_{\varepsilon}$. We may also assume that the support of $V$ is contained in $B_{0}\left(R_{\varepsilon}\right)$. We reconsider the error term in asymptotic expansion of the resolvent kernel as follows:

$$
\begin{equation*}
G_{0}(x, y)-G_{\beta}(x, y)=\frac{\beta}{2 \pi} \int_{0}^{\infty} e^{-\beta|x-y| \sinh z} d z \tag{3.20}
\end{equation*}
$$

Estimations in (3.10) and (3.11) imply

$$
\begin{equation*}
\beta \int_{0}^{\infty} e^{-\beta|x-y| \sinh z} d z \leq-\beta \log \left(\frac{\beta|x-y|}{2}\right)-\beta \gamma+2 \beta \exp \left(\frac{\beta|x-y|}{2}\right) \tag{3.21}
\end{equation*}
$$

Noting that the third term in (3.12) is uniformly for $x, y \in \mathbb{R}^{2}$ satisfying $|x-y| \leq 2 R_{\varepsilon}$, we obtain

$$
\begin{aligned}
\left|G_{0} V(x)-G_{\beta} V(x)\right| & \leq \int_{B_{0}\left(R_{\varepsilon}\right)}\left(-\beta \log \left(\frac{\beta|x-y|}{2}\right)+2 \beta \exp \left(\frac{\beta|x-y|}{2}\right)\right) V(y) d y \\
& \leq-\beta \log \beta \int_{B_{0}\left(R_{\varepsilon}\right)} V(y) d y \\
& +\beta \int_{B_{0}\left(R_{\varepsilon}\right)}\left(\left|\log \left(\frac{|x-y|}{2}\right)\right|+2 \exp \left(\frac{\beta|x-y|}{2}\right)\right) V(y) d y
\end{aligned}
$$

If we choose $\beta$ sufficiently small according to $R_{\varepsilon}$, the right hand side of the previous formula is bounded by $\varepsilon$. Therefore we have proved that $\left\|G_{0} V-G_{\beta} V\right\|_{\infty} \rightarrow 0$ as $\beta \rightarrow$ 0 .

We next consider the behavior of operators $\left\{L_{\beta}\right\}_{\beta \geq 0}$.

Lemma 3.4. The operator $L_{\beta}$ defined by $f \rightarrow G_{\beta}(V f)$ is a compact operator from $L^{\infty}\left(\mathbb{R}^{2}\right)$ to $L^{\infty}\left(\mathbb{R}^{2}\right)$. Furthermore it is continuous with respect to norm when $\beta \rightarrow 0$.

Proof. Since the multiple operator $V$ is compact and the resolvent operator $G_{\beta}(\beta>0)$ is bounded, $L_{\beta}$ is a compact operator. Furthermore, $L_{\beta}$ converges to $L_{0}$ in the space of bounded operators on $L^{\infty}\left(\mathbb{R}^{2}\right)$. Indeed, for any $f \in L^{2}\left(\mathbb{R}^{2}\right)$ with $\|f\|_{\infty}=1, V(x) f(x)$ has a compact support. Thus, using the same argument of Lemma 3.3 we obtain

$$
\left\|L_{\beta} f-L_{0} f\right\|_{\infty}=\left\|G_{\beta}(V f)-G_{0}(V f)\right\|_{\infty} \rightarrow 0 \quad(\beta \rightarrow 0)
$$

Moreover, the space of compact operators is a closed subspace in the space of bounded operators. Hence, $L_{0}$ is also a compact operator and we obtain the desired result.

Next we consider the operator defined on $L^{2}\left(\mathbb{R}^{2}\right)$.
Lemma 3.5. Define the operator $K_{\beta}: L^{2}\left(\mathbb{R}^{2}\right) \rightarrow L^{2}\left(\mathbb{R}^{2}\right)$ by

$$
K_{\beta} f(x)=V^{\frac{1}{2}} G_{\beta}\left(V^{\frac{1}{2}} f\right)(x)
$$

Then $e_{\beta}$ is an eigenvalue of $L_{\beta}$ if and only if it is an eigenvalue of $K_{\beta}$.
Proof. Suppose the function $g$ satisfies $L_{\beta} g=e_{\beta} g$ for some $e_{\beta}>0$. For $h=V^{\frac{1}{2}} g$, it follows that $K_{\beta} h=e_{\beta} h$.

The following lemma is based on Theorem 6.1 of Klaus and Simon [19].
Lemma 3.6. Suppose $\mu$ is critical and $\mu=V \cdot m$ for positive $V \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$. Denote by $e_{\beta}$ the principal eigenvalue of $K_{\beta}$. Then $e_{\beta}$ has the expansion as follows:

$$
\begin{equation*}
e_{\beta}=1+C_{1} \beta \log \beta+\ldots \quad\left(C_{1}>0\right) \tag{3.22}
\end{equation*}
$$

Proof. Note that $\lambda \mu$ is supercritical for $\lambda>1$. For $\lambda>1$ and $\beta>0$, consider the equation

$$
\begin{equation*}
(H-\lambda V) u_{\lambda}=-\beta u_{\lambda} . \tag{3.23}
\end{equation*}
$$

(3.23) is rewritten as follows:

$$
(H+\beta) u_{\lambda}=\lambda V u_{\lambda}
$$

Noting that $(H+\beta)^{-1}$ is the resolvent operator with order $\beta$, we have

$$
\lambda^{-1} u_{\lambda}=G_{\beta}\left(V u_{\lambda}\right)
$$

Moreover, if we substitute $w_{\lambda}=V^{\frac{1}{2}} u_{\lambda}$, we have

$$
\lambda^{-1} w_{\lambda}=V^{1 / 2} G_{\beta}\left(V^{1 / 2} w_{\lambda}\right)
$$

Thus (3.23) is equivalent to the equation $K_{\beta} w_{\lambda}=\lambda^{-1} w_{\lambda}$. Since $G_{\beta}(x, y)$ satisfies asymptotic expansion (3.7), the operator $K_{\beta}$ has asymptotic expansion as follows:

$$
\begin{equation*}
K_{\beta}=K_{0}+\beta \log \beta D_{1}+\beta D_{2}+\ldots \tag{3.24}
\end{equation*}
$$

where $D_{1}$ and $D_{2}$ are operators defined by

$$
\begin{aligned}
& D_{1} f(x)=\frac{V^{\frac{1}{2}}(x)}{2 \pi} \cdot \int_{\mathbb{R}^{d}} V^{\frac{1}{2}}(y) f(y) d y \\
& D_{2} f(x)=\frac{\gamma-\log 2}{2 \pi} V^{\frac{1}{2}}(x) \cdot \int_{\mathbb{R}^{d}} V^{\frac{1}{2}}(y) f(y) d y+\frac{V^{\frac{1}{2}}(x)}{2 \pi} \cdot \int_{\mathbb{R}^{d}} \log |x-y| V^{\frac{1}{2}}(y) f(y) d y
\end{aligned}
$$

The expansion (3.24) is the same form as that of 4-dimensional Brownian motion treated in Klaus and Simon [19]. Moreover $V^{1 / 2} h_{0}$ satisfies $K_{0}\left(V^{1 / 2} h_{0}\right)=V^{1 / 2} h_{0}$ where $h_{0}$ is the ground state of $\mathscr{E}^{\mu}$ defined in Lemma 3.2. Since $h_{0}$ is positive, $\left(V^{1 / 2}, V^{1 / 2} h_{0}\right) \neq 0$ and the proof of Theorem 6.1 in [19] implies

$$
e_{\beta}=1+C_{1} \beta \log \beta+\ldots
$$

for $C_{1} \neq 0$. Moreover, noting that $e_{\beta}=\lambda^{-1}<1$, we see $C_{1}>0$.
Now we obtain the behavior of the principal eigenvalue of the compact operator $L_{\beta}$. We consider the decomposition of the operator $\left(1-L_{\beta}\right)^{-1}$ as follows:

Lemma 3.7. Denote by $P_{\beta}$ the projection operator associated with $L_{\beta}$. Then it holds that

$$
\left(1-L_{\beta}\right)^{-1}=\left(1-e_{\beta}\right)^{-1} P_{\beta}+Q_{\beta}\left(1-P_{\beta}\right)
$$

where $Q_{\beta}$ is norm continuous with finite limit as $\beta \rightarrow 0$.
Proof. Since $L_{\beta}$ is a compact operator for $\beta>0, L_{\beta}$ has a spectral decomposition. Denote by $e_{\beta}^{\prime}$ the second largest eigenvalue of $L_{\beta}$. Then $e_{0}^{\prime}<e_{0}=1$ and $\left\|Q_{\beta}\right\|_{\infty}$ is smaller than $2\left(1-e_{0}^{\prime}\right)^{-1}$ if $\beta$ is sufficiently small.

The previous lemma and (3.17) imply that

$$
\lim _{\beta \rightarrow 0}\left\|G_{\beta}^{\mu} V-\left(1-e_{\beta}\right)^{-1} P_{\beta} G_{\beta} V\right\|_{\infty}<\infty
$$

Moreover $P_{0} \psi=h_{0} l(\psi)$, where $l$ is a suitable linear functional and $h_{0}$ is the harmonic function satisfying $(H-V) h_{0}=0$. Thus, combining (3.22), we have

$$
\begin{equation*}
\lim _{\beta \rightarrow 0}-\beta \log \beta G_{\beta}^{\mu}(V)=C_{1} h_{0}(x) \tag{3.25}
\end{equation*}
$$

where $C_{1}$ is some positive constant. Note that $C_{1} \neq 0$ follows because $\left(V, h_{0}\right) \neq 0$. Now we obtain the behavior of $G_{\beta}^{\mu} V$ as $\beta \rightarrow 0$. Applying the Tauberian theorem in [24], we can prove Theorem 3.3:

Proof of Theorem 3.3 First we fix $x \in \mathbb{R}^{2}$ and define the class of measures $\left\{v_{\beta}^{x}\right\}_{\beta>0}$ on $\mathbb{R}_{+}:=[0, \infty)$ by

$$
v_{\beta}^{x}(A)=-\beta \log \beta \int_{\beta^{-1} A} P_{s}^{\mu} V(x) d s=-\beta \log \beta \int_{A} P_{\beta^{-1} s}^{\mu} V(x) \cdot \beta^{-1} d s
$$

Denote the Lebesgue measure on $\mathbb{R}_{+}$by $v$. We have only to prove that $v_{\beta}^{x}([0,1])$ converges to $C_{1} h_{0}(x) v([0,1])$ as $\beta \rightarrow 0$ up to positive constant multiple. Indeed, it follows that

$$
v_{\beta}^{x}[0,1]=-\beta \log \beta \int_{0}^{\beta^{-1}} P_{s}^{\mu} V(x) d s
$$

and substituting $\beta^{-1}=t$, we obtain the desired result. First, (3.25) implies

$$
\begin{aligned}
-\beta \log \beta \int_{0}^{\infty} e^{-\beta t} P_{t}^{\mu} V(x) d t & =-\beta \log \beta \int_{0}^{\infty} e^{-s} P_{\beta^{-1} s}^{\mu} V(x) \cdot \beta^{-1} d s \\
& =\int_{0}^{\infty} e^{-s} d v_{\beta}^{x}(s) \rightarrow C_{1} h_{0}(x) . \quad(\beta \rightarrow 0)
\end{aligned}
$$

