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# Dimitar I. Hadjiev <br> The first passage problem for generalized OrnsteinUhlenbeck processes with non-positive jumps 

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PROCESSES WITH NON-POSITIVE JUMPS *

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1. Introduction. Let $(\Omega, F, P)$ be a probability space. We consider a cadlag stationary random process $S_{t}, t \geq 0$, with independent increments and non-positive jumps $\Delta S_{t}=S_{t}-S_{t-}=S_{t}-\lim _{S \uparrow t} S_{S} \leqq 0$, that is defined on this space and satisfies $S_{0}=0$. It is well known ([3]) that the characteristic function of $S_{t}$ has the form

$$
\begin{equation*}
E \exp \left(i u S_{t}\right)=\exp t\left\{i b u-c u^{2}+\underset{(-\infty, 0)}{\int} F(d x)\left(e^{i u x}-1-i u x . l_{\{x>-1\}}\right)\right\} \tag{1.1}
\end{equation*}
$$

where $-\infty<b<\infty, C \geq 0$, and the Lévy measure $F($.$) satisfies$

$$
\begin{equation*}
\int_{(-\infty, 0)}^{\int} F(d x) \quad l \Lambda x^{2}<\infty . \tag{1.2}
\end{equation*}
$$

Following Skorokhod ([8]) one can use the analytical continuation of (l.1) to the half-plane $\operatorname{Re}(i u)>0$ and obtain the Laplace transform of $S_{t}$ by substituting u instead of iu. Thus, we have
(1.3) $\quad E \exp \left(u S_{t}\right)=\exp t \psi(u), u \geqq 0$,
where
(1.4) $\quad \psi(u)=b u+c u^{2}+\int_{(-\infty, 0)} F(d x)\left(e^{u x}-1-u x . l_{\{x \geqslant-1\}}\right)$.

For arbitrary $\lambda>0$ and $-\infty<x<\infty$ we define the random process $X_{t}, t \geqslant 0$, by the formula

$$
\begin{equation*}
x_{t}=e^{-\lambda t}\left(x+\underset{(0, t]}{s} e^{\lambda v} d S_{v}\right) \tag{1.5}
\end{equation*}
$$

the stochastic integral w.r.t. the semi-martingale $S$ being understood in the usual sense.

Definition. The random process X will be called the starting at x generalized Ornstein-Uhlenbeck process with parameter $\lambda>0$.

[^0]Certainly, the process $X$ is characterized by the triplet (b,c,F(.)) as well. With $b=0, c=1 / 2$ and $F()=$.0 our definition yields the standard Wiener process $S$ and the usual Ornstein-Uhlenbeck process X .

Given a real number $\mu>x$, let us introduce the first passage time

$$
\begin{equation*}
T_{\mu}(x)=\inf \left\{t \geqslant 0: X_{t} \geqq \mu\right\} \tag{1.6}
\end{equation*}
$$

As far as $\Delta X_{t}=\Delta S_{t} \leq 0$, if $T_{\mu}(x)<\infty$ one gets immediately the equality $X_{T_{\mu}(x)}=\mu$.

The purpose of this paper is to determine the distribution of $T_{\mu}(x), \mu>x$, by means of Laplace transform

$$
\begin{equation*}
\gamma_{\mu}(\theta, x)=E \exp \left(-\theta T_{\mu}(x)\right), \quad \theta>0 \tag{1.7}
\end{equation*}
$$

It should be noted that generally speaking, we have no equation for the transition density of X and the usual Darling-Siegert approach to the first passage problem of diffusion processes ([2]) is not applicable in our case. Our approach is based on martingale techniques and depends essentially on the existence of suitable martingales on the process X (see Theorem 1 below). Besides the new generality of the explicit representation for $\gamma_{\mu}(\theta, x)$ (Section 4), this approach gives us in particular the possibility to obtain ones again and in a natural way the interesting result of Novikov ([6] ) concerning the first passage times of a stable process $S$ through one-sided non-linear boundaries. The basic tool in this special case is the suitable time-change (Section 6) that transfers the linear problems for $X_{t}, t \geqslant 0$, into some non-linear problems for $S_{t}, t \geqq 0$, and conversely. We make use of the reconversion in order to give an example of optimal stopping problem that admits a solution in terms of $T_{\mu}(x)$.
2. The process X . For the next we need to calculate the conditional Laplace transforms of the process X that was defined in (1.5). Let us introduce the $\sigma$-algebras $F_{t}^{X}=\sigma\left(X_{s}, 0 \leq s \leq t\right) ; t \geqq 0$, and the functions $L(u ; t, s)=E\left\{\exp \left(u X_{t}\right) \mid F_{s} X_{s}, s<t, u>0\right.$.

