# Concept of Sub-Independence and Characterizations of 2SQLindley and 2DQLindley Distributions 

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#### Abstract

Amer et al. [1] considered the distributions of the sum and the difference of two independent and identically distributed random variables with the common Quasi Lindley distribution. They derived, very nicely, the above mentioned distributions and provided certain important mathematical and statistical properties as well as simulations and applications of the new distributions. Wang and Ma [2] considered the sum of the gamma random variables under the assumption of independence of the summands and presented very interesting results. In this short note, we like to show that the assumption of "independence" can be replaced with a much weaker assumption of "sub-independence" in both papers. Then we present certain characterizations of the distributions derived by Amer et al. [1], called 2SQLindley and 2DQLindley distributions.


Keywords: quasi lindley distribution, gamma distribution, sub-independence, identically distributed random variables, characterizations of distributions

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## 1. Introduction

As we have done in a couple of our previous papers, to make this short note self contained, we will copy some parts of our work [3] here.

We may in some occasions have asked ourselves if there is a concept between "uncorrelatedness" and "independence" of two random variables. It seems that the concept of "sub-independence" is the one: it is much stronger than uncorrelatedness and much weaker than independence. The notion of sub-independence seems important in the sense that under usual assumptions, Khintchine's Law of Large Numbers and LindebergLevy's Central Limit Theorem as well as other important theorems in probability and statistics hold for a sequence of s.i. (sub-independent) random variables. While subindependence can be substituted for independence in many cases, it is difficult, in general, to find conditions under which the former implies the latter. Even in the case of two discrete identically distributed rv's (random variables) X and Y , the joint distribution can assume many forms consistent with sub-independence.

Limit theorems as well as other well-known results in probability and statistics are often based on the distribution of the sums of independent (and often identically distributed) random variables rather than the joint distribution of the summands. Therefore, the full force of independence of the summands will not be
required. In other words, it is the convolution of the marginal distributions which is needed, rather than the joint distribution of the summands which, in the case of independence, is the product of the marginal distributions. The concept of sub-independence, which is weaker than that of independence, is shown to be sufficient to yield the conclusions of these theorems and results. This is precisely the reason for the statement: "why assume independence when you can get by with sub-independence".

The concept of sub-independence can help to provide solution for some modeling problems where the variable of interest is the sum of a few components. Examples include household income, the total profit of major firms in an industry, and a regression model $Y=g(X)+\varepsilon$ where $g(X)$ and $\varepsilon$ are uncorrelated, however, they may not be independent. For example, in Bazargan et al. [4], the return value of significant wave height $(Y)$ is modeled by the sum of a cyclic function of random delay $D, \widehat{g}(\mathrm{D})$, and a residual term $\hat{\varepsilon}$. They found that the two components are at least uncorrelated but not independent and used subindependence to compute the distribution of the return value.

Let $X$ and $Y$ be two $r v^{\prime} s$ (random variables) with joint and marginal $c d f^{\prime} s$ (cumulative distribution functions) $F_{X, Y}, F_{X}$ and $F_{Y}$ respectively. Then $X$ and $Y$ are said to be independent if and only if

$$
F_{X, Y}(x, y)=F_{X}(x) F_{Y}(y), \text { for all }(x, y) \in \mathbb{R}^{2},(1.1)
$$

or equivalently, if and only if

$$
\begin{equation*}
\varphi_{X, Y}(s, t)=\varphi_{X}(s) \varphi_{Y}(t), \text { for } \operatorname{all}(s, t) \in \mathbb{R}^{2} \tag{1.2}
\end{equation*}
$$

where $\varphi_{X, Y}(s, t), \varphi_{X}(s)$ and $\varphi_{Y}(t)$, respectively, are the corresponding joint and marginal $c f^{\prime} s$ (characteristic functions). Note that (1.1) and (1.2) are also equivalent to

$$
\begin{equation*}
P(X \in A \text { and } Y \in B)=P(X \in A) P(Y \in B) \tag{1.3}
\end{equation*}
$$

