

The harmonic functions of $(A, B)^1$

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1. Introduction

The non-negative harmonic functions of a transient Markov process yield a great deal of information about the 'behaviour at infinity' of the process, and can be used to h -transform the process to behave in a certain way at infinity. The traditional analytic way of studying the non-negative harmonic functions is to construct the Martin boundary of the process (see, for example, Meyer[4], Kunita and T. Watanabe[3], and Kemeny, Snell & Knapp[2], Williams[7] for the chain case). However, certain conditions on the process need to be satisfied, one of the most basic of which is that there exists a reference measure η such that $U_\lambda(x, \cdot) \ll \eta$ for all $\lambda > 0$, all $x \in E$, the state space of the Markov process. (Here, $(U_\lambda)_{\lambda > 0}$ is the resolvent of the process.)

An interesting example proposed by Erwin Bolthausen arises when we take a standard one-dimensional Brownian motion $(B_t)_{t \geq 0}$, and define

$$A_t \equiv \int_0^t I_{(0, \infty)}(B_s) ds.$$

Then the process $(X_t)_{t \geq 0} \equiv ((A_t, B_t))_{t \geq 0}$ with values in $E = \mathbb{R}^+ \times \mathbb{R}$ is a Feller-Dynkin process, and is transient. However, if this process starts at any point (a, y) with $y < 0$, then for any $t > 0$ we have $\mathbb{P}^{(a, y)}(A_t = a) > 0$, while $\mathbb{P}^{(a, y)}(A_t = b) = 0$ for $b \neq a$. It is easy then to see that for this example there can be no reference measure η with respect to which all the resolvent kernels have a density, so the question of discovering the Martin boundary is ill-posed. Nonetheless, there are harmonic functions for X , and invariant functions too. We shall characterize all invariant functions h , and shall give representations of all harmonic functions, though the exact class remains mysterious.

To be precise about our definitions, if (P_t) denotes the semigroup of X , we shall say that a function $h: E \rightarrow \mathbb{R}^+$ is

invariant if $P_t h = h$ for all $t \geq 0$; (1.i)

harmonic if $h(X_t)$ is a P^x -local martingale for all $x \in E$; (1.ii)

excessive if $P_t h \leq h$ for all $t \geq 0$ and $P_t h \uparrow h$ as $t \downarrow 0$. (1.iii)

We have always invariant \Rightarrow harmonic \Rightarrow excessive, because h is invariant if and only if $h(X_t)$ is a P^x -martingale for all x , and h is excessive if and only if $h(X_t)$ is a P^x -supermartingale for all x .

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The main result is the following:

THEOREM 1. *Suppose that $h: E \rightarrow \mathbb{R}^+$ is harmonic. Then there exists some measurable $\rho: \mathbb{R}^+ \rightarrow \mathbb{R}^+$ such that for $t \geq 0 \geq x$*

$$h(t, x) = h(t, 0) - x\rho(t). \quad (2)$$

The function ρ satisfies the integrability condition

$$g(a, 0) \equiv \int_0^\infty \rho(a+s) \frac{ds}{\sqrt{(2\pi s)}} < \infty \quad \text{for all } a \geq 0. \quad (3)$$

Defining

$$g(t, x) \equiv \begin{cases} \int_0^\infty \rho(t+s) e^{-x^2/2s} \frac{ds}{\sqrt{(2\pi s)}} & (x > 0) \\ g(t, 0) - x\rho(t) & (x \leq 0) \end{cases} \quad (4)$$

the function g is invariant, and

(i) if h is invariant, then we have the representation

$$h(t, x) - g(t, x) = h_i(t, x) \equiv \int_0^\infty \exp(-\theta^2 t/2) \cosh(\theta x^+) \mu(d\theta) \quad (5)$$

for some non-negative measure μ on \mathbb{R}^+ satisfying the integrability condition

$$\int_0^\infty \exp(c\theta) \mu(d\theta) < \infty \quad \text{for all } c \in \mathbb{R}; \quad (6)$$

