# Spectral Representation of Multiply Self-decomposable Stochastic Processes and Applications\*

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In the present paper we study multiply selfdecomposable probability measures (SDPM) and processes and prove their integral representations. Similarly, the multiple s-selfdecomposability case is treated. Our results extend some of known results due to Urbanik K.,Jurek., Rosinski J.and Rajput B.S. As an application we construct a Damped-mixed stable price process in option pricing.

### 1 Introduction, notation and preliminaries

The main aim of this paper is to prove that each multiply self-decomposable process (MSDP) on an Euclidean space admits a stochastic integral w.r.t.a MSD random measure (RM). Moreover, we will consider similar problems for multiply s-self-decomposable processes(MsSDP). Through the paper we shall denote by  $\mathcal{X}$  a fixed d-dimensional (d=1,2,...) Euclidean space with the usual inner product <,> and norm  $\|.\|$ .

Let  $\mathcal{P}(\mathcal{X})$  denote the class of all probability measures (PM) on the  $\sigma$ -field  $\mathcal{B}(\mathcal{X})$  of Borel subsets of  $\mathcal{X}$  equipped with the weak convergence. Given a positive number

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c we define on  $\mathcal{X}$  the following two families of mappings  $T_c$  and  $U_r$  as follows:

$$\begin{cases}
T_r x = rx, \\
U_r = \max(0, ||x|| - r) \frac{x}{||x||}, U_r(0) = 0.
\end{cases}$$
(1)

Further, for a PM  $\mu \in \mathcal{P}(\mathcal{X})$  and a mapping T on  $\mathcal{X}$  let  $T\mu$  denote the image of  $\mu$  under T.

Recall (cf.Loéve [12] and Sato [25]) that a PM  $\mu \in \mathcal{P}(\mathcal{X}$  is called SD if for each 0 < c < 1 there exists a PM  $\mu_c$  such that

$$\mu = T_c \mu * \mu_c \tag{2}$$

where \* denotes the ordinary convolution of PM's.

The concept of shrinking SDPM (shortly, s-SDPM) was introduced by Medgyessy [14] and studied by Jurek [3], [4], [6]. Namely, a PM  $\mu$  is called s-SD if it is ID and for each 0 < c < 1 there exists a PM  $\mu_c$  such that

$$\mu = U_c \mu^c * \mu_c \tag{3}$$

where the power is taken in the convolution sense.

It is known [3], [4], [6], [29], [17]that if  $\mu$  is SD (resp., s-SD) then  $\mu, \mu_c$  are both ID. Let us denote by  $ID(\mathcal{X})$  the class of all IDPM's on the space. The class of all SDPM's (resp., s-SDPM's) on  $\mathcal{X}$  is denoted by  $\mathcal{L}(\mathcal{X})$  (resp.,  $\mathcal{U}(\mathcal{X})$ ).

Let  $\mathcal{L}_n(\mathcal{X})$ , n = 1, 2, ... (resp.,  $\mathcal{U}_n(\mathcal{X})$ , n = 1, 2, ...) denote the class of all n-times SDPM's (resp., n-times s-SDPM's) which were first introduced by Urbanik<sup>1</sup> [29] (resp.,Jurek[4] and then studied further by many other authors (cf., for example [4], [17], [25]...).

They are defined recursively as follows: A p.m.  $\mu \in \mathcal{L}_n(\mathcal{X}), n = 2, 3, ...$  if and only if  $\mu \in \mathcal{L}_1(\mathcal{X})$  and for each  $c \in (0, 1)$  the component  $\mu_c$  in (2) belongs to  $\mathcal{L}_{n-1}(\mathcal{X})$ .

It has been proved by Nguyen ([18], Proposition 1.1) that  $a \ p.m. \ \mu \in \mathcal{L}_n(\mathcal{X}), n = 1, 2, ..., if and only if, for every <math>c \in (0, 1)$  there exists  $a \ p.m. \ \nu := \mu_{c,n} \in ID(\mathcal{X})$  such that the following equality holds:

$$\mu = *_{k=0}^{\infty} (T_{c^k} \nu)^{*r_{k,n}} \tag{4}$$

where the power is taken in the convolution sense and, for  $n=1,2,\ldots$ ;  $k=0,1,2,\ldots$  we put

$$r_{k,n} = \binom{n+k-1}{k}. (5)$$

<sup>&</sup>lt;sup>1</sup>It should be noted, that our notation  $\mathcal{L}_n(\mathcal{X})$  used here and in references [17], [18] is other than that in Urbanik and other Authors [4], [25]. In particular, in our notation,  $\mathcal{L}_1(\mathcal{X})$  denotes the set of all SDPM's on  $\mathcal{X}$  while in [4], [25] this class was denoted by  $\mathcal{L}_0(\mathcal{X})$ .