Since the stochastic integral in (1.5) might be looked at as an integral taken in the sense of convergence in probability ([4]), a simple argument leads to the following result.

Proposition 1. For any $0 \leqq s<t$ and $u \geqq 0$ one has

$$
\begin{equation*}
L(u ; t, s)=\exp \left\{e^{-\lambda(t-s)} X_{s} \cdot u+\int_{s}^{t} \psi\left(u \cdot e^{-\lambda(t-v)}\right) d v\right\} \tag{2.1}
\end{equation*}
$$

Proof. With an arbitrary subdivision $s=t_{0}<t_{1}<\ldots<t_{n}=t, \quad \varepsilon=\max _{i \leqslant n}\left|t_{i}-t_{i-1}\right|$ and $Y_{t}=\int_{(0, t)} e^{\lambda v} d S_{v}$, we get

$$
\begin{aligned}
& E\left\{\exp \left(u Y_{t}\right) \mid F_{s} X_{\}}=\exp \left(u Y_{S}\right) \cdot E\left\{\exp \left(u \underset{(s, t]}{\delta} e^{\lambda v} d S_{v}\right) \mid F_{s} X_{i}\right.\right. \\
& =\exp \left(u Y_{s}\right) \quad \lim _{\varepsilon \downarrow 0} \prod_{i=1}^{n} E \exp \left(u \cdot e^{\lambda t_{i-1}} .\left(S_{t_{i}}-S_{t_{i-1}}\right)\right) \\
& =\exp \left(u_{s}\right) \quad \lim _{\varepsilon \neq 0} \prod_{i=1}^{n} \exp \left\{\psi\left(u . e^{\lambda t_{i}-1}\right)\left(t_{i}-t_{i-1}\right)\right\} \\
& =\exp \left(u Y_{s}+\int_{s}^{t} \psi\left(u \cdot e^{\lambda v}\right) d v\right)
\end{aligned}
$$

as a consequence of (1.3) and the independent increments property of $S$.
Now starting with (1.5) we have

$$
\begin{aligned}
L(u ; t, s) & =\exp \left(e^{-\lambda t} \cdot x u\right) E\left\{\exp \left(u \cdot e^{-\lambda t} \cdot Y_{t}\right) \mid F_{s} X_{\}}\right. \\
& =\exp \left\{e^{-\lambda t} \cdot x u+e^{-\lambda t} \cdot Y_{s} u+\int_{s}^{t} \psi\left(u \cdot e^{-\lambda(t-v)}\right) d v\right\}
\end{aligned}
$$

and the latter obviously implies (2.1).
Corollary 1. The Laplace transform of $X_{t}$ has the form

$$
E \exp \left(u x_{t}\right)=\exp \left\{e^{-\lambda t} \cdot x u+\int_{0}^{t} \psi\left(u \cdot e^{-\lambda(t-v)}\right) d v\right\}, u \geqq 0
$$

Corollary 2. The process X is a cadlag Markov process. (Certainly, X has also the strong Markov property.)
3. The martingale $M$. We are going to introduce a martingale $M_{t}(\theta), t \geqq 0$, depending on the process $X$ trajectories. To this end, one observes that because of (1.2) the quantity $F[-1,-z]$ is finite for every $z, 0<z \leqq 1$. Thus, the measure

$$
\mathrm{G}(\mathrm{dz})=\mathrm{F}[-1,-\mathrm{z}] \mathrm{dz}
$$

on $(0,1]$ is well defined. We need the following assumption.
Hypothesis G. Either $c>0$ or the measure G(.) satisfies the condition

$$
\begin{equation*}
\lim _{z \rightarrow 0+} z^{k} \cdot G(z, 1]=C>0 \tag{3.1}
\end{equation*}
$$

for some constant $k$, $0<k<1$.
Next, one defines successively

$$
\begin{equation*}
g(y)=-\frac{1}{\lambda} \int_{1}^{y} \frac{\psi(u)}{u} d u, \quad y>0 \tag{3.2}
\end{equation*}
$$

and

$$
\begin{equation*}
M_{t}(\theta)=e^{-\theta t} \cdot \int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} \cdot \exp \left\{x_{t} \cdot y+g(y)\right\} d y, \quad t \geqq 0 \tag{3.3}
\end{equation*}
$$

The next statement is crucial because it permits an essential use of the martingale theory later on.