(ii) if h is harmonic, then

$$\begin{aligned} h(t, x) - g(t, x) &= h_i(t, x) + \iint_{(t, \infty) \times \mathbb{R}} p_{s-t}(x^+, y) \nu(ds, dy) \\ &= h_i(t, x) + G\nu(t, x), \end{aligned} \quad (7)$$

where h_i is as at (5), and ν is a non-negative measure on E , symmetric under the map $(t, x) \mapsto (t, -x)$, and $p(\cdot, \cdot)$ is the Brownian transition density.

Any function h represented as $h = g + h_i$, with g given by (4) and h_i by (5) is an invariant function.

Remarks. (i) Not every function $h = h_i + g + G\nu$ is necessarily harmonic, because the potential $G\nu$ is in general excessive but not harmonic. It seems to be hard in general to describe the measures ν for which $G\nu$ is harmonic, but certainly if ν is concentrated on a finite set, then $G\nu$ is harmonic.

(ii) The description of the h -transformed process is peculiar. The representation (5) corresponds to picking a drift θ according to law $\mu(d\theta)/\mu(\mathbb{R}^+)$, and h -transforming according to that while B is positive. The effect of g is to transform the process below 0 into a Bessel (3) process, pushed away from the origin at $g(A_t, 0)\rho(A_t) > 0$. Since both the upward-drifting Brownian motion and the Bessel (3) process are transient, eventually one or other prevails, and the particle either drifts off to $+\infty$ at a linear rate, or else goes out to $-\infty$ like a Bessel (3) process, with the value of A frozen forever.

(iii) Notice that, while g will certainly be continuous in $[0, \infty) \times (0, \infty)$, there is no reason why it need be continuous in $[0, \infty) \times [0, \infty)$ as ρ may be quite badly behaved. However, if ρ is continuous at t , it is easy to show that

$$\lim_{\epsilon \downarrow 0} \frac{\partial g}{\partial x}(t, \epsilon) = -\rho(t) = \frac{\partial g}{\partial x}(t, 0-)$$

so that the x -derivative of g is continuous across the boundary.

The plan of the proof of Theorem 1 is as follows.

In Section 2, we establish (2), and that g defined by (3)–(4) is invariant and dominated by h . We thereby reduce the problem to a situation where the harmonic function $\tilde{h} \equiv h - g$ is constant along any line $\{a\} \times (-\infty, 0)$. Then $\tilde{h}(A_t, B_t) = \tilde{h}(A_t, B_t^+)$, and we time change by the inverse of A to obtain the result that $\tilde{h}(t, |B_t|)$ is a local martingale; this reduces the problem to a characterization of harmonic functions for space-time Brownian motion, and we deal with this in Section 3.

However, it turns out that in the case where h is invariant, not only is $\tilde{h}(A, |B_t|)$ a local martingale, but it is also a martingale. We prove this in Section 4.

2. The basic decomposition

Let us observe that if we start at (a, y) with $y < 0$ and stop at $H_0 = \inf\{u : B_u = 0\}$, then

$$h(A_{t \wedge H_0}, B_{t \wedge H_0}) = h(a, B_{t \wedge H_0}) \quad \text{is a martingale,}$$

and so

$$h(a, y) = h(a, 0) - y\rho(a), \quad y \leq 0,$$

for some $\rho(a) \geq 0$.