Moreover we see that the $C_{1} h_{0}(x)=C_{1} h_{0}(x) \int_{0}^{\infty} e^{-s} v(d s)$ and thus it follows that

$$
\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-s} d v_{\beta}^{x}(s)=C_{1} h_{0}(x) \int_{0}^{\infty} e^{-s} d v(s)
$$

For integer $n$, we consider

$$
\begin{align*}
\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-n s} d v_{\beta}^{x}(s) & =\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-n s} \cdot(-\beta \log \beta) P_{\beta^{-1} s}^{\mu} V(x) \cdot \beta^{-1} d s \\
& =\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-t} \cdot(-\beta \log \beta) P_{(n \beta)^{-1} t}^{\mu} V(x) \cdot(n \beta)^{-1} d t \tag{3.26}
\end{align*}
$$

Noting that

$$
\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-t} \cdot(-n \beta \log n \beta) P_{(n \beta)^{-1} t}^{\mu} V(x) \cdot(n \beta)^{-1} d t=\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-t} d v_{n \beta}^{x}(t)=C_{1} h_{0}(x)
$$

we obtain

$$
\lim _{\beta \rightarrow 0} \int_{0}^{\infty} e^{-t} \cdot(-\beta \log \beta) P_{(n \beta)^{-1} t}^{\mu} V(x) \cdot(n \beta)^{-1} d t=\frac{C_{1} h_{0}(x)}{n}=C_{1} h_{0}(x) \int_{0}^{\infty} e^{-n s} d v(s) .
$$

Thus, for $f \in\left\{e^{-n s}\right\}_{n \in \mathbb{N}}$, it follows that

$$
\begin{equation*}
\int_{0}^{\infty} f(s) d v_{\beta}^{x}(s) \rightarrow C_{1} h_{0}(x) \int_{0}^{\infty} f(s) d v(s) \quad(\beta \rightarrow 0) \tag{3.27}
\end{equation*}
$$

Since polynomials in $\left\{e^{-n s}\right\}_{n \in \mathbb{N}}$ are dense in $C_{\infty}\left(\mathbb{R}_{+}\right)$by the Stone-Weierstrass theorem, (3.27) is valid for $f \in C_{\infty}\left(\mathbb{R}_{+}\right)$. Finally we prove (3.27) for $f(s)=1_{[0,1]}(s)$,
namely the main assertion $v_{\beta}^{x}[0,1] \rightarrow C_{1} h_{0}(x) v[0,1]$. Here $1_{[0,1]}(s)$ is the indicator function of the set $[0,1]$. Set the sequences of functions $\left\{f_{n}\right\}_{n \geq 1}$ and $\left\{g_{n}\right\}_{n \geq 1}$ by

$$
\begin{aligned}
& f_{n}(s)= \begin{cases}1 & (0 \leq s \leq 1) \\
1-n(s-1) \quad(1 \leq s \leq 1+1 / n) \\
0 & (1+1 / n \leq s),\end{cases} \\
& g_{n}(s)= \begin{cases}1 & (0 \leq s \leq 1-1 / n) \\
n(1-s) \quad(1-1 / n \leq s \leq 1) \\
0 & (1 \leq s)\end{cases}
\end{aligned}
$$

Note that these functions satisfies

$$
\begin{aligned}
0 & \leq g_{n}(s) \leq 1_{[0,1]}(s) \leq f_{n}(s) \leq 1, \quad\left(s \in \mathbb{R}_{+}\right) \\
\lim _{n \rightarrow \infty} f_{n}(s) & =1_{[0,1]}(s) \quad\left(s \in \mathbb{R}_{+}\right) \\
\lim _{n \rightarrow \infty} g_{n}(s) & =1_{[0,1]}(s) \quad\left(s \in \mathbb{R}_{+} \backslash\{1\}\right) .
\end{aligned}
$$

Then it follows that

$$
v_{\beta}^{x}[0,1]=\int_{\mathbb{R}_{+}} 1_{[0,1]}(s) v_{\beta}^{x}(d s) \leq \int_{\mathbb{R}_{+}} f_{n}(s) v_{\beta}^{x}(d s)
$$

Noting that $f_{n} \in C_{\infty}\left(\mathbb{R}_{+}\right)$and $\lim _{\beta \rightarrow 0} \int_{0}^{\infty} f_{n} d v_{\beta}^{x}=C_{1} h_{0}(x) \int_{0}^{\infty} f_{n} d v$, we obtain for arbitrary $n \in \mathbb{N}$

$$
\begin{equation*}
\limsup _{\beta \rightarrow 0} v_{\beta}^{x}[0,1] \leq C_{1} h_{0}(x) \int_{0}^{\infty} f_{n} d v \tag{3.28}
\end{equation*}
$$

Letting $n \rightarrow \infty$ in (3.28), the dominated convergence theorem implies the right hand side converges to $v[0,1]$ and thus we have

$$
\begin{equation*}
\limsup _{\beta \rightarrow 0} v_{\beta}^{x}[0,1] \leq C_{1} h_{0}(x) v[0,1] . \tag{3.29}
\end{equation*}
$$

Conversely, it follows that

$$
v_{\beta}^{x}[0,1] \geq \int_{0}^{\infty} g_{n}(s) v_{\beta}^{x}(d s)
$$

and thus we obtain

$$
\begin{equation*}
\liminf _{\beta \rightarrow 0} v_{\beta}^{x}[0,1] \geq C_{1} h_{0}(x) \int_{0}^{\infty} g_{n}(s) v(d s) . \tag{3.30}
\end{equation*}
$$

Note that $\lim _{n \rightarrow \infty} g_{n}(s)=1_{[0,1]}(s)$ for $s \neq 1$. However, using $v(\{1\})=0$, this convergence is valid for $v$-a.e. $s$. Thus, applying the dominated convergence theorem in the right hand side of (3.30), we have

$$
\begin{equation*}
\liminf _{\beta \rightarrow 0} v_{\beta}^{x}[0,1] \geq C_{1} h_{0}(x) v[0,1] . \tag{3.31}
\end{equation*}
$$

Combining (3.29) and (3.31), we have $\lim _{\beta \rightarrow 0} v_{\beta}^{x}[0,1]=C_{1} h_{0}(x) v[0,1]$.

Remark 3.1. (i) Theorem 3.3 is an extension of the result for 4-dimensional Brownian motion. The result by Takeda [34] is an extension of d-dimensional Brownian motion for $d \geq 5$. Thus the growth of $\mathbb{E}_{x}\left[\exp \left(A_{t}^{\mu}\right)\right]$ seems to depend on $d / \alpha$.
(ii) Since we use an analytic method in the proof, we assume the absolute continuity of $\mu$. However, Takeda [34] treated general measures. Thus, it seems to extend the result to the case of general measures.

### 3.3 Differentiability of spectral functions

In the previous section we define the spectral bound of $\mu$ by formula (3.14). For $\lambda \geq 0$, we define the spectral function $C(\lambda)$ by

$$
\begin{equation*}
C(\lambda)=-\inf \left\{\mathscr{E}(u, u)-\lambda \int_{\mathbb{R}^{d}} u^{2} d \mu \mid \int_{\mathbb{R}^{d}} u^{2}(x) d x=1\right\} \tag{3.32}
\end{equation*}
$$

Suppose $\mu$ is critical and $\mu=V \cdot m$ for $V \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$. Lemma 3.6 shows that (3.22) is valid for $e_{\beta}=\lambda^{-1}$ and $\beta=C(\lambda)$ for $\lambda \geq 1$. Thus, if we apply the inverse function theorem, we can determine the behavior of the spectral function. In [35], Takeda and Tsuchida gave the criterion for the differentiability of the spectral function as follows:

Theorem 3.4. Let $\left\{X_{t}\right\}_{t \geq 0}$ be the rotationally invariant $\alpha$-stable process on $\mathbb{R}^{d}$ and $\mu$ be a Green-tight, smooth measure in the Kato class. The spectral function $C(\lambda)$ is differentiable if and only if $1<d / \alpha \leq 2$.

We justify their result considering the concrete behavior of $C(\lambda)$ in the case of $d=2, \alpha=1$.
Theorem 3.5. The spectral function $C(\boldsymbol{\lambda})$ satisfies

$$
\begin{equation*}
C(\lambda) \asymp c \frac{1-\lambda^{-1}}{-\log \left(1-\lambda^{-1}\right)} \quad \lambda \rightarrow 1 . \tag{3.33}
\end{equation*}
$$

In particular, it follows

$$
\lim _{\lambda \rightarrow 1} \frac{1-\lambda^{-1}}{-(\lambda-1) \log \left(1-\lambda^{-1}\right)}=\lim _{\lambda \rightarrow 1} \frac{\lambda^{-1}}{\log \left(1-\lambda^{-1}\right)}=0
$$

and hence $C(\lambda)$ is differentiable at $\lambda=1$, which has consistency with Theorem 3.4.
In order to prove this theorem, we begin with the inverse function theorem containing logarithm.