Theorem 1. Under the hypothesis $G$ for any positive $\theta$ the random process $M_{t}(\theta)$, $t \geqq 0$, is a martingale w.r.t. $\quad F_{t^{\prime}}^{X} t \geqq 0$.
Proof. First, we observe that our hypothesis $G$ implies the convergence of the integral in (3.3). In fact, we have

$$
g(y)=-\frac{b}{\lambda}(y-1)-\frac{c}{2 \lambda}\left(y^{2}-1\right)-\frac{1}{\lambda} g_{1}(y)-\frac{1}{\lambda} g_{2}(y),
$$

where

$$
g_{1}(y)=\int_{1}^{y} \frac{\psi_{1}(u)}{u} d u \quad, \quad g_{2}(y)=\int_{1}^{y} \frac{\psi_{2}(u)}{u} d u
$$

and

$$
\psi_{1}(u)=\int_{(-\infty,-1)}^{\int} F(d x)\left(e^{u x}-1\right), \psi_{2}(u)=\int_{[-1,0)} F(d x)\left(e^{u x}-1-u x\right), u \geq 0
$$

The convergence of the integral at $y=0$ is obvious, because $\infty>\lim _{y \neq 0} g(y) \geqq-\infty$. Now let us denote $d_{1}=\underset{(-\infty,-1)}{\int} F(d x) \geqq 0, d_{2}=\underset{[-1,0)}{\int} F(d x) x^{2} \geqq 0$. In consequence of (1.2) one gets $0 \leqq d_{1}+d_{2}<\infty$. Our function $\psi_{1}$ satisfies $0 \geqq \psi_{1}(u)+-d_{1}$ and $0 \geqq \frac{\Psi_{1}(\mathrm{u})}{\mathrm{u}} \uparrow 0$ as $\mathrm{u} \uparrow \infty$. This means that $\left.\left|g_{1}(\mathrm{y})\right| \leqq \int_{1} \frac{\mid \psi_{1}(\mathrm{u})}{\mathrm{u}} \right\rvert\, \mathrm{du} \leqq \mathrm{d}_{1} \ln \mathrm{y}$. On the other hand $0<e^{u x}-1-u x \leqq \frac{u x}{2}, u>0,-1 \leqq x<0$, and in this way one obtains the inequalities $0 \leqq \frac{\psi_{2}(u)}{u} \leqq \frac{u}{2} . d_{2}<\infty$ and $0 \leqq g_{2}(y) \leqq \frac{d}{4} 2\left(y^{2}-1\right)$.

If $c>0$, the corresponding term $-\frac{c}{2 \lambda}\left(y^{2}-1\right)$ in $g(y)$ ensures the convergence. If $c=0$, by the equality $\frac{\psi_{2}(u)}{u}=\int_{0}^{1}\left(1-e^{-u z}\right) G(d z)$, where obviously $0 \leqq \int_{0}^{1} \mathrm{zG}(\mathrm{dz})=\frac{\mathrm{d}_{2}}{2}<\infty$, the hypothesis (3.1) and the corollary of Theorem 4.15 in [1] one gets $\frac{l_{i m}}{u \rightarrow \infty} u^{-k} \cdot \frac{\psi_{2}(u)}{u} \geqq C . \Gamma(1-k)>0$. Consequently, $\frac{\psi_{2}(u)}{u} \geq C_{2} \cdot u^{k}$ for any $\mathrm{C}_{2}$ belonging to the interval $\left(0, \mathrm{C} . \Gamma(1-\kappa)\right.$ ) and $\mathrm{u} \geqq \mathrm{u}_{2}\left(\mathrm{C}_{2}\right)>0$ (sufficiently large). This implies $g_{2}(y) \geqq C_{2} \cdot y^{1+\kappa}+C_{1}, y>u_{2}\left(C_{2}\right)$, and the convergence of our integral too.