We now define a function

$$g(a, y) = \lim_{n \rightarrow \infty} E^{(a, y)}[n\rho(A(H_{-n}))].$$

Is this well defined? First, observe that if we set

$$g_n(a, y) = E^{(a, y)}[n\rho(A(H_{-n}))]$$

then

$$0 \leq g_n(a, y) \leq E^{(a, y)}[h(A(H_{-n}), -n)] \leq h(a, y)$$

using Fatou's Lemma for the last inequality: so there is no problem about finiteness of the g_n . Next, for $y \geq 0$,

$$E^{(a, y)}\phi(A(H_{-n})) = \int_0^\infty \frac{y e^{-y^2/2t}}{\sqrt{(2\pi t^3)}} dt \int_0^\infty \frac{d\nu}{n} e^{-\nu/n} \int_0^\infty \frac{\nu e^{-\nu^2/2s}}{\sqrt{(2\pi s^3)}} ds \phi(a+t+s) \quad (8)$$

because the Brownian motion has to get down to 0 (at time t) and then keep going till it hits $-n$. The local time at 0 when this happens is $V \sim \exp(1/2n)$, so the amount of time spent above 0 at that time is the same in law as the first passage time to $V/2$ for BM. That is where formula (8) above comes from!

Thus for $y \geq 0$

$$E^{(a, y)}[n\rho(A(H_{-n}))] = \int_0^\infty q(y, t) dt \int_0^\infty d\nu e^{-\nu/n} \int_0^\infty q(\nu, s) ds \rho(a+t+s)$$

where we abbreviate the Brownian first-passage density to $q(\cdot, \cdot)$. Now it is clear that as $n \rightarrow \infty$, this remains bounded by $h(a, y)$ and increases to a limit.

$$\begin{aligned} g(a, y) &= \int_0^\infty q(y, t) dt \int_0^\infty \rho(a+t+s) \frac{ds}{\sqrt{(2\pi s)}} \\ &= \int_0^\infty \rho(a+u) \frac{e^{-y^2/2u}}{\sqrt{(2\pi u)}} du \\ &\uparrow \int_0^\infty \rho(a+u) \frac{du}{\sqrt{(2\pi u)}} \quad (y \downarrow 0). \end{aligned} \quad (9)$$

Thus the integrability condition

$$\int_0^\infty \rho(a+s) ds/\sqrt{s} < \infty \quad \text{for all } a \geq 0 \quad (10)$$

is necessary for h to be harmonic, and is sufficient for us to define a function $g(a, y)$ (at least for $y \geq 0$) by (9).

This justifies the definition of g for $y \geq 0$. For $y < 0$, let us assume $n \geq |y|$, and then

$$g_n(a, y) = \frac{|y|}{n} n\rho(a) + \left(1 - \frac{|y|}{n}\right) g_n(a, 0) \uparrow -y\rho(a) + g(a, 0) \quad (\text{as } n \uparrow \infty).$$

Thus the function g is well defined, and

$$g(a, y) = g(a, 0) - y\rho(a) \quad \text{for } y \leq 0.$$

If we set $g_n(a, y) = g_n(a, -n)$ for $y \leq -n$, then the g_n increase everywhere, and we shall next prove that the limit g is invariant. For this,

$$\begin{aligned} E^{(a, y)}[g(A_t, B_t)] &= \uparrow \lim_n E^{(a, y)}[g_n(A_t, B_t)] \\ &= \uparrow \lim_n \{E^{(a, y)}[g_n(A_t, B_t) : t \leq H_{-n}] + E^{(a, y)}[g_n(A_t, B_t) : t > H_{-n}]\}. \end{aligned}$$

The second term is negligible since

$$E^{(a, y)}[g_n(A_t, B_t) : t > H_{-n}] \leq E^{(a, y)}[h(A_t, B_t) : t > H_{-n}] \downarrow 0 \quad \text{as } n \rightarrow \infty$$

since $h(A_t, B_t) \in L^1$. Also,

$$\begin{aligned} E^{(a, y)}[g_n(A_t, B_t) : t \leq H_{-n}] &= E^{(a, y)}[n\rho(A(H_{-n})) : t \leq H_{-n}] \\ &= g_n(a, y) - E^{(a, y)}[n\rho(A(H_{-n})) : t > H_{-n}]. \end{aligned}$$