Proposition 3.5. Let $f$ and $g$ be analytic function near $\xi=0$ and real-valued for $\xi$ real. Suppose $f(0)=f^{\prime}(0)=0$ and $g(0)=g^{\prime}(0)=0$ and $g^{\prime \prime}(0)<0$ and $\eta=f(\xi)+$ $g(\xi) \log \xi$. Then $\xi(\eta)$ satisfies

$$
\begin{gathered}
\xi(\eta)=\sum_{n \geq 1, m, k \geq 0} C_{n, m, k} \sigma^{n} \tau^{m} \omega^{k}, \\
\sigma=\left(\frac{\eta}{\log \eta^{-1}}\right)^{\frac{1}{2}}, \quad \tau=\frac{1}{\log \eta^{-1}}, \quad \omega=\frac{\log \left(\log \eta^{-1}\right)}{\log \eta^{-1}}
\end{gathered}
$$

Proof. Without loss of generality we may assume that

$$
f(\xi)=\sum_{n=2}^{\infty} a_{n} \xi^{n}, \quad g(\xi)=-2 \xi^{2} \log \xi+\log \xi \sum_{n=3}^{\infty} b_{n} \xi^{n}
$$

It follows that

$$
\begin{equation*}
\eta=\sum_{n=2}^{\infty} a_{n} \xi^{n}+\log \xi \sum_{n=3}^{\infty} b_{n} \xi^{n} \tag{3.34}
\end{equation*}
$$

and substitute $\xi=\sigma(1+z)$. (3.34) is rewritten as follows:

$$
\begin{align*}
\eta & =\sum_{n=2}^{\infty} a_{n} \sigma^{n}(1+z)^{n}-2 \sigma^{2}(1+z)^{2} \log (\sigma(1+z))+\log (1+z) \sum_{n=3}^{\infty} b_{n} \sigma^{n}(1+z)^{n} \\
& =\sigma^{2}\left(\sum_{n=0}^{\infty} a_{n+2} \sigma^{n}(1+z)^{n+2}-2(1+z)^{2} \log (\sigma(1+z))\right. \\
& \left.+\sigma \log (\sigma(1+z)) \sum_{n=0}^{\infty} b_{n+3} \sigma^{n-3}(1+z)^{n}\right) \tag{3.35}
\end{align*}
$$

Noting that $\sigma^{2}=\eta /\left(\log \eta^{-1}\right)$, we have

$$
\begin{align*}
\log \eta^{-1} & =\sum_{n=0}^{\infty} a_{n+2} \sigma^{n}(1+z)^{n+2}-2(1+z)^{2} \log (\sigma(1+z)) \\
& +\sigma \log (\sigma(1+z)) \sum_{n=0}^{\infty} b_{n+3} \sigma^{n-3}(1+z)^{n} \tag{3.36}
\end{align*}
$$

Multiplying $\tau$ in the both side of (3.36), we obtain

$$
\begin{align*}
1 & =\tau \sum_{n=0}^{\infty} a_{n+2} \sigma^{n}(1+z)^{n+2}-2(1+z)^{2} \tau \log (\sigma(1+z)) \\
& +\sigma \tau \log (\sigma(1+z)) \sum_{n=0}^{\infty} b_{n+3} \sigma^{n}(1+z)^{n+3} \tag{3.37}
\end{align*}
$$

$\log (\sigma(1+z))$ is rewritten as follows:

$$
\begin{aligned}
\log \sigma+\log (1+z) & =\frac{1}{2}\left(\log \eta-\log \left(\log \eta^{-1}\right)\right)+\log (1+z) \\
& =-\frac{1+\omega}{2 \tau}+\log (1+z)
\end{aligned}
$$

Thus, we obtain

$$
\begin{aligned}
1 & =\tau \sum_{n=0}^{\infty} a_{n+2} \sigma^{n}(1+z)^{n+2}+(1+z)^{2}(1+\omega) \\
& -2(1+z)^{2} \tau \log (1+z)-\left(\frac{\sigma}{2}+\frac{\sigma \omega}{2}-\sigma \tau \log (1+z)\right) \sum_{n=0}^{\infty} b_{n+3} \sigma^{n}(1+z)^{n+3}
\end{aligned}
$$

Set $F(z, \sigma, \tau, \omega)$ by

$$
\begin{aligned}
& F(z, \sigma, \tau, \omega)=z^{2}+2 z+(1+z)^{2} \omega-2 \tau(1+z)^{2} \log (1+z) \\
& \quad+\tau \sum_{n=0}^{\infty} a_{n+2} \sigma^{n}(1+z)^{n+2}-\left(\frac{\sigma}{2}+\frac{\sigma \omega}{2}-\sigma \tau \log (1+z)\right) \sum_{n=0}^{\infty} b_{n+3} \sigma^{n}(1+z)^{n+3}
\end{aligned}
$$

Then we see that $F(z, \sigma, \tau, \omega)=0$ and $\partial F / \partial z(0,0,0,0) \neq 0$, and thus we conclude

$$
z=\sum_{n, m, k \geq 0} d_{n, m, k} \sigma^{n} \tau^{m} \omega^{k}
$$

for $z$ sufficiently near 0 . Noting that $\xi=\sigma(1+z)$, we have the desired result.
Now we prove Theorem 3.5.
Proof of Theorem 3.5 In order to apply the inverse function theorem, we need to substitute $\beta=\xi^{2}$ and $e_{\beta}=\eta$ in (3.22). Thus we obtain

$$
\eta=1+C_{1} \xi^{2} \log \xi+C_{2} \xi^{2}+\ldots
$$

and $C_{1}>0$. Applying Proposition 3.5, we obtain the expansion of $\xi$ as follows:

$$
\begin{equation*}
\xi=\sum_{n \geq 1, m, k \geq 0} d_{n, m, k} \sigma^{n} \tau^{m} \omega^{k} \tag{3.38}
\end{equation*}
$$

where

$$
\sigma=\left(\frac{1-\eta}{\log (1-\eta)^{-1}}\right)^{\frac{1}{2}} \quad \tau=\frac{1}{\log (1-\eta)^{-1}} \quad \omega=\frac{\log \left(\log (1-\eta)^{-1}\right)}{\log (1-\eta)^{-1}}
$$

Since $\eta \rightarrow 1$ as $\xi \rightarrow 0$, the principal term of (3.38) is $d_{1,0,0} \sigma$. Noting that $\eta=e_{\beta}=\lambda^{-1}$ and $C(\lambda)=\xi^{2}$, we obtain the desired result.

## Chapter 4

## Continuity of harmonic functions for non-local Markov generators

We treated harmonic functions for Schrödinger operator in Chapters 2 and 3. Here we consider the harmonic function for Markov generator $\mathscr{L}$. We assume that $\mathscr{L}$ martingale problem is well-posed. We denote the associated jump process by $\left\{X_{t}\right\}_{t \geq 0}$. That is to say, the following two conditions are satisfied;

- For all $x \in \mathbb{R}^{d}, \mathbb{P}_{x}\left(X_{0}=x\right)=1$.
- For all $f \in C_{b}^{2},\left\{f\left(X_{t}\right)-f(x)-\int_{0}^{t} \mathscr{L} f\left(X_{s}\right) d s\right\}_{t \geq 0}$ is a $\mathbb{P}_{x}$-martingale.

We define $\mathscr{L}$-harmonic function as a bounded function $u$ such that $\left\{u\left(X_{t}\right)\right\}_{t \geq 0}$ is a martingale. Note that $u$ satisfies $\mathscr{L} u=0$. If $\mathscr{L}$ is a uniform elliptic second order operator, $\left\{X_{t}\right\}_{t \geq 0}$ is the diffusion process and harmonic function $u$ is continuous. In this chapter, we consider the same problem for non-local Markov generators. Thus, $\mathscr{L}$ is given by

$$
\mathscr{L} u(x)=\int_{\mathbb{R}^{d} \backslash\{0\}}\left(u(x+h)-u(x)-h \cdot \nabla u(x) 1_{\{|h| \leq 1\}}\right) n(x, h) d h,
$$

where $n(x, h)$ is a non-negative measurable function on $\mathbb{R}^{d} \times \mathbb{R}^{d}$. The function $n(x, h)$ is called a jump measure or a jump density from $x$ to $x+h$. When $n(x, h)$ satisfies

$$
\frac{C_{1}}{|h|^{d+\alpha}} \leq n(x, h) \leq \frac{C_{2}}{|h|^{d+\alpha}}
$$

for some constants $0<C_{1} \leq C_{2}$ and $0<\alpha<2$, the associated jump process is called the $\alpha$-stable-like process in Chen and Kumagai [9]. Bass and Levin [5] is one of the earliest papers which deals with this class of jump processes with non-smooth $n(x, h)$. They establish elliptic Harnack inequalities and the Hölder continuity for harmonic
functions. In [26], Song and Vondraček consider a wider class of jump processes, such as a sum of symmetric stable processes with different order.

In general, the exponent $\alpha$ depends on $x$ and $h$, and the associated operator is of variable order. One of the difficulties of variable order case is that we cannot use scaling property, unlike the stable-like case. It is a very delicate problem whether harmonic functions for operators of variable order are continuous or not. Indeed, Barlow and coauthors [3] prove that harmonic functions are not necessarily continuous in general. To guarantee the continuity of harmonic functions, we impose some conditions on $n(x, h)$. Bass and Kassmann [4] established sufficient conditions for the continuity of harmonic functions. These conditions are divided to two parts:
(SSJ) Singularity of small jumps, i.e. how the amount of jumps with size $r$ grows as $r$ tends to 0 .
(QRI) Quasi rotationally invariance, i.e. how the process jumps in any direction to some extent.

Husseini and Kassmann [14] treated weaker conditions than those of Bass and Kassmann [4]. In this chapter we reconsider conditions (SSJ) and (QRI) to extend their results. This chapter is organized as follows: In Section 1, we review a main result by Bass and Kassmann [4] which describes the sufficient condition for Hölder continuity of harmonic functions. We refer where the conditions (SSJ) and (QRI) contributes to the Hölder continuity. In Section 2, we consider the condition (SSJ) and extend the result of Bass and Kassmann [4]. In Section 3, we weaken the condition (QRI) and prove the continuity of harmonic functions.

### 4.1 Preceding results on continuous harmonic functions

Let $\mathscr{L}$ be a non-local operator of form (4.1) given by

$$
\begin{equation*}
\mathscr{L} u(x)=\int_{\mathbb{R}^{d} \backslash\{0\}}\left(u(x+h)-u(x)-h \cdot \nabla u(x) 1_{\{|h| \leq 1\}}\right) n(x, h) d h, \tag{4.1}
\end{equation*}
$$

and suppose the jump measure $n(x, h)$ satisfies

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d}}\left(|h|^{2} \wedge 1\right) n(x, h) d h<\infty \tag{4.2}
\end{equation*}
$$

Let $A$ be a Borel set. For the process $\left\{X_{t}\right\}_{t \geq 0}$, denote the first exiting time of $A$ by $\tau_{A}$ and the first hitting time of $A$ by $T_{A}$ respectively; i.e.

$$
\tau_{A}:=\inf \left\{t>0: X_{t} \notin A\right\}, \quad T_{A}:=\inf \left\{t>0: X_{t} \in A\right\}
$$

Now we give the definition of harmonic functions.
Definition 4.1. The bounded function $u$ on $\mathbb{R}^{d}$ is called harmonic with respect to $\mathscr{L}$ on $D \subset \mathbb{R}^{d}$ if $\left\{u\left(X_{t \wedge \tau_{D}}\right)\right\}_{t \geq 0}$ is a martingale.

Roughly speaking, a harmonic function $u$ satisfies $\mathscr{L} u(x)=0$ for $x \in D$.

The following definition of $S(x, r), L(x, r)$ and $N(x, r)$ are taken from Bass and Kassmann [4].

$$
\begin{aligned}
& S(x, r)=\int_{|h| \geq r} n(x, h) d h, \\
& L(x, r)=S(x, r)+\frac{1}{r}\left|\int_{r \leq|h| \leq 1} h n(x, h) d h\right|+\frac{1}{r^{2}} \int_{|h| \leq r}|h|^{2} n(x, h) d h \\
& N(x, r)=\inf \left\{\int_{h \in A-x} n(x, h) d h: A \subset B(x, 2 r) \text { and }|A| \geq \frac{|B(x, r)|}{3 \cdot 2^{d}}\right\},
\end{aligned}
$$

where $B(x, r)$ and $|A|$ stand for an open ball in $\mathbb{R}^{d}$ centered at $x$ with radius $r$ and the Lebesgue measure of $A$ respectively. The following assumption is also taken from Bass and Kassmann [4].