The formulas (3) and (4) lead to the following interpolation of classes  $\mathcal{L}_n(\mathcal{X})$  (cf. Nguyen [17]): For each  $\alpha > 0$  we put

$$\begin{pmatrix} \alpha \\ k \end{pmatrix} = \begin{cases} 1 & k = 0, \\ \alpha(\alpha - 1)...(\alpha - k + 1)/k! & k = 1, 2, ... \end{cases}$$
 (6)

and introduce the class  $\alpha$ -times SDPM's, shortly,  $\alpha - SDPM's$  as the following:

**Definition 1.1 (cf. Nguyen [19])** A p.m.  $\mu \in \mathcal{L}_{\alpha}(\mathcal{X}), \alpha > 0$ , if and only if, for every  $c \in (0,1)$  there exists a p.m.  $\nu := \mu_{c,\alpha} \in ID(\mathcal{X})$  such that the following equality holds:

$$\mu = *_{k=0}^{\infty} (T_{c^k} \nu)^{r_{k,\alpha}} \tag{7}$$

where, the power is taken in the convolution sense and, for  $\alpha > 0, k = 0, 1, 2, ...$  we put

$$r_{k,\alpha} = \binom{\alpha + k - 1}{k} \tag{8}$$

It should be noted [17] that the infinite convolution (8) is weakly convergent if and only if

$$\int_{\mathcal{X}} \log^{\alpha}(1 + ||x||)\nu(dx) < \infty \tag{9}$$

In the sequel, we shall denote by  $ID_{log^{\alpha}}(\mathcal{X})$  the subclass of  $ID(\mathcal{X})$  of all distributions for which the condition (9) is satisfied.

Now, let us quote the following important integral representation for SDPM's due to Vervaat-Jurek [5]:

**Theorem 1.1 (Jurek-Vervaat)** A p.m.  $\mu \in \mathcal{L}_1(\mathcal{X})$  if and only if there exists a  $\mathcal{X}$ -valued Lévy process  $\{X(.)\}$  of the class  $ID_{log}$  such that

$$\mu \stackrel{d}{=} \int_0^\infty exp(-t)X(dt) \tag{10}$$

The integrator  $\{X(.)\}$  is called the background driving Lévy process (shortly, BDLP) of  $\mu$  (cf.Jurek [5, 6]) and the r.v. X(1) is called BD r.v.. Further, Nguyen N.H.[16], obtained the following pretty generalization of Theorem 1.1 to the case of  $\alpha$ -SDPM's for each  $\alpha > 0$ .

**Theorem 1.2 (Nguyen N.H** [16]) A p.m.  $\mu \in \mathcal{L}_{\alpha}(\mathcal{X})$  if and only if there exists a  $\mathcal{X}$ -valued Lévy process  $\{X_{\alpha}(t)\}$  of the class  $ID_{log^{\alpha}}(\mathcal{X})$  such that

$$\mu \stackrel{d}{=} \int_0^\infty exp(-t^{\frac{1}{\alpha}})X_\alpha(dt) \tag{11}$$

In the sequel we shall need the following representation of ch.f.'s of ID and MSDPM's on  $\mathcal{X}$ :

**Theorem 1.3 (cf[20], [24])** A p.m.  $\mu$  is ID if and only if its ch.f.  $\hat{\mu}(y), y \in \mathcal{X}$  is of the unique form:

$$\begin{cases} -log\hat{\mu}(y) = i < z, y > + < \Sigma y, y > \\ -\int_{\mathcal{X}} \left( e^{i < y, x > -1 - i\tau(x) \right) M(dx) \end{cases}$$

where  $z \in \mathcal{X}$  is fixed;  $\Sigma$  is a quadratic form on  $\mathcal{X}$  and M is a Lévy measure on  $\mathcal{X}$  characterized by the property that M(0) = 0, M is finite ouside of very neighborhood of the origin and

$$\int_{U_1} \frac{\|x\|^2}{1 + \|x\|^2} M(dx) < \infty;$$

the function  $\tau(x)$  is defined by

$$\tau(x) = \begin{cases} ||x|| & x \in U_1; \\ 1 & ||x|| > 1, \end{cases}$$

 $U_1$  being the closed unit Ball in  $\mathcal{X}$ .