Secondly, applying Fubini's lerma and (2.1) for $0 \leqq s \leqq t$ (and with $z=y e^{-\lambda(t-s)}$ we get

$$
\begin{aligned}
& E\left\{M_{t}(\theta) \left\lvert\, F_{s}^{X_{S}}=e^{-\theta t} \cdot \int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} E\left\{\exp \left(X_{t} \cdot y+g(y)\right) \mid F_{s}^{X}\right\} d y\right.\right. \\
= & e^{-\theta s} \cdot \int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} \exp \left\{g(y)-\theta(t-s)+e^{-\lambda(t-s)} y \cdot X_{s}+\int_{s}^{t} \psi\left(y e^{-\lambda(t-v)}\right) d v\right\} d y \\
= & e^{-\theta s} \cdot \int_{0}^{\infty} z^{\frac{\theta}{\lambda}-1} \exp \left\{z X_{s}+g\left(z e^{\lambda(t-s)}\right)+\int_{0}^{t-s} \psi\left(z e^{\lambda v}\right) d v\right\} d z .
\end{aligned}
$$

But the function $f(u, z)=g\left(z e^{\lambda u}\right)+\int_{0}^{u} \psi\left(z e^{\lambda v}\right) d v, u \geqq 0$, satisfies the condition

$$
\frac{\partial f(u, z)}{\partial u}=g^{\prime}\left(z e^{\lambda u}\right) \cdot z \lambda e^{\lambda u}+\psi\left(z e^{\lambda u}\right)=g^{\prime}(y) \cdot \lambda y+\psi(y) \equiv 0
$$

with $y=z e^{\lambda u}$, in view of (3.2). Therefore,

$$
f(u, z)=\text { const }=f(0, z)=g(z)
$$

and we get $E\left\{M_{t}(\theta) \mid F_{s} X_{s}=X_{s}\right.$, that completes the proof.
Remark 1. We emphasize the fact that Theorem 1 is valid for every process $X$ with $S$ containing a Gaussian component ( $c>0$ ). If the process $S$ has no Gaussian component $(c=0)$, the condition (3.1) is nevertheless fulfilled for a class of measures $F($. that includes the stable processes $S$ with parameter $\alpha$ satisfying $1<\alpha<2$. Because of its importance, we consider this special case in Section 5.
4. The Laplace transform of $T_{\mu}(x)$. Now we are in a position to derive an explicite expression for the Laplace transform $\gamma_{\mu}(\theta, x)$. Due to the particular structure of
the martingale $\mathrm{M}(\theta)$ we have the following result.
Theorem 2. Under the hypothesis G the next equality holds:

$$
\begin{equation*}
\gamma_{\mu}(\theta, x)=\frac{\int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} \exp (x y+g(y)) d y}{\int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} \exp (\mu y+g(y)) d y}, \theta>0 \tag{4.1}
\end{equation*}
$$

Proof. We put $T_{\mu}(x) \Lambda t$ instead of $t$ in (3.3) and we make use of the well known martingale property that

$$
E M_{T_{\mu}}(x) \Lambda t(\theta)=E M_{0}(\theta)=\int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} \exp (x y+g(y)) d y
$$

Next, one observes that
and, moreover, when $T_{\mu}(x)=\infty$ then

$$
0 \leqq M_{T_{\mu}}(x) \Lambda t(\theta)=M_{t}(\theta) \leqq e^{-\theta t} \int_{0}^{\infty} y^{\frac{\theta}{\lambda}-1} \exp (\mu y+g(y)) d y
$$

as well. Therefore,


The right-hand sides of our equalities give directly (4.1).
Remark 2. For the validity of Theorem 2 we need not (and we did not use) any fact about the finiteness of $T_{\mu}(x)$. It is well known that $T_{\mu}(x)<\infty \quad P-a . s$. if and only if $\lim _{\theta \downarrow 0} \gamma_{\mu}(\theta, x)=1$. The latter equality is easily verified when there exists $\lim _{y \downarrow 0} g(y)>-\infty$ or when $\underset{(-\infty,-1)}{\int} F(d x)|x|<\infty$.
5. The case of stable process $S$ with parameter $1<\alpha \leq 2$. Now we turn to the particular case when the following hypothesis is satisfied.