Assumption 4.1. Let $R_{0}$ be a constant such that $0<R_{0}<1$. Suppose the following conditions hold.
(SSJ-1) There exist $\kappa_{1}>0$ and $\sigma>0$ such that for all $x \in \mathbb{R}^{d}, r \in\left(0, R_{0}\right)$ and $1<\lambda<1 / r$,

$$
\frac{S(x, \lambda r)}{S(x, r)} \leq \kappa_{1} \lambda^{-\sigma}
$$

(QRI-1) There exists $\kappa_{2}>0$ such that for all $x \in \mathbb{R}^{d}, r \in\left(0, R_{0}\right)$ and $y \in B(x, 2 r)$,

$$
N(x, r) \geq \kappa_{2} L(y, r / 2) .
$$

In [4], Bass and Kassmann prove the Hölder continuity of harmonic functions under the assumptions (SSJ-1) and (QRI-1). Precise statement of their theorem is the following.

Theorem 4.1. ([4] Theorem2.2.) Let $R_{0}$ be a constant such that $0<R_{0}<1$. Suppose $0<R<R_{0}$ and let $u$ be bounded in $\mathbb{R}^{d}$ and harmonic in $B\left(z_{0}, R\right)$ with respect to $\mathscr{L}$. Then, under the assumptions (SSJ-1) and (QRI-1), there exist $v \in(0,1)$ and $C>0$ depending only on $\kappa_{1}, \kappa_{2}$, and $\sigma$ such that

$$
\begin{equation*}
|u(z)-u(y)| \leq C\|u\|_{\infty}\left(\frac{|z-y|}{R}\right)^{v}, \quad z, y \in B\left(z_{0}, R / 3\right) \tag{4.3}
\end{equation*}
$$

Define $L_{0}\left(x_{0}, r\right)$ and $N_{0}\left(x_{0}, r\right)$ as follows:

$$
L_{0}\left(x_{0}, r\right):=\sup _{x \in B\left(x_{0}, r\right)} L(x, r) . \quad N_{0}\left(x_{0}, r\right):=\inf _{x \in B\left(x_{0}, r\right)} N(x, r)
$$

For the proof of Theorem 4.1, we begin with three lemmas.
Proposition 4.1. ([4] Proposition 3.1.) There exists $C_{1}$ such that

$$
\mathbb{P}_{x_{0}}\left(\tau_{B\left(x_{0}, r\right)}<t\right) \leq C_{1} t L_{0}\left(x_{0}, r\right) .
$$

Proof. Let $u$ be a $C^{2}\left(\mathbb{R}^{d}\right)$ function with bounded first and second partial derivatives that satisfies the following;

- $u(x)=\left|x-x_{0}\right|^{2} / r^{2}$ for $\left|x-x_{0}\right| \leq r$.
- $1 \leq u(x) \leq C_{2}$ for $\left|x-x_{0}\right|>r$.
- $|\nabla u| \leq C_{3} / r$.
- The second partial derivatives are bounded by $C_{4} / r^{2}$.

Since $\mathbb{P}_{x_{0}}$ is a solution to the martingale problem for $\mathscr{L}$ started $x_{0}$,

$$
\mathbb{E}_{x_{0}}\left[u\left(X_{\left.t \wedge \tau_{B\left(x_{0}, r\right)}\right)}\right)\right]-u\left(x_{0}\right)=\mathbb{E}_{x_{0}}\left[\int_{0}^{t \wedge \tau_{B\left(x_{0}, r\right)}} \mathscr{L} u\left(X_{s}\right) d s\right]
$$

The left hand side is larger than $\mathbb{P}_{x_{0}}\left(\tau_{B\left(x_{0}, r\right)} \leq t\right)$, while the right hand side is bounded from above by

$$
t \sup _{x \in B\left(x_{0}, r\right)}|\mathscr{L} u(x)| \leq\left(C_{2}+C_{3}+C_{4}\right) t L_{0}\left(x_{0}, r\right)
$$

We thus obtain the assertion.
Denote the left limit of $X_{t}$ by $X_{t-}$ and set $\Delta X_{t}:=X_{t}-X_{t-}$. Let $Q$ be a Borel set and $U:=\inf \left\{t:\left|\Delta X_{t}\right| \geq r\right\}$ for $r<R_{0}$. We then have the following.

Proposition 4.2. ([4] Proposition 3.3.) Under the condition (SSJ-1), we have

$$
\begin{equation*}
\mathbb{P}_{x}\left(\left|\Delta X_{U \wedge \tau_{Q}}\right| \geq \lambda r\right) \leq \kappa_{1} \lambda^{-\sigma} \tag{4.4}
\end{equation*}
$$

for $r<R_{0}, 1<\lambda<1 / r$ and $x \in \mathbb{R}^{d}$.
Proof. By the Lévy system formula, for every bounded stopping time $S$ and disjoint Borel sets $B$ and $C$, we have

$$
\begin{equation*}
\mathbb{E}_{y}\left[\sum_{s \leq S} 1_{\left(X_{s-} \in B, X_{s} \in C\right)}\right]=\mathbb{E}_{y}\left[\int_{0}^{S} 1_{B}\left(X_{s}\right) n\left(X_{s}, C-X_{s}\right) d s\right] . \tag{4.5}
\end{equation*}
$$

Using this formula and the condition (SSJ-1), for arbitrary $t>0$ we have

$$
\begin{align*}
\mathbb{P}_{x}\left(\left|\Delta X_{U \wedge \tau_{Q} \wedge t}\right|\right. & \geq \lambda r)=\mathbb{E}_{x}\left[\sum_{s \leq U \wedge \tau_{Q} \wedge t} 1_{\left(\left|\Delta X_{s}\right| \geq \lambda r\right)}\right] \\
& =\mathbb{E}_{x}\left[\int_{0}^{U \wedge \tau_{Q} \wedge t} \int_{|h| \geq \lambda r} n\left(X_{s}, h\right) d h d s\right] \\
& =\mathbb{E}_{x}\left[\int_{0}^{U \wedge \tau_{Q} \wedge t} S\left(X_{s}, \lambda r\right) d s\right] \leq \kappa_{1} \lambda^{-\sigma} \mathbb{E}_{x}\left[\int_{0}^{U \wedge \tau_{Q} \wedge t} S\left(X_{s}, r\right) d s\right] \tag{4.6}
\end{align*}
$$

By the Lévy system formula, the right hand side of (4.6) is estimated as follows:

$$
\kappa_{1} \lambda^{-\sigma} \mathbb{E}_{x}\left[\int_{0}^{U \wedge \tau_{Q} \wedge t} S\left(X_{s}, r\right) d s\right]=\kappa_{1} \lambda^{-\sigma} \mathbb{E}_{x}\left[\int_{0}^{U \wedge \tau_{Q} \wedge t} \int_{|h| \geq r} n\left(X_{s}, h\right) d h d s\right]
$$

$$
=\kappa_{1} \lambda^{-\sigma} \mathbb{E}_{x}\left[\sum_{s \leq U \wedge \tau_{Q} \wedge t} 1_{\left(\left|\Delta X_{s}\right| \geq r\right)}\right]=\kappa_{1} \lambda^{-\sigma} \mathbb{P}_{x}\left(\left|\Delta X_{U \wedge \tau_{Q} \wedge t}\right| \geq r\right) \leq \kappa_{1} \lambda^{-\sigma} .
$$

Thus, we obtain

$$
\mathbb{P}_{x}\left(\left|\Delta X_{U \wedge \tau_{Q} \wedge t}\right| \geq \lambda r\right) \leq \kappa_{1} \lambda^{-\sigma}
$$

Now letting $t \rightarrow \infty$ and using the dominated convergence theorem, the assertion (4.4) follows.

Proposition 4.3. ([4] Proposition 3.2.) Suppose the condition (QRI-1) holds. Suppose $0<r<R_{0}, A \subset B\left(x_{0}, r\right), y \in B\left(x_{0}, r / 2\right)$, and $|A| /\left|B\left(x_{0}, r\right)\right| \geq 1 /\left(3 \cdot 2^{d}\right)$. Then, there exists $\kappa_{5}>0$ not depending on $x_{0}, r$, or $A$ such that

$$
\mathbb{P}_{y}\left(T_{A}<\tau_{B\left(x_{0}, r\right)}\right) \geq \kappa_{5}
$$

Proof. Set $\tau=\tau_{B\left(x_{0}, r\right)}$. If $\mathbb{P}_{y}\left(T_{A}<\tau\right) \geq 1 / 4$, there is nothing to prove. Therefore, we may assume $\mathbb{P}_{y}\left(T_{A}<\tau\right) \leq 1 / 4$ without loss of generality. Note that $B\left(x_{0}, r\right) \supset$ $B(y, r / 2)$. Combining with Proposition 4.1, we have

$$
\mathbb{P}_{y}\left(\tau \leq t_{0}\right) \leq \mathbb{P}_{y}\left(\tau_{B(y, r / 2)} \leq t_{0}\right) \leq C_{1} t_{0} L_{0}(y, r / 2)
$$

Therefore, for sufficiently small $C_{2}$ if we set

$$
t_{0}=\frac{C_{2}}{L_{0}(y, r / 2)}
$$

then $\mathbb{P}_{y}\left(\tau \leq t_{0}\right) \leq 1 / 2$.
Moreover, if $x \in B\left(x_{0}, r\right)$, we have

$$
\begin{equation*}
n(x, A-x):=\int_{A-x} n(x, h) d h \geq N_{0}\left(x_{0}, r\right) \tag{4.7}
\end{equation*}
$$

Combining this estimate with the Lévy system formula (4.5), we obtain

$$
\begin{aligned}
\mathbb{P}_{y}\left(T_{A} \leq \tau\right) & \geq \mathbb{E}^{y}\left[\sum_{s \leq T_{A} \wedge \tau \wedge t_{0}} 1_{\left\{X_{s-} \neq X_{s}, X_{s} \in A\right\}}\right] \\
& \geq \mathbb{E}_{y}\left[\int_{0}^{T_{A} \wedge \tau t_{0}} n\left(X_{s}, A-X_{s}\right) d s\right] \geq N_{0}\left(x_{0}, r\right) \mathbb{E}_{y}\left(T_{A} \wedge \tau \wedge t_{0}\right)
\end{aligned}
$$

Furthermore, we have

$$
\begin{aligned}
\mathbb{E}_{y}\left(T_{A} \wedge \tau \wedge t_{0}\right) & \geq t_{0} \mathbb{P}_{y}\left(T_{A} \geq \tau \geq t_{0}\right) \\
& \geq t_{0}\left[1-\mathbb{P}_{y}\left(T_{A}<\tau\right)-\mathbb{P}_{y}\left(\tau<t_{0}\right)\right] \geq \frac{t_{0}}{4}
\end{aligned}
$$

Therefore, we conclude

$$
\mathbb{P}_{y}\left(T_{A}<\tau\right) \geq \frac{t_{0}}{4} N_{0}\left(x_{0}, r\right)=\frac{C_{2}}{4} \cdot \frac{N_{0}\left(x_{0}, r\right)}{L_{0}(y, r / 2)} \geq \frac{C_{2} \kappa_{2}}{4}
$$

and the assertion follows, where in the last inequality we used (QRI-1).