All that is needed now is to prove that for $t > 0$ fixed

$$E^{(a, y)}[n\rho(A(H_{-n})) : H_{-n} < t] \equiv \Psi_n(t) \xrightarrow{n \rightarrow \infty} 0.$$

Since $\Psi_n(\cdot)$ is clearly increasing, it will be sufficient to prove that, with $\lambda > 0$ fixed,

$$\int_0^\infty \lambda e^{-\lambda t} E^{(a, y)}[n\rho(A(H_{-n})) : H_{-n} < t] dt \equiv E^{(a, y)}[n\rho(A(H_{-n})) : H_{-n} < T] \rightarrow 0 \quad (\text{as } n \rightarrow \infty) \quad (11)$$

where T is exponentially distributed with parameter λ , independent of B . Now for $y > -n$,

$$P^{(a, y)}[H_{-n} < T] = \exp(-(y+n)\sqrt{(2\lambda)}),$$

and conditional on $\{H_{-n} < T\}$, $(B_t: 0 \leq t \leq H_{-n})$ is identical in law to

$$(B_t - \theta t: 0 \leq t \leq \sigma_{-n}),$$

where $\theta \equiv \sqrt{2\lambda}$ and σ_{-n} is the first time $B_t - \theta t$ reaches $-n$ (see Williams[6]).

So now we want to compute the law of

$$A(H_{-n}) \equiv \zeta_n \equiv \int_0^{\sigma_{-n}} I_{[0, \infty)}(B_t - \theta t) dt.$$

However, if $\psi(y) = E^y[e^{-\alpha \zeta_n}]$, then ψ must satisfy

$$\frac{1}{2}\psi'' - \theta\psi' - \alpha\psi = 0 \quad \text{in } (0, \infty)$$

$$\frac{1}{2}\psi'' - \theta\psi' = 0 \quad \text{in } (-n, 0)$$

$$\psi(-n) = 1$$

together with the condition that ψ is C^1 at 0. A few simple calculations yield the solution

$$\begin{aligned} \psi(y) &= c e^{-\gamma y} \quad (y \geq 0) \\ &= c \left\{ 1 + \frac{\alpha}{2\theta} (1 - e^{2\theta y}) \right\} \quad (-n \leq y \leq 0), \end{aligned} \quad (12)$$

where $\gamma \equiv \sqrt{(\theta^2 + 2\alpha) - \theta}$, $c^{-1} \equiv 1 + (\gamma/2\theta)(1 - e^{-2\theta n})$. For the time being, we restrict our attention to starting values $y \geq 0$. What (12) tells us is that the P^y -distribution of ζ_n is the same as the time taken for $B_t - \theta t$ to drop from y to $-V_n$, where V_n is exponentially distributed with mean $(1 - e^{-2\theta n})/2\theta$. This comes as no surprise to anyone who has understood the path decompositions of Williams, and the excursion theory of drifting Brownian motion (see VI.55 in Rogers and Williams[5]).

Thus for $y \geq 0$, with $q_n \equiv 2\theta(1 - e^{-2\theta n})^{-1}$,

$$\begin{aligned} E^{(a, y)}[n\rho(A(H_{-n})) | H_{-n} < T] \\ = \int_0^\infty q_n e^{-q_n x} dx \int_0^\infty (x+y) e^{-(x+y-\theta t)^2/2t} \frac{dt}{\sqrt{(2\pi t^3)}} n\rho(a+t). \end{aligned}$$