Hypothesis $H_{\alpha}$. Either $F()=$.0 and $c>0$ (we characterize this by posing $\alpha=2$ ), or $c=0$ and $F(d x)=\frac{\sigma \cdot d x}{|x|^{\alpha+1}} 1_{\{x<0\}}$ for some $\sigma>0$ and $1<\alpha<2$.

Using standard arguments (see [8], §25, Theorem 4) one obtains the equivalent form of $H_{\alpha}$ in the terms of our function $\psi: H_{\alpha}, l<\alpha \leqq 2$, means that

$$
\begin{equation*}
\psi(\mathrm{u})=\overline{\psi(u)}=\overline{\mathrm{b}} \mathrm{u}+\bar{\sigma} \mathrm{u}^{\alpha} \tag{5.1}
\end{equation*}
$$

with some $\overline{\mathrm{b}},-\infty<\overline{\mathrm{b}}<\infty$, and $\bar{\sigma}>0$. In this situation by (3.2) we get

$$
\begin{equation*}
g(y)=\bar{g}(y)=-\frac{\bar{b}}{\lambda}(y-1)-\frac{\bar{\sigma}}{\alpha \lambda}\left(y^{\alpha}-1\right), \tag{5.2}
\end{equation*}
$$

and the martingale $M(\theta)$ is well defined via (3.3).
Following Novikov we introduce the function

$$
\mathrm{H}(\nu, \alpha, x)=\frac{1}{\Gamma(-\alpha \nu)} \int_{0}^{\infty} y^{-\alpha \nu-1} \exp \left(x y-\frac{1}{\alpha} y^{\alpha}\right) d y,
$$

which turns to be analytic in the half-plane Rev < 1 . All the essential properties of $\mathrm{H}(\nu, \alpha, \mathrm{x})$ are collected in the supplement of [6].

Next we obtain a special case of Theorem 2.
Proposition 2. Under the hypothesis $H_{\alpha}, l<\alpha \leqq 2$, the following equality holds for $\theta>0$ :

$$
\begin{equation*}
\gamma_{\mu}(\theta, x)=\frac{H\left(-\frac{\theta}{\alpha \lambda}, \alpha,\left(\frac{\lambda}{\bar{\sigma}}\right)^{\frac{1}{\alpha}}\left(x-\frac{\bar{b}}{\lambda}\right)\right)}{H\left(-\frac{\theta}{\alpha \lambda}, \alpha,\left(\frac{\lambda}{\bar{\sigma}}\right)^{\frac{1}{\alpha}}\left(\mu-\frac{\bar{b}}{\lambda}\right)\right)} . \tag{5.3}
\end{equation*}
$$

Moreover, this formula defines also an analytical continuation of the Laplace transform $\gamma_{\mu}(\theta, x)$ to the half-plane $\operatorname{Re} \theta>-\alpha \lambda \cdot \nu_{\alpha}(\bar{\mu})$, where $\bar{\mu}=\left(\frac{\lambda}{\bar{\sigma}}\right)^{\frac{1}{\alpha}}\left(\mu-\frac{\bar{b}}{\lambda}\right)$ and $\nu_{\alpha}(z)$ is the smallest positive zero of $H(\nu, \alpha, z)$ with $(\alpha, z)$ fixed. Proof. Applying the change of variables $y=\left(\frac{\lambda}{\bar{\sigma}}\right)^{\frac{1}{\alpha}}{ }^{2}$ we see the formula (5.3) is another form of (4.1) for $\theta>0$. As far as the right-hand side of (5.3) is analytic in $\theta$ in the half-plane $\operatorname{Re} \theta>-\alpha \lambda \cdot \nu_{\alpha}(\bar{\mu})$ (see [6] ), the left-hand side can be analytically continued in $\theta$ to this half-plane.