These three propositions are the keys to prove Theorem 4.1.
Proof of Theorem 4.1. Suppose $\|u\|_{\infty} \leq K$ and choose $0<\rho_{0}<\frac{1}{4}$ appropriately. Set

$$
r_{n}:=\frac{R}{12} \rho_{0}^{n-1} \quad(n \in \mathbb{N})
$$

Then, we have $B\left(z_{1}, 2 r_{1}\right) \subset B\left(z_{0}, R / 2\right)$ for arbitrary $z_{1} \in B\left(z_{0}, R / 3\right)$. Set $B_{n}, M_{n}$ and $m_{n}$ as follows.

$$
B_{n}:=B\left(z_{1}, r_{n}\right), \quad M_{n}:=\sup _{x \in B_{n}} u(x), \quad \quad m_{n}:=\inf _{x \in B_{n}} u(x) .
$$

Let $a<1$ and $\theta_{1}>2 K$ be positive constants chosen later appropriately. Set

$$
s_{n}=\theta_{1} a^{n} \quad(n \in \mathbb{N})
$$

We will prove by induction that

$$
\begin{equation*}
M_{n}-m_{n} \leq s_{n} \quad(\forall n \in \mathbb{N}) \tag{4.8}
\end{equation*}
$$

We assume that $M_{k}-m_{k} \leq s_{k}$ for $1 \leq k \leq n$ and consider the case of $k=n+1$. Choose $y, z \in B_{n+1}$ such that

$$
u(z) \geq M_{n+1}-\varepsilon, \quad u(y) \leq m_{n+1}+\varepsilon
$$

where $\varepsilon$ is an arbitrary positive number. Set

$$
A_{n}=\left\{z \in B_{n}: u(z) \leq\left(M_{n}+m_{n}\right) / 2\right\}
$$

Without loss of generality, we may assume $\left|A_{n}\right| /\left|B_{n}\right| \geq 1 / 2$. (If $\left|A_{n}\right| /\left|B_{n}\right|<1 / 2$, consider the function $-u(x)$ instead of $u(x)$.)
Choose a compact set $D_{n}$ such that

$$
D_{n} \subset A_{n}, \quad\left|D_{n}\right| /\left|B_{n}\right| \geq 1 / 3
$$

Note that Proposition 4.3 implies $\mathbb{P}_{z}\left(T_{D_{n}}<\tau_{B_{n}}\right) \geq \kappa_{5}$. Let $p_{n}:=\mathbb{P}_{z}\left(T_{D_{n}}<\tau_{B_{n}}\right)$ and define $\tau_{n}:=\tau_{B_{n}}$. Harmonic property of $u$ and the optional sampling theorem imply that

$$
\begin{align*}
u(z)-u(y) & =\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y)\right] \\
& =\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): T_{D_{n}}<\tau_{n}, X_{\tau_{n}} \in B_{n-1}\right] \\
& +\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): T_{D_{n}} \geq \tau_{n}, X_{\tau_{n}} \in B_{n-1}\right] \\
& +\sum_{i=1}^{n-2} \mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): X_{\tau_{n}} \in B_{n-i-1} \backslash B_{n-i}\right] \\
& +\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): X_{\tau_{n}} \notin B_{1}\right] \\
& =: I_{1}+I_{2}+I_{3}+I_{4} . \tag{4.9}
\end{align*}
$$

Here we define $u\left(X_{\infty}\right):=\lim _{t \rightarrow \infty} u\left(X_{t}\right)$ and assume $X_{\tau_{n}} \in B_{n-1}$ when $\tau_{n}=\infty$. Now we estimate each term of (4.9). First, we consider $I_{1}$ and $I_{2}$.

$$
\begin{align*}
& I_{1} \leq \frac{M_{n}-m_{n}}{2} \mathbb{P}_{z}\left(T_{D_{n}}<\tau_{n}\right) \leq \frac{s_{n}}{2} p_{n}  \tag{4.10}\\
& I_{2} \leq\left(M_{n-1}-m_{n-1}\right) \mathbb{P}_{z}\left(\tau_{n} \leq T_{D_{n}}\right) \leq s_{n-1}\left(1-p_{n}\right) \tag{4.11}
\end{align*}
$$

In order to estimate $I_{3}$ and $I_{4}$, we first define

$$
F_{i}:=\mathbb{P}_{z}\left(X_{\tau_{n}} \notin B_{n-i}\right)
$$

and consider this value. If $X_{\tau_{n}} \notin B_{n-i}$, the process cannot have a jump larger than $2 r_{n}$ before the time $\tau_{n}$ and $\left|X_{\tau_{n}}-X_{\tau_{n}-}\right| \geq r_{n-i}-r_{n}$. Note that $\rho_{0}<\frac{1}{4}$ and therefore $2 r_{n} \leq r_{n-i}-r_{n}$. Then we can apply Proposition 4.2 and conclude that

$$
F_{i} \leq \kappa_{1}\left(\frac{2 r_{n}}{r_{n-i}-r_{n}}\right)^{\sigma}=\kappa_{1}\left(\frac{2 \rho_{0}^{i}}{1-\rho_{0}^{i}}\right)^{\sigma} \leq C_{2} \rho_{0}^{i \sigma}
$$

Then we can obtain the following estimate;

$$
\begin{align*}
I_{3} & \leq \sum_{i=1}^{n-2}\left(M_{n-i-1}-m_{n-i-1}\right)\left(F_{i}-F_{i+1}\right) \leq \sum_{i=1}^{n-2} s_{n-i-1}\left(F_{i}-F_{i+1}\right) \\
& =s_{n-2} F_{1}-s_{1} F_{n-1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i} \leq s_{n-2} F_{1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i} \\
& \leq s_{n-2} C_{2} \rho_{0}^{\sigma}+s_{n} C_{2}\left(a^{-1}-1\right) \sum_{i=2}^{n-2}\left(\frac{\rho_{0}^{\sigma}}{a}\right)^{i}  \tag{4.12}\\
I_{4} & \leq 2 K \cdot F_{n-1} \leq \theta_{1} C_{2} \rho_{0}^{(n-1) \sigma} \tag{4.13}
\end{align*}
$$

Combining the estimates (4.10)-(4.13), we conclude that

$$
\begin{align*}
u(z)-u(y) & \leq \frac{1}{2} s_{n} p_{n}+s_{n-1}\left(1-p_{n}\right)+s_{n-2} F_{1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i}+2 K \cdot F_{n-1} \\
& \leq \frac{1}{2} s_{n} p_{n}+s_{n-1}\left(1-p_{n}\right)+s_{n-2} C_{2} \rho_{0}^{\sigma} \\
& +s_{n} C_{2}\left(a^{-1}-1\right) \sum_{i=2}^{n-2}\left(\frac{\rho_{0}^{\sigma}}{a}\right)^{i}+\theta_{1} C_{2} \rho_{0}^{(n-1) \sigma} \tag{4.14}
\end{align*}
$$

In the rest of the proof, we assume $n$ is large enough. Choose $a$ such that $\rho_{0}^{\frac{\sigma}{2}}<a<1$, then we have

$$
\sum_{i=2}^{n-2}\left(\frac{\rho_{0}^{\sigma}}{a}\right)^{i} \leq C_{3}, \quad \text { and } \quad \theta_{1} \rho_{0}^{(n-1) \sigma} \leq C_{4} s_{n} \rho_{0}^{n \sigma / 3}
$$

Moreover, recall that there exists $\kappa_{5}>0$ such that $p_{n} \geq \kappa_{5}$ for arbitrary $n \in \mathbb{N}$. Combining these facts, we obtain that
R.H.S. of $(4.14) \leq s_{n}\left(a^{-1}-\frac{\kappa_{5}}{2}+a^{-2} C_{2} \rho_{0}^{\sigma}+C_{5}\left(a^{-1}-1\right)+C_{6} \rho_{0}^{n \sigma / 3}\right)$

$$
\begin{equation*}
\leq s_{n}\left(1-\frac{\kappa_{5}}{2}+a^{-2} C_{2} \rho_{0}^{\sigma}+C_{7}\left(a^{-1}-1\right)+C_{6} \rho_{0}^{n \sigma / 3}\right) \tag{4.15}
\end{equation*}
$$

If we choose $a$ sufficiently close to 1 and $\rho_{0}$ sufficiently small, there exists $n_{0} \in \mathbb{N}$ such that for arbitrary $n \geq n_{0}$,

$$
\begin{equation*}
\text { R.H.S. of }(4.15) \leq s_{n}\left(1-\frac{\kappa_{5}}{4}\right) \leq s_{n+1} \text {. } \tag{4.16}
\end{equation*}
$$

Furthermore, if we choose $\theta_{1}$ such that $\theta_{1} a^{n_{0}} \geq 2 K$, then the induction hypothesis is satisfied for $1 \leq n \leq n_{0}$. Therefore, we have proved (4.8) and this implies the Hölder continuity of $u$. Indeed, set $z=z_{1}$ and choose $y \in B\left(z_{0}, \frac{R}{3}\right)$. If there exists $n \in \mathbb{N}$ such that

$$
\begin{equation*}
\frac{R}{12} \rho_{0}^{n} \leq|z-y|<\frac{R}{12} \rho_{0}^{n-1} \tag{4.17}
\end{equation*}
$$

then, we have

$$
|u(z)-u(y)| \leq \theta_{1} a^{n}=C_{8} K\left(\exp \left(n \log \rho_{0}\right)\right)^{\frac{\log a}{\log \rho_{0}}} \leq C_{9}\|u\|_{\infty}\left(\frac{|z-y|}{R}\right)^{\frac{\log a}{\log \rho_{0}}}
$$

If there is no $n$ that satisfies (4.17), then $|z-y| \geq R / 12$. Therefore, there exists $C_{10}>0$ such that

$$
|u(z)-u(y)| \leq 2\|u\|_{\infty} \leq C_{10}\|u\|_{\infty}\left(\frac{|z-y|}{R}\right)^{\frac{\log a}{\log \rho_{0}}}
$$

Noting that $z=z_{1} \in B\left(z_{0}, R / 3\right)$ is arbitrary, this completes the proof of the Hölder continuity of $u$ in $B\left(z_{0}, R / 3\right)$.