Since the q_n decrease to 2θ , we have the upper bound

$$q_1 \int_0^\infty e^{-2\theta x} dx \int_0^\infty (x+y) e^{-(x+y-\theta t)^2/2t} \frac{dt}{\sqrt{(2\pi t^3)}} n\rho(a+t),$$

which is finite if and only if it is finite for $y = 0$. Putting $y = 0$,

$$\begin{aligned} &\int_0^\infty e^{-2\theta x} dx \int_0^\infty x e^{-(x-\theta t)^2/2t} \frac{dt}{\sqrt{(2\pi t^3)}} n\rho(a+t) \\ &= \int_0^\infty dx \int_0^\infty x e^{-(x+\theta t)^2/2t} \frac{dt}{\sqrt{(2\pi t^3)}} n\rho(a+t) \\ &= \int_0^\infty \frac{dt}{\sqrt{t}} \rho(a+t) \int_0^\infty \zeta \exp(-(\zeta + \theta\sqrt{t})^2/2) \frac{d\zeta}{\sqrt{(2\pi)}} \\ &\leq \int_0^\infty \frac{dt}{\sqrt{t}} \rho(a+t) \int_0^\infty (\zeta + \theta\sqrt{t}) \exp(-(\zeta + \theta\sqrt{t})^2/2) \frac{d\zeta}{\sqrt{(2\pi)}} \\ &\leq \int_0^\infty \frac{dt}{\sqrt{(2\pi t)}} \rho(a+t) < \infty \end{aligned}$$

from (3). Since $P[H_{-n} < T] \rightarrow 0$ as $n \rightarrow \infty$, we deduce (10), at least when $y \geq 0$. But for $y < 0$, $n > -y$,

$$E^{(a, y)}[n\rho(A(H_{-n})) : H_{-n} < T] = E^{(a, 0)}[n\rho(A(H_{-n})) : H_{-n} < T] \cdot P^{(a, y)}[H_0 < H_{-n} \wedge T] \\ + n\rho(a)P^{(a, y)}[H_{-n} < H_0 \wedge T],$$

and it is a simple matter now to deduce that this goes to 0. Hence the function g is invariant for X .

3. Representing $h - g$

Since g is invariant and $g \leq h$, it follows that $\tilde{h} \equiv h - g$ is harmonic, that is, $\tilde{h}(X_t)$ is a local martingale. But since $\tilde{h}(a, x) = \tilde{h}(a, x^+)$, if we let

$$T_n \equiv \inf\{t : \tilde{h}(X_t) > n\}$$

then certainly $X(T_n)$ must be in $\mathbb{R}^+ \times \mathbb{R}^+$. So if we define the time change $\tau_t \equiv \inf\{u : A_u > t\}$ then $\tilde{h}(X(\tau_t))$ is a local martingale in the (\mathcal{F}_{τ_t}) -filtration, reduced by the stopping times $A(T_n)$. However, $(X(\tau_t))_{t \geq 0}$ has the same distribution as $((t, |B_t|))_{t \geq 0}$, so if we redefine

$$\tilde{h}(t, y) = \tilde{h}(t, -y) \quad (y \leq 0)$$

we obtain the conclusion that

$$\tilde{h}(t, B_t) \text{ is a local martingale.}$$

So we now address the task of characterizing all functions $h : \mathbb{R}^+ \times \mathbb{R} \rightarrow \mathbb{R}^+$ which are symmetric in that $h(t, y) = h(t, -y)$ for all (t, y) and such that $h(t, B_t)$ is a local martingale under every P^x . Observe that this means that h is excessive.

To finish things reasonable directly, we appeal to results of Meyer [4]. Using the results III.T13, I.T16, III.T9, we may represent excessive h_e as

$$h_e(x) = h_0(x) + \int g(x, x') \mu(dx') \quad (13)$$

for some measure μ such that

$$u(x) \equiv \int g(x, x') \mu(dx') < 0 \quad \text{for all } x, \quad (14)$$

where

$$g((s, y), (s', y')) \equiv I_{\{s' > s\}} p_{s'-s}(y, y'),$$

with $p(\cdot, \cdot)$ denoting the Brownian transition density. The function h_0 in the representation (13) has the property

$$h_0(x) = E^x[h_0(X(\tau_K))] \quad \text{for all compact } K \quad (15)$$

where $\tau_K \equiv \inf\{t > 0 : X_t \notin K\}$. This property is what Meyer [4] calls harmonic, but this does not agree with our definition. We do, however, have the following result.