In what follows, if  $\mu$  is ID with the ch.f. given by (11) then we will identify it with the triple  $[z, \Sigma, M]$ . Thus,we have

**Theorem 1.4 (Nguyen [19],Theorem 2.4)** A p.m.  $\mu \in \mathcal{L}_{\alpha}(\mathcal{X}), \alpha > 0$  if and only if  $\mu = [z, \Sigma, M]$ , where  $z, \Sigma$  are the same as in Theorem (1.2) and the Lévy's measure M is given by

$$M(A) = \begin{cases} \int_{\mathcal{X}} v_{\alpha}(x) \\ (\int_{0}^{\infty} \chi_{A}(e^{-u}x)u^{\alpha-1}du)m(dx) \end{cases}$$
 (12)

where m is a finite measure on  $\mathcal{X}$  vanishing at the origin; A is a Borel subset of the real line separated from 0; the weight function  $v_{\alpha}(x)$  is defined by

$$v_{\alpha}^{-1}(x) = \int_{0}^{\infty} \frac{e^{-2t}x^{2}}{1 + e^{-2t}x^{2}} t^{\alpha - 1} dt$$
 (13)

**Theorem 1.5 (Nguyen[17],[19])** A p.m.  $\mu$  is mixed-stable i.e.  $\mu \in \mathcal{L}_{\infty}(\mathcal{X})$  if and only if  $\mu = [z, \Sigma, M]$ , where  $z, \Sigma$  are the same as in Theorem (1.4) and the Lévy's measure M of  $\mu$  is given by

$$M(A) = \int_{V_1} \int_0^\infty \chi_A(tx) \frac{dt}{t^2|x|+1} h(x) \nu(dx)$$
 (14)

where  $\nu$  is a PM on the open unit ball  $V_1 := \{x \in \mathcal{X} : ||x|| < 1\}$  and h(x) is a nonnegative continuous weight function on  $V_1$ .

## 2 Mappings $\{T_c^{(\alpha)}\}$ and classes $\{\mathcal{L}_{\alpha}(\mathcal{X})\}$

In this section we introduce families of mappings  $\{T_c^{(\alpha)}\}$ , where  $0 < c < 1; \alpha > 0$  acting on the whole class  $ID(\mathcal{X})$  and show that they play the same role as mappings  $T_c$  in the definition of  $\alpha$ -SDPM's. To begin with let us consider the following particular cases:

### **2.1** $\alpha = n = 1, 2, ...$

Let  $\mu \in \mathcal{L}_n(\mathcal{X})$ , n = 1, 2, ... By Proposition 1.1 [17], for every 0 < c < 1 the equation (??) holds. Putting

$$T_c^n \mu = *_{k=1}^{\infty} (T_{c^k} \nu)^{r_{k,n}} \tag{15}$$

and taking into account (??) we have

$$\mu = T_c^{(n)} \mu * \mu_{c,n} \tag{16}$$

Conversely, it is also true. Namely, by induction one can prove that if a PM  $\mu$  satisfies equation (16) for each 0 < c < 1 and for a PM  $\mu_{c,n}$ , then it belongs to  $\mathcal{L}_n(\mathcal{X})$ .

#### **2.2** $0 < \alpha < 1$ .

This case was treated in [17]. Namely, for such  $\alpha$  the mapping  $T_{c,\alpha}$  is defined in [17]. Then, by Theorem 2.1 [17], it follows that a PM  $\mu$  belongs to  $\mathcal{L}_{\alpha}(\mathcal{X})$  if and only if for every 0 < c < 1 there exists a PM  $\mu_{c,\alpha}$  such that

$$\mu = T_c^{\alpha} \mu * \mu_{c,\alpha} \tag{17}$$

### **2.3** The general case $\alpha > 0$ :

It is easy to show that

$$1 = \sum_{k=1}^{\infty} (-1)^{k-1} r_{k,\alpha} = \sum_{k=1}^{\infty} |r_{k,\alpha}|$$
 (18)

Consequently, the mapping  $T_{c,\alpha}:ID(\mathcal{X})\to ID(\mathcal{X})$  given by

$$T_{c,\alpha}\mu = *_{k=1}^{\infty} T_{c^k}\mu^{|r_{k,\alpha}|} \tag{19}$$

for any 0 < c < 1 and  $\alpha > 0$  is well-defined. Furthermore, the following general theorem holds: **Theorem 2.1** A PM  $\mu \in \mathcal{L}_{\alpha}(\mathcal{X})$ ,  $\alpha > 0$ , if and only if for each 0 < c < 1 there exists a PM  $\mu_{c,\alpha}$  such that the equation (17) holds for each  $\alpha > 0$ .