Remark 4.1. There is a minor error in the proof of Theorem 2.2 in [4] by Bass and Kassmann. One of the keys to prove this theorem is the oscillation arguments; i.e. we choose appropriate sequences $\left\{r_{n}\right\}$ and $\left\{s_{n}\right\}$ decreasing to 0 as $n \rightarrow \infty$ and show that

$$
\begin{equation*}
|u(z)-u(y)| \leq s_{n} \quad\left(|z-y| \leq r_{n}\right) \tag{4.18}
\end{equation*}
$$

by induction. Where the choice $r_{n}=\theta_{2} / 4^{n}$ is specified, it is necessary to take $r_{n}=\theta_{2} \rho^{n}$ for some $0<\rho<1 / 4$. They choose $r_{n}=\theta_{2} / 4^{n}$, where $\theta_{2}$ is a sufficiently small positive constant. With this choice, the proof does not necessarily work well. We may resolve this by choice $r_{n}=\theta_{2} \rho^{n}$ with appropriate $0<\rho<1 / 4$.

### 4.2 Condition (SSJ) and continuous harmonic functions

We consider the following assumption instead of (SSJ-1);
(SSJ-2) There exist $\kappa_{3}>0$ and $\gamma>1$ such that for all $x \in \mathbb{R}^{d}, r \in\left(0, R_{0}\right)$ and $1<\lambda<1 / r$,

$$
\frac{S(x, \lambda r)}{S(x, r)} \leq \kappa_{3}(\log \lambda)^{-\gamma}
$$

It is clear that (SSJ-1) implies (SSJ-2). In this section, we will prove the equivalence of the conditions (SSJ-1) and (SSJ-2) with $\gamma>0$. The following proposition plays a crucial role in proving this equivalence.
Theorem 4.2. Let $f:(0, \infty) \rightarrow(0, \infty)$ be a non-increasing function satisfying

$$
\begin{equation*}
\frac{f(\lambda r)}{f(r)} \leq g(\lambda) \quad \text { for all } \quad r \in(0,1) \quad \text { and } \quad \lambda \in\left(1, \frac{1}{r}\right) \tag{4.19}
\end{equation*}
$$

for a function $g:(1, \infty) \rightarrow(0, \infty)$ such that $\liminf _{\lambda \rightarrow \infty} g(\lambda)<1$. There exist constants $c>0, \sigma>0$ such that (4.19) holds for $g(\lambda)=c \lambda^{-\sigma}$, i.e.

$$
\frac{f(\lambda r)}{f(r)} \leq c \lambda^{-\sigma} \quad \text { for all } \quad r \in(0,1) \quad \text { and } \quad \lambda \in\left(1, \frac{1}{r}\right)
$$

Proof. Let $r \in(0,1)$ and $\lambda \in(1,1 / r)$. Choose $b>1$ so that $g(b)<1$. Let $n \in \mathbb{Z}_{+}$be such that $b^{n} \leq \lambda<b^{n+1}$. Assume first that $n \geq 1$.
Since $b^{n}<1 / r$, we see that $b^{n-1} r<b^{-1}<1$ and so we can apply (4.19) with $b^{n-1} r \in$ $(0,1)$ and $b \in\left(1,1 / b^{n-1} r\right)$ to get

$$
\begin{equation*}
\frac{f\left(b^{n} r\right)}{f\left(b^{n-1} r\right)}=\frac{f\left(b \cdot b^{n-1} r\right)}{f\left(b^{n-1} r\right)} \leq g(b) \tag{4.20}
\end{equation*}
$$

Iterating (4.20), it follows that

$$
\begin{aligned}
\frac{f\left(b^{n} r\right)}{f(r)} & \leq g(b)^{n}=g(b) \exp ((n-1) \log g(b))=g(b) \exp \left((n-1) \log b^{-1} \cdot \frac{\log g(b)}{\log b^{-1}}\right) \\
& =g(b) \cdot\left(b^{-(n-1) \sigma}\right)=g(b) \cdot b^{2 \sigma} \cdot b^{-(n+1) \sigma}
\end{aligned}
$$

where $\sigma:=-\log g(b) / \log b>0$. Since $f$ is non-increasing, we finally obtain

$$
\frac{f(\lambda r)}{f(r)} \leq \frac{f\left(b^{n} r\right)}{f(r)} \leq b^{2 \sigma} g(b) b^{-(n+1) \sigma} \leq b^{2 \sigma} g(b) \lambda^{-\sigma}
$$

The case $n=0$ follows from the fact that $f$ is non-increasing:

$$
\frac{f(\lambda r)}{f(r)} \leq 1<b^{\sigma} \lambda^{-\sigma}
$$

Thus, it is enough to choose $c:=g(b) b^{2 \sigma} \vee b^{\sigma}$.
Applying Theorem 4.2 with $g(\lambda)=\kappa_{3}(\log \lambda)^{-\gamma}$, (SSJ-2) implies (SSJ-1) and thus we can see the equivalence of (SSJ-1) and (SSJ-2) with $\gamma>0$. Consequently, we can obtain the Hölder continuity of harmonic functions even if we change the condition (SSJ-1) into (SSJ-2) with $\gamma>0$. Thus, combining this equivalence and Theorem 4.1, we can obtain the following corollary.
Corollary 4.1. Suppose $0<R<R_{0}<1$ and let $u$ be bounded in $\mathbb{R}^{d}$ and harmonic in $B\left(z_{0}, R\right)$ with respect to $\mathscr{L}$. Then, under the assumptions (SSJ-2) with $\gamma>0$ and (QRI-1), there exist $v \in(0,1)$ and $C>0$ depending only on $R_{0}, \kappa_{2}, \kappa_{3}$ and $\gamma$ such that

$$
\begin{equation*}
|u(z)-u(y)| \leq C\|u\|_{\infty}\left(\frac{|z-y|}{R}\right)^{v}, \quad z, y \in B\left(z_{0}, R / 3\right) \tag{4.21}
\end{equation*}
$$

### 4.3 Condition (QRI) and continuous harmonic functions

We next consider the following condition which is weaker than (QRI-1);
(QRI-2) There exists $\kappa_{4}>0$ such that for all $x \in \mathbb{R}^{d}, r \in\left(0, R_{0}\right)$ and $y \in B(x, 2 r)$,

$$
N(x, r) \geq \frac{\kappa_{4}}{|\log r|} L(y, r / 2) .
$$

If we change the condition (QRI-1) into (QRI-2), we can prove the uniform continuity of harmonic functions. The precise statement is as follows;

Theorem 4.3. Let $0<R<R_{0}<1$ and suppose (SSJ-1) or (SSJ-2) with $\gamma>0$ and in addition (QRI-2). Then there exists a monotone increasing continuous function $P$ : $\mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$with $P(0)=0$ such that for any $\mathscr{L}$-harmonic function on $B\left(z_{0}, R\right)$

$$
|u(z)-u(y)| \leq\|u\|_{\infty} P(|z-y|), \quad z, y \in B\left(z_{0}, R / 2\right) .
$$

Although Husseini and Kassmann [14] also obtain similar theorem under the assumptions (SSJ-1) and (QRI-2), it seems that there is a gap in their proof. (See Remark 4.3 for details.) Thus, we will give a full proof there. We first consider an example where the assumption (SSJ-1) and (QRI-2) are satisfied.

Example 4.1. Let $\kappa_{0}$ be a given positive constant. If the jump density function $n(x, h)$ satisfies for $0<|h| \leq 2$

$$
\begin{equation*}
\frac{\kappa_{0}}{|h|^{d+\alpha}} \log \frac{3}{|h|} \leq n(x, h) \leq \frac{\kappa_{0}-1}{|h|^{d+\alpha}} \log \frac{3}{|h|} \tag{4.22}
\end{equation*}
$$

with $1 \leq \alpha<2$ and $\sup _{x \in \mathbb{R}^{d}} S(x, 1)<\infty$, then the uniform continuity of harmonic function holds.

Proof. Due to Theorem 4.3, it is enough to verify (4.2), (SSJ-1) and (QRI-2). (4.2) is easy to verify. By (4.22), we have

$$
\begin{aligned}
\int_{|h| \leq 1}|h|^{2} n(x, h) d h & \leq \int_{|h| \leq 1}|h|^{2} \cdot \frac{\kappa_{0}^{-1}}{|h|^{d+\alpha}} \log \frac{3}{|h|} d h \\
& =C_{1} \int_{0}^{1} \frac{1}{\psi^{\alpha-1}} \log \frac{3}{\psi} d \psi=C_{2} \int_{0}^{\infty}(\chi+\log 3) e^{(\alpha-2) \chi} d \chi<\infty
\end{aligned}
$$

where we use $\psi=e^{-\chi}$ in the last equality. Thus, combining with the assumption $\sup _{x \in \mathbb{R}^{d}} S(x, 1)<\infty$, we obtain

$$
\sup _{x \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d}}\left(|h|^{2} \wedge 1\right) n(x, h) d h \leq \sup _{x \in \mathbb{R}^{d}} S(x, 1)+\int_{|h| \leq 1} \frac{\kappa_{0}^{-1}}{|h|^{d+\alpha-2}} \log \frac{3}{|h|} d h<\infty .
$$

Let $r$ be $0<r<R_{0}<1$. To see (SSJ-1) follow, we begin with the estimate of

$$
\int_{r \leq|h| \leq 1} \frac{1}{|h|^{d+\alpha}} \log \frac{3}{|h|} d h .
$$

Using the spherical coordinate, we can obtain that

$$
\begin{aligned}
\int_{r \leq|h| \leq 1} \frac{1}{|h|^{d+\alpha}} \log \frac{3}{|h|} d h & =C_{2} \int_{r}^{1} \frac{1}{\psi^{\alpha+1}} \log \frac{3}{\psi} d \psi \\
& =C_{2}\left(\frac{1}{\alpha} r^{-\alpha} \log \frac{1}{r}+\frac{\alpha \log 3-1}{\alpha^{2}}\left(r^{-\alpha}-1\right)\right) \\
& =C_{3}\left(r^{-\alpha} \log \frac{1}{r}+\left(\log 3-\frac{1}{\alpha}\right)\left(r^{-\alpha}-1\right)\right) .
\end{aligned}
$$