PROPOSITION. *If h_0 satisfies (15) then h_0 is invariant.*

Proof. For ease of notation, we consider only starting points on $\{0\} \times \mathbb{R}$, and (for the purposes of this proof only) abbreviate $(0, y)$ to y . Thus for any $t > 0$ and $-K < y < N$,

$$h_0(y) = E^y[h_0(X(H_{-K} \wedge H_N \wedge t))] \\ = E^y[h_0(X_t) : t < H_{-K} \wedge H_N] + E^y[h_0(X(H_N)) : H_N < t \wedge H_{-K}] \\ + E^y[h_0(X(H_{-K})) : H_{-K} < t \wedge H_N].$$

The first of these three terms tends to $E^y[h_0(X_t)]$ as $N, K \rightarrow \infty$, so we must show that the last two tend to zero. By letting $K \rightarrow \infty$, we learn that $E^y[h_0(X(H_N)) : H_N < t] \leq h_0(y)$ and so $\int_0^t q(N-y, s) h(s, N) ds \leq h_0(y)$. But if we consider starting from $y+1$, we learn that

$$\begin{aligned} h_0(y+1) &\geq \int_0^t (N-y-1) \exp\{-(N-y-1)^2/2s\} \frac{h(s, N)}{\sqrt{(2\pi s^3)}} ds \\ &\geq \frac{N-y-1}{N-y} \int_0^t (N-y) \exp\{-(N-y)^2/2s + (N-y)/2s\} \frac{h(s, N)}{\sqrt{(2\pi s^3)}} ds \\ &\geq \frac{N-y-1}{N-y} e^{(N-y)/2t} \int_0^t q(N-y, s) h(s, N) ds. \end{aligned}$$

It follows immediately that

$$E^y[h_0(X(H_N)) : H_N < t] \rightarrow 0 \quad (\text{as } N \rightarrow \infty)$$

and the invariance of h_0 is established.

Now it is well known that the invariant functions h_i for space-time Brownian motion are all of the form

$$h_i(t, x) = \int \exp(\theta x - \theta^2 t/2) \nu(d\theta)$$

for some ν satisfying the integrability condition (6) and that h_i is symmetric in x if and only if ν is symmetric.

We have now established that any harmonic function of (A_t, B_t) can be represented as at (7) for some invariant h_i , and the potential of some measure ν . Not every measure ν will give a harmonic function, of course and it appears difficult to characterise the ν for which we do get a harmonic function. One example is where we take ν to be the unit mass at $(1, 0)$, and then we have the harmonic function

$$h(a, y) = (2\pi(1-a))^{-1/2} \exp\{-(y^+)^2/2(1-a)\} I_{\{a < 1\}}.$$

This is a particularly interesting example, because if we h -transform using h , we obtain a process (A_t, Y_t) satisfying

$$dY_t = dW_t - \frac{Y_t^+}{1-A_t} dt, \quad dA_t = I_{\{Y_t > 0\}} dt,$$

and it is not clear how this process behaves; it is like Brownian motion when $Y < 0$, and like Brownian bridge when $Y > 0$, but does it have finite or infinite lifetime? As A approaches 1, the excursions of Y from 0 into $(0, \infty)$ get shorter and shorter, so it is conceivable that the lifetime of the process could be infinite, as the bulk of time is spent with $Y < 0$. We shall prove that this is not in fact the case, by observing first that $Z_t \equiv Y_t(1-A_t)^{-1}$ satisfies

$$dZ_t = \frac{dW_t}{1-A_t}$$

and so if σ is inverse to the continuous increasing process

$$\gamma_t \equiv \int_0^t (1-A_s)^{-2} ds,$$

then $\tilde{Z}_t = Z(\sigma_t)$ is a Brownian motion. We now express the time change in terms of \tilde{Z} . We have

$$\dot{\sigma}_t = (1 - A(\sigma_t))^2 = \left(1 - \int_0^{\sigma_t} I_{\{Y_u > 0\}} du\right)^2 = \left(1 - \int_0^t \dot{\sigma}_s I_{\{\tilde{Z}_u > 0\}} du\right)^2.$$