**Proof.** The "if" part is similar to the proof of Theorem 2.1 [17]. To prove the "only if" part one may assume that  $\alpha = \beta + n$ , where  $0 < \beta < 1, n = 1, 2, ...$  But it is clear by virtue of the cases 2.1 and 2.2 and by noticing that the mappings  $T_{c,n}$  and  $T_{c,\beta}$  commute with each other.

Theorem 2.2 ( $\alpha$ -differentiability of  $\alpha$ -SDPM's on  $\mathcal{X}$ ) For every  $\alpha > 0$  and every  $PM \ \mu \in \mathcal{L}_{\alpha}(\mathcal{X})$  there exists a weak limit, denoted by  $\mathcal{D}^{\alpha}\mu$ , which belongs to  $ID_{log^{\alpha}}(\mathcal{X})$  and satisfies the equation

$$\mathcal{D}^{\alpha}\mu = \lim_{t \to 0} \mu_{c,\alpha}^{t^{-\alpha}} \tag{20}$$

where  $t = -\log c$ ,  $\mu_{c,\alpha}$  is as in (7) and (17).

**Proof.** See Nguyen (Theorem 2.4 [17]).

**Definition 2.1 (cf. Nguyen [17)** The limit measure  $D^{\alpha}\mu$  in Theorem (2.2) is called the  $\alpha$ -derivative of  $\mu$ .

The following Theorem is obvious:

**Theorem 2.3** For each  $\alpha > 0$  the operator  $D^{\alpha}$  stands for an algebraic isomorphism between  $\mathcal{L}_{\alpha}(\mathcal{X})$  and  $ID_{log^{\alpha}}(\mathcal{X})$ .

# 3 Mappings $\{U_c^{(\alpha)}\}\$ and classes $\{\mathcal{U}_{\alpha}(\mathcal{X})\}\$

Following verbatim the proof of cases 2.1, 2.2 and 2.3 we have the Theorem:

**Theorem 3.1** For any 0 < c < 1 and  $\alpha > 0$  and for every  $PM \mu \in ID(\mathcal{X})$  we put

$$U_{c,\alpha}\mu = *_{k=1}^{\infty} T_{c^k} \mu^{\left|\binom{\alpha}{k}\right| c^k} = *_{k=1}^{\infty} U_{c^k} \mu^{\left|\binom{\alpha}{k}\right|} \tag{21}$$

Then we get a mapping  $U_{c,\alpha}$  which stands for a well defined continuous isomorphism of the convolution algebra  $ID(\mathcal{X})$ . Moreover, restricted to  $ID(\mathcal{X})$ , it stands for an analogue of the shrinking mapping  $U_c$  in (1.1).

**Definition 3.1** A PM  $\mu \in ID(\mathcal{X})$  is said to be of the class  $U_{\alpha}(\mathcal{X}), \alpha > 0$ , or equivalently,  $\alpha$ -s-SD, if for each 0 < c < 1 the following formula holds:

$$\mu = U_{c,\alpha}\mu * \mu_{c,\alpha} \tag{22}$$

for some PM  $\mu_{c,\alpha} \in ID(\mathcal{X})$ 

From the above definition we have:

**Theorem 3.2** A PM  $\mu = [z, \Sigma, M] \in \mathcal{U}_{\alpha}(\mathcal{X}), \alpha > 0$  if and only if the Lévy measure M satisfies the following condition:

$$\sum_{k=0}^{\infty} {\binom{\alpha}{k}} |c^k T_{c^k} M \ge 0 \tag{23}$$

for each 0 < c < 1, or, equivalently,

$$\left[\left|\sum_{k=0}^{\infty}\right| \binom{\alpha}{k} | U_{c^k} M \ge 0\right] \tag{24}$$

**Definition 3.2 (cf. Jurek**[8]) Given  $\alpha > 0$  let  $G_{\alpha}$  denote a Gamma r.v. with distribution  $\tau_{\alpha}$ . Let  $\mathcal{U}^{<\alpha>}$  denote the class of all distributions of  $\int_{(0,1)} t dY_{\rho}(\tau_{\alpha}(t))$ , where  $Y_{\rho}(.)$  is a Lévy process with  $\mathcal{L}(Y_{\rho}(1)) = \rho$ .