We now estimate $S(x, r)$ from below. The above formula and (4.22) imply

$$
\begin{align*}
S(x, r) & =\int_{|h| \geq r} n(x, h) d h \geq \int_{r \leq|h| \leq 1} n(x, h) d h \\
& \geq C_{4}\left(r^{-\alpha} \log \frac{1}{r}+\left(\log 3-\frac{1}{\alpha}\right)\left(r^{-\alpha}-1\right)\right) \geq C_{4} r^{-\alpha} \log \frac{1}{r}, \tag{4.23}
\end{align*}
$$

where we used the assumption $1 \leq \alpha<2$ in the last inequality. Moreover, if we estimate $S(x, r)$ from above, we obtain that

$$
\begin{align*}
S(x, r)= & S(x, 1)+\int_{r \leq|h| \leq 1} n(x, h) d h \\
& \leq C_{5}+C_{6}\left(r^{-\alpha} \log \frac{1}{r}+\left(\log 3-\frac{1}{\alpha}\right)\left(r^{-\alpha}-1\right)\right) \\
& \leq C_{5}+C_{6} r^{-\alpha} \log \frac{3}{r} \leq C_{7} r^{-\alpha} \log \frac{3}{r} . \tag{4.24}
\end{align*}
$$

Combining (4.23) and (4.24), we see that

$$
\frac{S(x, \lambda r)}{S(x, r)} \leq \frac{C_{3}(\lambda r)^{-\alpha} \log (3 / \lambda r)}{C_{1} r^{-\alpha} \log (1 / r)} \leq C_{4} \lambda^{-\alpha} \frac{\log (3 / \lambda r)}{\log (1 / r)} \leq C_{5} \lambda^{-\alpha}
$$

where we used the assumption $r<R_{0}<1$ in the last inequality. We thus obtain the condition (SSJ-1).
We next prove (QRI-2). Using (4.22) and (4.24), we have for $0<r<R_{0}$,

$$
\begin{equation*}
N(x, r) \geq \frac{\kappa_{0}}{(2 r)^{d+\alpha}} \log \frac{3}{2 r} \cdot \frac{1}{3 \cdot 2^{d}}|B(x, r)| \geq C_{6} \frac{1}{r^{\alpha}} \log \frac{1}{r} . \tag{4.25}
\end{equation*}
$$

and

$$
\begin{aligned}
L(y, r / 2) & =S(y, r / 2)+\frac{2}{r}\left|\int_{\frac{r}{2} \leq|h| \leq 1} h n(x, h) d h\right|+\frac{4}{r^{2}} \int_{|h| \leq \frac{r}{2}}|h|^{2} n(x, h) d h \\
& \leq C_{3}\left(\frac{r}{2}\right)^{-\alpha} \log \frac{6}{r}+\frac{C_{7}}{r} \int_{\frac{r}{2}}^{1} \frac{1}{\psi^{\alpha}} \log \frac{3}{\psi} d \psi+\frac{C_{8}}{r^{2}} \int_{0}^{\frac{r}{2}} \frac{1}{\psi^{\alpha-1}} \log \frac{3}{\psi} d \psi
\end{aligned}
$$

$$
=: L_{1}+L_{2}+L_{3}
$$

It is clear that $L_{1} \leq C_{10} r^{-\alpha} \log (1 / r) . L_{3}$ satisfies the upper estimate as follows:

$$
L_{3} \leq \frac{C_{13}}{r^{2}} \int_{0}^{\frac{r}{2}} \frac{1}{\psi^{\alpha-1}} \log \frac{1}{\psi} d \psi \leq C_{14} r^{-\alpha} \log \frac{1}{r}
$$

Moreover, $L_{2}$ satisfies the upper estimate as follows:
Case1. $1<\alpha<2$.

$$
L_{2} \leq \frac{C_{11}}{r}\left(\frac{-\log (r / 2)+\log 3}{\alpha-1}\left(\frac{2}{r}\right)^{\alpha-1}+\frac{1}{(\alpha-1)^{2}}\right) \leq C_{15} r^{-\alpha} \log \frac{1}{r} .
$$

Case2. $\alpha=1$.

$$
L_{2} \leq \frac{C_{11}}{r}\left(\left(\log \frac{r}{2}\right)^{2}+\log 3 \cdot\left(-\log \frac{r}{2}\right)\right) \leq \frac{C_{16}}{r}\left(\log \frac{1}{r}\right)^{2} .
$$

Consequently, we conclude that

$$
L(y, r / 2) \leq \begin{cases}C_{11} r^{-\alpha} \log (1 / r) & (\alpha \neq 1)  \tag{4.26}\\ C_{12} r^{-1}(\log (1 / r))^{2} & (\alpha=1)\end{cases}
$$

Combining (4.25) and (4.26), we see that there exists $\kappa_{4}$ such that the condition (QRI-2) holds.

Remark 4.2. As we see in the proof above, when $1<\alpha<2$, (QRI-1) holds. Consequently, we have the Hölder continuity of harmonic functions.

Now we prove Theorem 4.3. We use the similar method as the case of Theorem 4.1. Note that Proposition 4.1 and 4.2 are valid under the assumption of Theorem 4.3. Moreover, the following lemma is obtained by easy modification of Proposition 4.3

Proposition 4.4. ([14] Lemma 2.5.) Suppose (QRI-2) holds. Then there exists $\kappa_{6}>0$ such that

$$
\mathbb{P}_{y}\left(T_{A}<\tau_{B\left(x_{0}, r\right)}\right) \geq \frac{\kappa_{6}}{|\log r|}
$$

for $r \in\left(0, R_{0}\right), y \in B\left(x_{0}, r / 2\right)$, and $A \subset B\left(x_{0}, r\right)$ such that $|A| \geq\left|B\left(x_{0}, r\right)\right| /\left(3 \cdot 2^{d}\right)$.
We now prove Theorem 4.3.
Proof of Theorem 4.3. It is sufficient to prove the assertion under the assumptions (SSJ-1) and (QRI-2). The idea of the proof is the same as Theorem 4.1. Suppose $\|u\|_{\infty}=K, z_{1} \in B\left(z_{0}, R / 2\right)$ and define

$$
s_{n}:=\frac{\theta_{1}}{(\log (n+1))^{\beta_{1}}}, \quad r_{n}:=\frac{\theta_{2}}{n^{\beta_{2} n}},
$$

where $\theta_{1}$ is chosen large enough to satisfy $\theta_{1}>2 K$, and $\theta_{2}$ is chosen small enough to satisfy $\theta_{2}<R / 4$. In addition, suppose $0<\beta_{1} \leq 1$ and $\beta_{2} \geq 1$ which are chosen
appropriately later. Note that $\lim _{n \rightarrow \infty} r_{n}=0$ and $\lim _{n \rightarrow \infty} s_{n}=0$. Set $B_{n}, M_{n}$ and $m_{n}$ as follows.

$$
B_{n}:=B\left(z_{1}, r_{n}\right), \quad M_{n}:=\sup _{z \in B_{n}} u(z), \quad m_{n}:=\inf _{z \in B_{n}} u(z) .
$$

As in the proof of Theorem 4.1, we will prove by induction that

$$
M_{n}-m_{n} \leq s_{n} \quad(\forall n \in \mathbb{N})
$$

We assume that $M_{k}-m_{k} \leq s_{k}$ for $1 \leq k \leq n$ and consider the case of $k=n+1$. Choose $y, z \in B_{n+1}$ such that

$$
u(z) \geq M_{n+1}-\varepsilon, \quad u(y) \leq m_{n+1}+\varepsilon
$$

where $\varepsilon$ is an arbitrary positive number. Define

$$
A_{n}=\left\{z \in B_{n}: u(z) \leq\left(M_{n}+m_{n}\right) / 2\right\}
$$

and assume that $\left|A_{n}\right| /\left|B_{n}\right| \geq 1 / 2$. Let $D_{n}$ be a compact set that satisfies

$$
D_{n} \subset A_{n}, \quad\left|D_{n}\right| /\left|B_{n}\right| \geq 1 / 3 .
$$

Let

$$
p_{n}:=\mathbb{P}_{z}\left(T_{D_{n}}<\tau_{n}\right), \quad F_{i}:=\mathbb{P}_{z}\left(X_{\tau_{n}} \notin B_{n-i}\right),
$$

where $\tau_{n}:=\tau_{B_{n}}$. The harmonic property of $u$ and the optional sampling theorem imply that

$$
\begin{align*}
u(z)-u(y) & =\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y)\right] \\
& =\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): T_{D_{n}}<\tau_{n}, X_{\tau_{n}} \in B_{n-1}\right] \\
& +\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): T_{D_{n}} \geq \tau_{n}, X_{\tau_{n}} \in B_{n-1}\right] \\
& +\sum_{i=1}^{n-2} \mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): X_{\tau_{n}} \in B_{n-i-1} \backslash B_{n-i}\right] \\
& +\mathbb{E}_{z}\left[u\left(X_{\tau_{n} \wedge T_{D_{n}}}\right)-u(y): X_{\tau_{n}} \notin B_{1}\right] \\
& =: I_{1}+I_{2}+I_{3}+I_{4} . \tag{4.27}
\end{align*}
$$

Here we define $u\left(X_{\infty}\right):=\lim _{t \rightarrow \infty} u\left(X_{t}\right)$ and assume $X_{\tau_{n}} \in B_{n-1}$ when $\tau_{n}=\infty$. Now we estimate each term of (4.27). First, we consider $I_{1}$ and $I_{2}$.

$$
\begin{align*}
& I_{1} \leq \frac{M_{n}-m_{n}}{2} \mathbb{P}_{z}\left(T_{D_{n}}<\tau_{n}\right) \leq \frac{s_{n}}{2} p_{n}  \tag{4.28}\\
& I_{2} \leq\left(M_{n-1}-m_{n-1}\right) \mathbb{P}_{z}\left(\tau_{n} \leq T_{D_{n}}\right) \leq s_{n-1}\left(1-p_{n}\right) \tag{4.29}
\end{align*}
$$

In order to estimate $I_{3}$ and $I_{4}$, we use $F_{i}$.

$$
I_{3} \leq \sum_{i=1}^{n-2}\left(M_{n-i-1}-m_{n-i-1}\right)\left(F_{i}-F_{i+1}\right) \leq \sum_{i=1}^{n-2} s_{n-i-1}\left(F_{i}-F_{i+1}\right)
$$

$$
\begin{align*}
& =s_{n-2} F_{1}-s_{1} F_{n-1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i} \\
& \leq s_{n-2} F_{1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i} .  \tag{4.30}\\
I_{4} & \leq 2 K \cdot F_{n-1} . \tag{4.31}
\end{align*}
$$

Combining the estimates from (4.28) to (4.31), we can obtain

$$
\begin{align*}
u(z)-u(y) & \leq \frac{1}{2} s_{n} p_{n}+s_{n-1}\left(1-p_{n}\right)+s_{n-2} F_{1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i}+2 K \cdot F_{n-1} \\
& \leq s_{n-1}-\frac{1}{2} s_{n} p_{n}+s_{n-2} F_{1}+\sum_{i=2}^{n-2}\left(s_{n-i-1}-s_{n-i}\right) F_{i}+2 K \cdot F_{n-1} \\
& =s_{n+1}\left(\frac{s_{n-1}}{s_{n+1}}-\frac{s_{n} p_{n}}{2 s_{n+1}}+\frac{s_{n-2}}{s_{n+1}} F_{1}+\sum_{i=2}^{n-2} \frac{s_{n-i-1}-s_{n-i}}{s_{n+1}} F_{i}+\frac{2 K}{s_{n+1}} \cdot F_{n-1}\right) \\
& =: s_{n+1}\left(J_{1}+J_{2}+J_{3}+J_{4}+J_{5}\right) . \tag{4.32}
\end{align*}
$$