Corollary 3. Since $\lim _{v \rightarrow 0} H(v, \alpha, x)=1,-\infty<x<\infty$, under the hypothesis $H_{\alpha}$ we get $\lim _{\theta \downarrow 0} \gamma_{\mu}(\theta, x)=1$ and, consequently, $T_{\mu}(x)<\infty \quad$ P-a.s.
6. The time change - two applications. Throughout this section we suppose the hypothesis $H_{\alpha}$ holds with some $\alpha, 1<\alpha \leqq 2$, and $\overline{\mathrm{b}}=0$ (see (5.1)). As a consequence we have

$$
\psi(u)=\widetilde{\psi}(u)=\bar{\sigma} \cdot u^{\alpha}, 1<\alpha \leqq 2,
$$

and the process X is stationary too (see (2.1)).

Let us introduce the real (increasing and continuous) function

$$
\delta(t)=(\alpha \lambda)^{-1}\left(e^{\alpha \lambda t}-1\right), t \geqq 0,
$$

which determines an one-to-one mapping of $[0, \infty)$ onto $[0, \infty)$, and the convers function

$$
\rho(t)=(\alpha \lambda)^{-1} \ln (1+\alpha \lambda t), t \geq 0
$$

Lemma 1. The distributions of $S_{t}, t \geqq 0$, and of $\tilde{S}_{t}=\int_{0}^{\rho(t)} e^{\lambda v_{d S}} d v^{\prime} t \geqq 0$, coinside.

## Proof. As in Proposition 1 one calculates

$E \exp \left(u \tilde{S}_{t}\right)=E \exp \left(u Y_{\rho(t)}\right)=\exp \left\{\bar{\sigma} u^{\alpha} . \delta(\rho(t))\right\}=\exp \left(\bar{\sigma} u^{\alpha} t\right), u>0$.
But under the hypothesis stated ( $\mathrm{H}_{\alpha}$ and $\overline{\mathrm{b}}=0$ ) the latter term is just $\mathrm{E} \exp \left(u S_{t}\right)$. The lemma is proved.

Now for any constants $a, b$ and $c$ such that $b \geqq 0$ and $a b^{\underline{1}}+c>0$, define the stopping time $\tau(a, b, c)$ w.r.t. $F_{t^{\prime}}^{S} t \geqq 0$, by the formula

$$
\begin{equation*}
\tau(a, b, c)=\inf \left\{t>0: S_{t} \geq a(t+b)^{\frac{1}{\alpha}}+c\right\} \tag{6.1}
\end{equation*}
$$

and pose

$$
\begin{equation*}
\tau_{\mu}(x)=\tau\left(\mu(\alpha \lambda)^{\frac{1}{\alpha}},(\alpha \lambda)^{-1},-x\right), \mu>x \tag{6.2}
\end{equation*}
$$

The following simple fact is valid in our situation.
Theorem 3. The stopping time $T_{\mu}(x)$ has the same distribution as $\rho\left(\tau_{\mu}(x)\right)$ does. Proof. We define similarly $\tilde{\tau}(a, b, c)$ and $\tilde{\tau}_{\mu}(x)$ by replacing $S_{t}$ by $\tilde{S}_{t}$ in (6.1) and (6.2). Next, starting with (1.6), we calculate

$$
\begin{aligned}
T_{\mu}(x) & =\inf \left\{t: x+Y_{t} \geqq \mu e^{\lambda t}\right\} \\
& =\inf \left\{\rho(s): Y_{\rho(s)} \geqq \mu e^{\lambda \rho(s)}-x\right\} \\
& =\inf \left\{\rho(s): \tilde{S}_{s} \geqq \mu(1+\alpha \lambda s)^{\frac{1}{\alpha}}-x\right\}=\rho\left(\tilde{\tau}_{\mu}(x)\right) .
\end{aligned}
$$

The statement of the theorem follows from Lemma 1 which says the distribution of $\tilde{\tau}_{\mu}(x)$ coinsides with the distribution of $\tau_{\mu}(x)$.