By virtue of the above formulas (23) and (24) and Jurek [8], formula (29) we have the following theorem

**Theorem 3.3** The following equation hold:

$$\mathcal{U}_{\alpha}(\mathcal{X}) = \mathcal{U}^{\langle \alpha \rangle} \tag{25}$$

which shows that definitions 3.1 and 3.2 are equivalent.

## $\textbf{4} \quad \textbf{Stochastic representation of classes} \,\, \mathcal{U}_{\alpha}(\mathcal{X})$

**Definition 4.1** Let T be a parameter set Z of all integers or R of all real numbers. A stochastic process  $X_t, t \in T$  is said to be ID, stable, mixed-stable,  $\alpha$ -SD,  $\alpha$ -SD if for any  $t_1, t_2, ..., t_n \in T$  and  $\lambda_1, \lambda_2, ..., \lambda_n, n = 1, 2, ...$  the r.v.  $\Sigma_1^n \lambda_j X_{t_j}$  is ID, stable, mixed-stable,  $\alpha$ -SD,  $\alpha$ -S-SD, respectively.

**Definition 4.2** Let  $\Lambda = \{\Lambda(A) : A \in \mathcal{S}\}$  be a real stochastic process defined on a probability space  $(\Omega, \mathcal{F}, \mathcal{P})$ , where  $\mathcal{S}$  stands for a  $\sigma$ -ring of subsets of an arbitrary non-empty set S satisfying the following condition: There exists an increasing sequence  $S_n, n = 1, 2, ...$  of sets in S with  $\bigcap_n S_n = S$ .

We call  $\Lambda$  to be an *independently scattered random measure*(RM), if, for every sequence  $\{A_n\}$  of disjoint sets in  $\mathcal{S}$ , the random variables  $\Lambda(A_n)$ , n = 1, 2, ... are independent, and, if  $\cup_n A_n$  belongs to  $\mathcal{S}$ , then we also have

$$\Lambda(\cup_n A_n) = \Sigma_n \Lambda(A_n) \quad a.s.,$$

where the series is assumed to be convergent a.s. In addition, if for every  $A \in \mathcal{S}$  the distribution of  $\Lambda(A)$  is ID, stable, mixed-stable, MSD, respectively, then we say that it is an ID, stable, mixed-stable, MSD RM.

By virtue of Theorem 2 each r.v.  $\Lambda(A)$ ,  $A \in S$  has the ch.f.

$$\begin{cases}
-log E exp(it\Lambda(A) = it\nu_0(A) + \frac{1}{2}t^2\nu_1(A) \\
-\int_{-\infty}^{\infty} (e^{itx} - 1 - it\tau(x))F_A(dx),
\end{cases}$$
(26)

where  $t \in \mathbb{R}, A \in S$  and  $-\infty < v_0(A) < \infty, 0 \le v_1(A) < \infty$  and  $F_A$  is a Lévy measure on  $\mathbb{R}$ . Moreover,  $v_0$  is a signed measure,  $v_1$  a measure and  $F_A$  a Lévy measure.

The above representation implies the following

Theorem 4.1 (Raiput and Rosinski [24], Proposition 2.1) The characteristic function (19) can be written in the unique form:

$$Eexp(it\Lambda(A)) = exp(\int_A K(t,s)\lambda(ds))$$
 (27)

where  $t \in \mathbb{R}, A \in S$  and

$$\begin{cases}
K(t,s) = ita(s) - 1/2t^2\sigma^2(s) \\
+ \int_A (e^{itx} - 1 - it\tau(x))\rho(s, dx),
\end{cases}$$
(28)

with

$$a(s) = \frac{dv_0}{d\lambda}(s) \tag{29}$$

and

$$\sigma^2(s) = \frac{dv_1}{d\lambda}(s) \tag{30}$$

and  $\rho$  is given by Lemma 2.3 in [3]. Moreover, we have

$$|a(s)| + \int_{\mathbb{R}} \min\{1, x^2\} \rho(s, dx) = 1 \quad a.e.[\lambda]$$
 (31)