Our goal is to prove that

$$
\begin{equation*}
J_{1}+J_{2}+J_{3}+J_{4}+J_{5} \leq 1 \tag{4.33}
\end{equation*}
$$

Now, we estimate each $J_{i}$. In the rest of the proof, we assume that $n$ is large enough. $J_{1}$ and $J_{2}$ are easily estimated.

$$
\begin{align*}
J_{1} & =\frac{s_{n-1}}{s_{n+1}}=\left(\frac{\log (n+2)}{\log n}\right)^{\beta_{1}}=\left(1+\frac{\log \left(1+\frac{2}{n}\right)}{\log n}\right)^{\beta_{1}} \\
& \leq\left(1+\frac{2}{n \log n}\right)^{\beta_{1}} \leq 1+\frac{2 \beta_{1}}{n \log n} \tag{4.34}
\end{align*}
$$

where we used the assumption $0<\beta_{1} \leq 1$ in the last inequality. Note that Lemma 4.4 implies $p_{n} \geq \kappa_{6} /\left|\log r_{n}\right|$. Then, we obtain

$$
\begin{align*}
J_{2}=-\frac{s_{n} p_{n}}{2 s_{n+1}} & \leq-\frac{1}{2} p_{n} \leq-\frac{\kappa_{6}}{2\left|\log r_{n}\right|} \\
& =-\frac{\kappa_{6}}{2\left(\left|\log \theta_{2}\right|+\beta_{2} n \log n\right)} \leq-\frac{\kappa_{7}}{\beta_{2} n \log n} \tag{4.35}
\end{align*}
$$

where $\kappa_{7}$ is a constant that does not depend on $\beta_{1}, \beta_{2}$ or $n$. In order to estimate $J_{3}, J_{4}$ and $J_{5}$, we begin with the estimate of $F_{i}:=\mathbb{P}^{z}\left(X_{\tau_{n}} \notin B_{n-i}\right)$. If $X_{\tau_{n}} \notin B_{n-i}$, the process cannot have a jump larger than $2 r_{n}$ before the time $\tau_{n}$ and $\left|X_{\tau_{n}}-X_{\tau_{n}-}\right| \geq r_{n-i}-r_{n}$. Note that $r_{n} / r_{n-1} \leq 1 / 4$ and therefore $2 r_{n}<r_{n-i}-r_{n}$. Then, we can apply Lemma 4.2 and obtain

$$
F_{i} \leq \kappa_{1}\left(\frac{2 r_{n}}{r_{n-i}-r_{n}}\right)^{\sigma} \quad(1 \leq i \leq n)
$$

Thus, noting that $0<\beta_{1} \leq 1 \leq \beta_{2}$, we have

$$
\begin{align*}
J_{5} & =\frac{2 K}{s_{n+1}} \cdot F_{n-1} \leq(\log (n+2))^{\beta_{1}} \cdot \kappa_{1}\left(\frac{2 r_{n}}{r_{1}-r_{n}}\right)^{\sigma} \\
& =(\log (n+2))^{\beta_{1}} \cdot \kappa_{1}\left(\frac{2}{n^{\beta_{2} n}-1}\right)^{\sigma} \leq \log (n+2) \cdot \kappa_{1}\left(\frac{4}{n^{n}}\right)^{\sigma} \\
& \leq C_{1} e^{-n \sigma} \log n \leq C_{2} \exp \left(-\frac{n \sigma}{2}\right) . \tag{4.36}
\end{align*}
$$

The estimates of $J_{3}$ and $J_{4}$ are a little complicated.

$$
\begin{align*}
J_{3} & =\frac{s_{n-2}}{s_{n+1}} \cdot F_{1} \leq\left(\frac{\log (n+2)}{\log (n-1)}\right)^{\beta_{1}} \cdot \kappa_{1}\left(\frac{2 r_{n}}{r_{n-1}-r_{n}}\right)^{\sigma} \\
& \leq\left(\frac{\log (n+2)}{\log (n-1)}\right)^{\beta_{1}} \cdot \kappa_{1}\left(\frac{2}{n^{\beta_{2}}-1}\right)^{\sigma} \leq\left(\frac{\log (n+2)}{\log (n-1)}\right)^{\beta_{1}} \cdot \kappa_{1}\left(\frac{4}{n^{\beta_{2}}}\right)^{\sigma} \\
& \leq C_{3}\left(1+\frac{3}{(n-1) \log (n-1)}\right)^{\beta_{1}} \cdot \frac{1}{n^{\beta_{2} \sigma}} \leq \frac{C_{3}}{n^{\beta_{2} \sigma}}+\frac{C_{4}}{n^{\beta_{2} \sigma+1} \log n} . \tag{4.37}
\end{align*}
$$

In order to estimate $J_{4}$, we begin with the following decomposition.

$$
\begin{align*}
J_{4} & =\sum_{i=2}^{n-2} \frac{s_{n-i-1}-s_{n-i}}{s_{n+1}} \cdot F_{i} \\
& =\sum_{i=2}^{n_{0}-1} \frac{s_{n-i-1}-s_{n-i}}{s_{n+1}} \cdot F_{i}+\sum_{i=n_{0}}^{n-2} \frac{s_{n-i-1}-s_{n-i}}{s_{n+1}} \cdot F_{i}=: K_{1}+K_{2}, \tag{4.38}
\end{align*}
$$

where $n_{0}$ is an integer that satisfies $n_{0}-1<\sqrt{n} \leq n_{0}$. Note that $r_{n+1} / r_{n} \leq 1 / 4$ and therefore

$$
\begin{equation*}
F_{i} \leq \kappa_{1}\left(\frac{2 r_{n}}{r_{n-i}-r_{n}}\right)^{\sigma} \leq \kappa_{1}\left(\frac{2}{4^{i}-1}\right)^{\sigma} \leq \frac{C_{5}}{4^{(i-1) \sigma}} \tag{4.39}
\end{equation*}
$$

Noting that $(\log x)^{-\beta_{1}}$ is a convex monotone decreasing function on $x>1$, we can estimate $K_{1}$ as follows;

$$
\begin{align*}
K_{1} & \leq(\log (n+2))^{\beta_{1}}\left(\frac{1}{(\log (n-\sqrt{n}))^{\beta_{1}}}-\frac{1}{(\log (n-\sqrt{n}+1))^{\beta_{1}}}\right) \cdot \sum_{i=2}^{n_{0}-1} \frac{C_{5}}{4^{(i-1) \sigma}} \\
& \leq C_{6}(\log (n+2))^{\beta_{1}} \cdot \beta_{1} \frac{1}{(n-\sqrt{n})(\log (n-\sqrt{n}))^{\beta_{1}+1}} \\
& \leq C_{6}(2 \log n)^{\beta_{1}} \cdot \beta_{1} \frac{2^{\beta_{1}+2}}{n(\log n)^{\beta_{1}+1}} \leq \frac{C_{7} \beta_{1}}{n \log n} . \tag{4.40}
\end{align*}
$$

The estimate of $K_{2}$ is much easier;

$$
\begin{equation*}
K_{2} \leq\left(\frac{\log (n+2)}{\log 2}\right)^{\beta_{1}} \sum_{i=n_{0}}^{\infty} \frac{C_{5}}{4^{(i-1) \sigma}} \leq C_{8} \log (n+2) \cdot \frac{1}{4 \sqrt{n} \sigma} \leq \frac{C_{9}}{2 \sqrt{n} \sigma} . \tag{4.41}
\end{equation*}
$$

From (4.40) and (4.41), we obtain

$$
\begin{equation*}
J_{4} \leq \frac{C_{7} \beta_{1}}{n \log n}+\frac{C_{9}}{2 \sqrt{n} \sigma} . \tag{4.42}
\end{equation*}
$$

Combining all the estimates of $J_{i}$ 's, we conclude that

$$
\begin{align*}
J_{1} & +J_{2}+J_{3}+J_{4}+J_{5} \\
& \leq 1+\frac{2 \beta_{1}}{n \log n}-\frac{\kappa_{7}}{\beta_{2} n \log n}+\frac{C_{3}}{n^{\beta_{2} \sigma}}+\frac{C_{4}}{n^{\beta_{2} \sigma+1} \log n}+\frac{C_{7} \beta_{1}}{n \log n}+\frac{C_{9}}{2 \sqrt{n} \sigma}+C_{2} \exp \left(-\frac{n \sigma}{2}\right) \\
& \leq 1+\left(C_{10} \beta_{1}-\frac{\kappa_{7}}{\beta_{2}}\right) \frac{1}{n \log n}+\frac{C_{3}}{n^{\beta_{2} \sigma}}+\frac{C_{11}}{n^{\beta_{2} \sigma+1} \log n} . \tag{4.43}
\end{align*}
$$

Now, choose $\beta_{2} \geq 1$ so that $\beta_{2} \sigma>1$, and then take $\beta_{1}$ small enough so that $C_{10} \beta_{1}-$ $\kappa_{7} / \beta_{2}<-\kappa_{7} / 2 \beta_{2}$. Then, we conclude that there exist $\kappa_{8}$ and $n_{1} \in \mathbb{N}$ such that the right hand side of (4.43) is bounded from above by $1-\kappa_{8} /(n \log n)$ for $n \geq n_{1}$. If we choose $\theta_{1}$ sufficiently large so that

$$
M_{n}-m_{n} \leq 2 K \leq s_{n} \quad\left(n \leq n_{1}\right),
$$

the induction hypothesis also holds for all $n \leq n_{1}$. Now we conclude $|u(z)-u(y)| \leq s_{n}$ for $z, y \in B\left(z_{0}, R / 2\right)$ such that $|z-y| \leq r_{n}$. Noting that $s_{1} \geq 2 K$, we can obtain the desired continuity.

Remark 4.3. In the proof, we should choose $s_{n}$ and $r_{n}$ very carefully. In Theorem 1.2 of [14], they set

$$
s_{n}=\frac{\theta_{1}}{n^{\zeta}}, \quad r_{n}=\frac{\theta_{2}}{4^{n}}
$$

where $\zeta>0$. However, this does not work well. The trouble is that $J_{2}$ may go to 0 as $n \rightarrow \infty$, while $J_{3}$ does not depend on $n$. To avoid this trouble, it is essential to choose $r_{n}$ so that $F_{1}$ also depends on $n$.

Recently Kassmann and Mimica prove Hölder continuity of harmonic functions for more general non-local operators in [18].

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