From Theorem 3 and Proposition 2 we deduce the following result of A.Novikov
(see [6], Theorem 1).
Theorem 4. For every $a, b, c$ with $b \geqq 0, a b^{\frac{1}{\alpha}}+c>0$, one has

$$
\begin{equation*}
E(\tau(a, b, c)+b)^{\nu}=b^{\nu} \cdot \frac{H\left(\nu, \alpha,-c b^{-\frac{1}{\alpha}} \cdot d\right)}{H(\nu, \alpha, a d)}, \text { if } b>0 \text { and } \nu<\nu_{\alpha}(a d) \text {, } \tag{6.3}
\end{equation*}
$$

and

$$
E\left(\tau(a, b, c)^{v}\right)=\left\{\begin{array}{cc}
\frac{(c d)^{\alpha v}}{H(v, \alpha, a d)} & , \text { if } \quad v<v_{\alpha}(a d),  \tag{6.4}\\
+\infty, & \text { if } \quad v \geqq v_{\alpha}(a d),
\end{array}\right.
$$

where $d=(\alpha \bar{\sigma})^{\frac{1}{\alpha}}$.
Proof. Assume $b>0$ and put $x=-c, \lambda=(\alpha b)^{-1}, \mu=a b^{\frac{1}{\alpha}}$. Then

$$
\mu-x=a b^{\frac{1}{\alpha}}+c>0, \quad \bar{\mu}=\left(\frac{\lambda}{\sigma}\right)^{\frac{1}{\alpha}} \cdot \mu=a d
$$

and by Proposition 2 with $v=-\frac{\theta}{\alpha \lambda}$ we get the equalities

$$
\begin{aligned}
& E(\tau(a, b, c)+b)^{\nu}=E\left(\tilde{\tau}_{\mu}(x)+\frac{1}{\alpha \lambda}\right)^{\nu} \\
= & b^{\nu} \cdot E\left(\alpha \lambda \tilde{\tau}_{\mu}(x)+1\right)^{\nu}=b^{\nu} \cdot E \exp \left\{\nu \ln \left(1+\alpha \lambda \tilde{\tau}_{\mu}(x)\right)\right\} \\
= & b^{\nu} \cdot E \exp \left\{-\theta \rho\left(\tilde{\tau}_{\mu}(x)\right)\right\}=b^{\nu} \cdot \frac{H\left(\nu, \alpha,-c b-\frac{1}{\alpha} \cdot d\right)}{H(\nu, \alpha, a d)}
\end{aligned}
$$

provided that $\theta>-\alpha \lambda \nu_{\alpha}(\mathrm{ad})$ (or $\nu<\nu_{\alpha}(\mathrm{ad})$ ). The rest statements of the theorem follow from the properties of $H(\nu, \alpha, x)$, the case $b=0$ being taken into account by letting $b \downarrow 0$ (or $\lambda \rightarrow+\infty$ ).

Remark 3. In the original theorem of Novikov (with $d=1$, see [6]) one makes use of the fact that

$$
(t+b)^{v} \cdot H\left(v, \alpha, \frac{S_{t}-c}{(t+b)^{\frac{1}{\alpha}}}\right), t \geqq 0, b>0,
$$

is a complex-valued martingale (w.r.t. $\mathrm{F}_{\mathrm{t}^{\prime}}^{\mathrm{S}} \mathrm{t} \geq 0$ ) for every complex $v$ with $\operatorname{Re} v<1$. This fact involves an analytical continuation in contrast to our Theorem 1 .

As a second example we consider an optimal stopping problem originally treated in more general setting in [5], [7] and [9]. This problem admits a simple solution in terms of stopping times $T_{\mu}(x)$.

Under the hypothesis stated at the beginning of this section ( $H_{\alpha}$ and $\overline{\mathrm{b}}=0$ ) the quantity

$$
\begin{equation*}
\mathrm{v}(\mathrm{x}, \mathrm{~b}, \tau)=\mathrm{E} \frac{\mathrm{x}+\mathrm{S}_{\tau}}{\mathrm{b}+\tau}, \mathrm{b}>0,-\infty<\mathrm{x}<\infty, \tag{6.5}
\end{equation*}
$$

is to be maximized on stopping times $\tau=\tau(\omega)$ w.r.t. $F_{t^{\prime}}, t \geq 0$.
By Lemma 1 we have

$$
v(x, b, \tau)=v(x, b, \tilde{\tau})=E \frac{x+\tilde{S}_{\tilde{\tau}}}{b+\tilde{\tau}},
$$