**Definition 4.3 (Urbanik and Woyczynski [27]) (a)** If f is a simple function on  $S, f = \Sigma_j x_j \chi_{A_j}, A_j \in S$  then we put, for each  $A \in \sigma(S)$ 

$$\int_{A} f d\Lambda = \Sigma_{j} \lambda(A \cap A_{j})$$

- (b) A measurable function  $f:(S,\sigma(S))\to (\mathbb{R},\mathcal{B}(\mathbb{R}))$  is said to be  $\Lambda$ -integrable if there exists a sequence  $\{f_n\}$  of simple functions as defined in (a) such that
  - (i)  $f_n \to f$  a.e. $[\lambda]$ ,
  - (ii) For every  $A \in \sigma(S)$ , the sequence  $\{\int_A f_n d\Lambda\}$  converges in prob., as  $n \to \infty$ .

If f is  $\Lambda$ -integrable, then we put

$$\{ \int_{A} f d\Lambda = P - \lim_{n \to \infty} \int_{A} f_{n} d\Lambda,$$

where  $\{f_n\}$  satisfies (i) and (ii).

Now, combining Theorems 3.2, 3.3, 4.1 we get the following:

**Theorem 4.2** Given  $\alpha > 0$ , let  $\Lambda(A)$ ,  $A \in \mathcal{S}$  be a  $\alpha$ -s.d.r.m. Then, the characteristic function of  $\Lambda(A)$  is of the unique form (20) where

$$\begin{cases}
K(t,s) = ita(s) - 1/2t^2\sigma^2(s) \\
+ \int_A (e^{itx} - 1 - it\tau(x))\rho(s, dx)
\end{cases}$$
(32)

with

$$a(s) = \frac{dv_0}{d\lambda}(s) \tag{33}$$

and

$$\sigma^2(s) = \frac{dv_1}{d\lambda}(s) \tag{34}$$

and  $\rho$  is given by Lemma 2.3 in [3]. Moreover, we have

$$|a(s)| + \int_{\mathbb{R}} \min\{1, x^2\} \rho(s, dx) = 1 \quad a.e.[\lambda].$$

**Proof.** By virtue of (13) it follows that for any  $A \in S$  and  $t \in \mathbb{R}$   $\Lambda(A)$  has the representation

$$-logEexp(it\Lambda(A)) = it\nu_0(A) + \frac{1}{2}t^2\nu_1(A) - \int_{-\infty}^{\infty} v_{\alpha}(x) \left(\int_0^{\infty} k(e^{-u}x, t)u^{\alpha - 1}du\right) m(A, dx)$$

$$(35)$$

which, by a similar argument of Proposition 2.1 in [3], implies that there exists a unique finite measure  $\nu$  on  $\sigma(S) \times \mathcal{B}(\mathcal{R})$  such that

$$\nu(A \times B) = m(A, B), \text{ for any } A \in \mathcal{S}, B \in \mathcal{B}(\mathcal{R}).$$

Moreover, for every  $A \in \sigma(S)$  we have  $\nu(A, \{0\}) = 0$ .

Now,we are in the position to present the following theorem whose proof is a simple combination of Theorem 6 and the Komogorov extension theorem and Threorem 5.2 in [3].

**Theorem 4.3** Given  $0 < \alpha \le \infty$  let  $\{X_t : t \in T\}$  be an  $\alpha - s.d.$  stochastic process defined on a probability space  $(\Omega', \mathcal{P}')$ . Then there exists an  $\alpha - s.d.r.m.$ , say  $\Lambda$ , defined on the probability space  $(\Omega, \mathcal{P})$  such that

$$(\Omega = \Omega' \times [0, 1], \mathcal{P} = \mathcal{P}' \times Leb),$$

Leb being the Lebesgue measure on [0,1] and

$${X_t: t \in T} = {\int_{\mathbb{S}} f_t(s) d\Lambda(s): t \in T} \quad a.s.\mathcal{P},$$

where  $\{f_t(s): t \in T, s \in S\}$  are some measurable functions on  $\mathbb{S}$ .