Thus if $\beta_t \equiv I_{\{\tilde{Z}_t > 0\}}$, a little elementary calculus gives us

$$\dot{\sigma}_t = \left(1 + \int_0^t \beta_u du\right)^2.$$

Now Hobson[1] proves that if $f: \mathbb{R}^+ \rightarrow \mathbb{R}^+$ is decreasing, and $\int^\infty t^{-1} \sqrt{f(t)} dt < \infty$, then $\liminf_{t \rightarrow \infty} A_t/t f(t) = +\infty$. Taking $f(t) = t^{-\epsilon}$, we see that $A_t \geq t^{1-\epsilon}$ for all large enough t , and so $\dot{\sigma}$ is integrable. The conclusion is that $\sigma_\infty < \infty$ and γ explodes in finite time, that is, A reaches 1 in finite time!

4. The invariant case

The first thing to prove is that the potential term $G\nu$ in (7) cannot be invariant for X . However, this is almost obvious if we re-express it as

$$G\nu(t, x) = \iint_{(t, \infty) \times (0, \infty)} \{p_{s-t}(x^+, y) + p_{s-t}(x^+, -y)\} \nu(ds, dy) + \int_{(t, \infty)} p_{s-t}(x^+, 0) \nu(ds, \{0\}).$$

Indeed, $\gamma(t, x; s, y) \equiv \{p_{s-t}(x^+, y) + p_{s-t}(x^+, -y)\} I_{\{s > t\}}$

is the density with respect to Lebesgue measure of the Green's function of X , at least for $y \geq 0$. This is because if we only view X at times when $B \geq 0$, we see a reflecting space-time Brownian motion in $\mathbb{R}^+ \times \mathbb{R}^+$. Hence $G\nu$ is in fact a potential with respect to the semi-group of X , and so is not invariant. Next we must prove that every function h_i of the form given in (5) is invariant for X . This is not difficult if we set it up correctly.

Take a Brownian motion W , with local time L at zero, and an independent stable $(\frac{1}{2})$ subordinator Z :

$$E \exp(-\alpha Z_t) = \exp(-t \sqrt{2\alpha})$$

and now consider the bivariate Markov process $(|W_t|, t + Z(L_t))_{t \geq 0}$.

Let $(\mathcal{G}_t)_{t \geq 0}$ be the filtration of this process. We then have that $M_t \equiv h_i(t, |W_t|)$ is a (\mathcal{G}_t) -martingale. If we now time change by

$$\sigma_t \equiv \inf\{u : u + Z(L_u) > t\}$$

then always $\sigma_t \leq t$, and so $M(\sigma_t)$ is a martingale in the filtration $(\mathcal{G}_{\sigma(t)})$. However, by the way it has been constructed, we have the identity in law as processes

$$(\sigma_t, |W_{\sigma(t)}|)_{t \geq 0} \stackrel{D}{=} (A_t, B_t^+),$$

and so $h_i(A_t, B_t^+) \equiv h_i(A_t, B_t)$ is a martingale, which means that h_i is also invariant for $(A_t, B_t) \equiv X_t$.

Assembling this finally, if h is invariant for X , then it is harmonic for X , and so has a representation of the form (7), by the previous section. However, we have seen that g is invariant for X , and have just proved that h_i is invariant for X , so we conclude

that the potential $G\nu$ is invariant for X . This can only happen if $\nu \equiv 0$. The theorem is proved.

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