using $\tilde{S}_{t}=Y_{\rho(t)}, t \geq 0$, and $\tilde{\tau}$ in the place of $S_{t}, t \geqq 0$, and $\tau$. Now taking $\lambda=\frac{1}{\alpha \mathrm{~b}}$ and $t=\delta(s), s \geq 0$, we get
$\frac{x+\tilde{S}_{t}}{b+t}=\frac{x+Y_{\rho(t)}}{b+t}=\frac{e^{\lambda \rho(t)} \cdot X_{\rho(t)}}{\frac{1}{\alpha \lambda}+t}=\frac{e^{\lambda s} \cdot x_{s}}{\frac{1}{\alpha \lambda} e^{\alpha \lambda S}}=\alpha \lambda e^{-(\alpha-1) \lambda s} \cdot X_{s}$.
Consequently, it is equivalent to consider the problem of maximizing the quantity

$$
\begin{equation*}
V(x, b, T)=\frac{1}{b} E e^{-\beta T} \cdot X_{T}, \beta=\frac{\alpha-1}{\alpha b}>0 \tag{6.6}
\end{equation*}
$$

on stopping times $T=T(\omega)$ w.r.t. $F_{s^{\prime}}^{X}, s \geqq 0$, provided that $T=\rho(\tau)$, because $V(x, b, T)=V(x, b, \tau)$.

By [7] for $\alpha=2$ and [5] for $1<\alpha<2$ one knows the solution of the original problem of maximizing (6.5) is one of the stopping times $\tau(a, b,-x)$ or the stopping time $\tau_{0}=0$.

Let us denote

$$
\left.\Psi(\mu)=\frac{\int_{0}^{\infty} y^{\alpha-2} \exp \left(\mu y-\overline{\sigma b y}{ }^{\alpha}\right) d y}{\int_{0}^{\infty} y^{\alpha-1} \exp (\mu y-\overline{\sigma b y}}{ }^{\alpha}\right) d y \quad,-\infty<\mu<\infty .
$$

As far as $\Psi(\mu)$ is positive, decreasing and continuous and $\Psi(0)=\Gamma\left(\frac{\alpha-1}{\alpha}\right)>0$, the equation $\mu=\Psi(\mu)$ has a unique solution $\tilde{\mu}($ moreover, $0<\tilde{\mu}<\Psi(0))$. The corresponding result in our case is given below without proof because it can be justified as in [5] and [7] (see also [9], Example 2, for the case $\alpha=2$ and $\lambda=1$ ). Theorem 5. For every real $x$ and $b>0$, either the stopping time $T_{\tilde{\mu}}(x)$, or the stopping time $\mathrm{T}_{\mathrm{x}}(\mathrm{x})=0$ maximizes the quantity (6.6). More precisely,

$$
\sup _{T} V(x, b, T)=V\left(x, b, T \tilde{\mu}_{\tilde{\mu}}\right)=\frac{\tilde{\mu}}{b} \gamma_{\tilde{\mu}}(\beta, x) \quad \text { if } x \leqq \tilde{\mu},
$$

and

$$
\sup _{T} V(x, b, T)=V(x, b, 0)=\frac{x}{b} \quad \text { if } x>\tilde{\mu} .
$$

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## References

1. J.-M. Bismut, Calcul des variations et processus de sauts, Z.Wahr.verw.Geb. 63 (1983) ,No. 2,pp.147-236.
2. D.A.Darling and A.J.F.Siegert, The first passage problem for a continuous Markov process, Ann.Math.Statist. 24 (1953), pp.624-639.
3. J.Jacod, Calcul stochastique et problèmes de martingales, Lecture Notes in Math. vol.714, Springer-Verlag, 1979.
4. E.Lukacs, Stochastic convergence, Acad.Press, 1975.
5. V.Mackevicius, On some problems of optimal stopping of stable stochastic processes, Lietuvos Mat.Rink. 12 (1972), No.1,pp.173-180 (in Russian).
6. A.A.Novikov, Martingale approach to the firs passage time problems for nonlinear boundaries, Proc.Steklov Math.Inst. 158 (1981), pp.130-152 (in Russian).
7. L.A.Shepp, Explicit solutions to some problems of optimal stopping, Ann.Math. Statist. 40 (1969) ,No.3,pp.993-1010.
8. A.V.Skorokhod, Random processes with independent increments, Nauka, 1964 (in Russian) .
9. H.M.Taylor, Optimal stopping in a Markov process, Ann.Math.Statist. 39 (1968), pp.1333-1344.

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