## 5 An Application in Option Pricing

If X is Lévy - stable random variable with index  $0 < \alpha < 1$ , then it does not have any integer moment, and for the case  $1 < \alpha < 2$  only the first integer moment exists. Therefore, to overcome this dificulties, following Cartea and Howinson [1], we introduce the following Damped - Lévy - mixed - stable process which will lead to a mathematical model for our purpose of option pricing. Suppose that  $X_i(t)$ , i = 1, 2 are independent Lévy -stable processes with indexes  $0 < \alpha_1 < \alpha_2 < 2$ , respectively such that the logarithm of the characteristic function of  $X_i(1)$  is given by

$$\psi_j(u) = \int_{-\infty}^{+\infty} (e^{iux} - 1 - iu\tau_{\alpha_j}(x))W_j(x)dx, \quad j = 1, 2.$$
 (36)

where

$$W_j(x) = \begin{cases} C_q |x|^{-1-\alpha_j} & \text{for } x < 0\\ C_p x^{-1-\alpha_j} & \text{for } x > 0 \end{cases}$$

and

$$\tau_{\alpha_j}(x) = \begin{cases} x & \text{for } \alpha_j > 1\\ sinx & \text{for } \alpha_j = 1\\ 0 & \text{for } \alpha_j < 1. \end{cases}$$

Here  $C_p, C_q > 0$  are scale constants,  $p, q \ge 0$  and p + q = 1. Following Cartea and Howinson [1] the exponential cut-off  $e^{-\lambda|x|}$  is introduced to obtain the Damped Lévy measures

$$W_j^{\lambda}(x) = \begin{cases} C_q |x|^{-1-\alpha} e^{-\lambda|x|}, & \text{for } x < 0 \\ C_p x^{-1-\alpha} e^{-\lambda|x|}, & \text{for } x > 0 \end{cases}$$
 (37)

Let  $W_j^{\lambda}$ , j=1,2, denote the Damped Lévy measures corresponding to Lévy processes  $X_j^{\lambda}(t)$ , j=1,2 with

$$\phi_j(u) = \int_{-\infty}^{+\infty} \left( e^{iux} - 1 - iu\tau_{\alpha_j}(x) \right) e^{-\lambda|x|} W_j(dx)$$
(38)

Putting, for  $t \ge 0, X(t) = X_1(t) + X_2(t)$  we get a Lévy process X(t) which is also a mixed - stable - Lévy

process with  $\Phi(u) = \Phi_1(u) + \Phi_2(u)$  where  $\Phi_j(u), j = 1, 2$  are given by (38). Putting

$$W_j^{\lambda}(x) = \begin{cases} C_q |x|^{-1-\alpha_j} e^{-\lambda|x|}, & \text{for } x < 0 \\ C_p x^{-1-\alpha_j} e^{-\lambda|x|}, & \text{for } x > 0 \end{cases}, j = 1, 2$$
 (39)

and taking into account (38) we infer that the logarithm of the characteristic function, denoted by  $\phi^{\lambda}(u)$ , for a Damped-Levy's process  $\{X^{\lambda}(t)\}$  is of the form

$$\phi^{\lambda}(u) = \phi_1^{\lambda}(u) + \phi_2^{\lambda}(u)$$

where

$$\phi_j^{(\lambda)}(u) = \int_{-\infty}^{+\infty} (e^{iux} - 1 - iu\tau_{\alpha_j}(x))e^{-\lambda|x|}W_j^{(\lambda)}(dx), \ j = 1, 2$$
 (40)

which implies that the Damped Lévy process  $X^{\lambda}(t):=X_1^{\lambda}(t)+X_2^{\lambda}(t)$  has the following property:

- (i)  $\{X^{\lambda}(t)\}$  is a Lévy process.
- (ii) It is not a stable process.
- (iii)  $\lim_{\lambda \to 0} X^{\lambda}(t) = X(t)$  (in distribution and in probability).
- (iv) The process  $\{X^{\lambda}(t)\}$  has finite moments of all orders. Moreover, its exponential moments exist.

**Definition 5.1 (Power - Jump Process)** For k = 2, 3, ... we define, following M. Corunera, D. Nualart and W.Schouten, the following Power - Jump process.

$$Z_t^{(k)} = \sum_{0 < s \le t} [\Delta X^{\lambda}(t)]^k, \ k \ge 2$$
 (41)

where

$$\Delta X \lambda(t) = X^{\lambda}(t) - X^{\lambda}(t^{-})$$

and, for convenience, we put

$$Z_t^{(k)} = X^{\lambda}(t).$$

Each process  $Z_t^{(k)}$ ,  $k \geq 2$  is called the k-th power - jump process.

We have

$$E[Z_t^{(1)}] = E[X^{\lambda}(t)] = ta^{\lambda}$$

$$E[Z_t^{(k)}] = E\left[\sum_{0 < s \le t} [\Delta X^{\lambda}(t)]^k\right] = t \int_{-\infty}^{+\infty} x^k \gamma(dx) = m_k t$$

#### Definition 5.2 Put

$$Y_t^{(k)} = Z_t^{(k)} - m_k t, \quad k = 1, 2, 3, \dots$$

Then,  $\{Y_t^{(k)}\}$  is a normal martingale and is called Teugels martingale of order k

By taking a suitable linear combination of the  $\{Y_t^{(k)}\}$ , one get pairwise strongly orthonormal martingale, say  $\{T^{(k)}\}$ , which correspond to the procedure of the orthonormal lization of the polynomial 1, x,  $x^2$  with respect to to measure

$$\mu(dx) = x^2 \gamma(dx) + c^2 \delta(dx)$$

The resulting process

$$\{T_t^{(k)}, t > 0\}$$

are called the orthonormalization k -th jump process.

**Theorem 5.1 (Nualart - Schouten)** Let the  $\sigma$ -fields  $(\mathcal{F}_t)$  be generated by  $\{K_t\}$  and  $L_T^2(T>0)$  be the space of all square integrable and  $(\mathcal{F}_t)$  - measurable  $\Phi=\Phi_t$ ,  $t\in[0,T]$  such that

$$\|\Phi\|^2 = E[\int_0^T |\Phi_t^2| dt] < \infty$$

Then each square integrable random variable  $F \in L^2_T$  has the representation

$$F = E[F] + \sum_{k=1}^{\infty} V_s^{(k)}$$

where  $\{V_s^{(k)}\}$  is a predictable process in  $M_T^2(l^2)$ , the space of predictable  $l^2$ -valued processes

Suppose that we work under the framework of the market with the k-th power - jump assets. Note that the value of the contingent claim at time t is given by

$$F(S_t) = exp(-r(T-t))E_Q[X|\mathcal{F}_t].$$

We call  $F(t, S_t)$  the price function. Note that we have choosen an equivalent martingale measure Q under which all discounted assets price processes are martingales.

#### The Black - Scholes Formula under Lévy Mixed Stable Shocks

Let  $0 < \alpha_1 < \alpha_2 < 2$ . Suppose that our stock price  $X(t) = X_1(t) + X_2(t)$  satisfies the condition that  $X_j(t), j = 1, 2$  are  $\alpha_k$  - stable Lévy processes under measure Q, respectively.

Let

$$X_{t+\Delta t} = X_t e^{\mu \Delta t + \sigma \Phi}$$

where  $\Phi$  is a parameter for Damped - $(\alpha_1, \alpha_2)$  mixed - stable - Lévy process  $X^{\lambda}(t)$ . Then as  $\Delta t \to 0$  the "Damped Black - Scholes" partial differential equation (PDE) satisfies  $\frac{\delta \hat{V}}{\delta t} = \psi(-\bar{\sigma}\xi) + i\xi[\psi(-i\bar{\sigma}) + \bar{D}_t - r(1-i\xi)\hat{V}]$  when  $\psi$  is the logarithm of the characteristic function of the Damped - Lévy - mixed - stable process and

$$V(S,t) \stackrel{\text{def}}{=} E_Q[e^{-r(T-t)}\Pi(S,T)]$$

where  $\Pi(S,T)$  stands for payoff at time t=T.

Notice that we can write

$$V(S,t) = e^{-r\delta t} E_O[e^{-r(T-t-\Delta t)}\Pi(S,T)] = e^{-r\delta t} E_O[V(St+\Delta t,t) + \delta t]$$

which is equivalent to the equation

$$rV(S,t)\Delta t = E_Q[\Delta V(S,t)]$$

where

$$\Delta(S,t) \stackrel{\text{def}}{:=} V(V + \Delta S, t + \Delta t) - V(S,t).$$

Finally, one can solve the Black -Scholes equation by a similar method as for the case of Lévy - stable process [1].

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