for all Borel sets $A, B$.
The concept of sub-independence, as far as we have gathered, was formally introduced by Durairajan (1979) and developed by Hamedani in the past 40 years, stated as follows: The $r v^{\prime} s X$ and $Y$ with $c d f^{\prime} s F_{X}$ and $F_{Y}$ are s.i. (sub-independent) if the $c d f$ of $X+Y$ is given by

$$
\begin{align*}
& F_{X+Y}(z)=\left(F_{X} * F_{Y}\right)(z) \\
& =\int_{\mathbb{R}} F_{X}(z-y) d F_{Y}(y), z \in \mathbb{R}, \tag{1.4}
\end{align*}
$$

or equivalently if and only if

$$
\varphi_{X+Y}(t)=\varphi_{X, Y}(t, t)=\varphi_{X}(t) \varphi_{Y}(t), \text { for all } t \in \mathbb{R}
$$

The drawback of the concept of sub-independence in comparison with that of independence has been that the former does not have an equivalent definition in the sense of (1.3) which some believe, to be the natural definition of independence. We found such a definition which is stated below. We shall give the definition for the continuous case (Definition 1.1).

We observe that the half-plane $H=\{(x, y): x+y<0\}$ can be expressed as a countable disjoint union of rectangles:

$$
H=\bigcup_{i=1}^{\infty} E_{i} \times F_{i}
$$

where $E_{i}$ and $F_{i}$ are intervals. Now, let $(X, Y): \Omega \rightarrow \mathbb{R}^{2}$ be a continuous random vector and for $\in \mathbb{R}$, let

$$
A_{c}=\{\omega \in \Omega: X(\omega)+Y(\omega)<c\}
$$

and

$$
\begin{aligned}
& A_{i}^{(c)}=\left\{\omega \in \Omega: X(\omega)-\frac{c}{2} \in E_{i}\right\}, \\
& B_{i}^{(c)}=\left\{\omega \in \Omega: Y(\omega)-\frac{c}{2} \in F_{i}\right\} .
\end{aligned}
$$

Definition 1.1. The continuous $r v^{\prime} s X$ and $Y$ are s.i. if for every $c \in \mathbb{R}$

$$
\begin{equation*}
P\left(A_{c}\right)=\sum_{i=1}^{\infty} P\left(A_{i}^{(c)}\right) P\left(B_{i}^{(c)}\right) \tag{1.6}
\end{equation*}
$$

To see that (1.6) is equivalent to (1.4), observe that (LHS of (1.6))

$$
\begin{equation*}
P\left(A_{c}\right)=P(X+Y<c)=P\left((X, Y) \in H_{c}\right) \tag{1.7}
\end{equation*}
$$

where $H_{c}=\{(x, y): x+y<c\}$. Now, if $X$ and $Y$ are s. $i$. then

$$
P\left(A_{c}\right)=\left(P_{X} \times P_{Y}\right)\left(H_{c}\right)
$$

where $P_{X}, P_{Y}$ are probability measures on $\mathbb{R}$ defined by

$$
P_{X}(B)=P(X \in B) \text { and } P_{Y}(B)=P(Y \in B) \text {, }
$$

and $P_{X} \times P_{Y}$ is the product measure.
We also observe that (RHS of (1.6))

$$
\begin{aligned}
& \sum_{i=1}^{\infty} P\left(A_{i}^{(c)}\right) P\left(B_{i}^{(c)}\right) \\
& =\sum_{i=1}^{\infty} P\left(X-\frac{c}{2} \in E_{i}\right) P\left(Y-\frac{c}{2} \in F_{i}\right) \\
& =\sum_{i=1}^{\infty} P\left(X \in E_{i}+\frac{c}{2}\right) P\left(Y \in F_{i}+\frac{c}{2}\right) \\
& =\sum_{i=1}^{\infty} P_{X} \times P_{Y}\left(E_{i}+\frac{c}{2}\right) \times\left(F_{i}+\frac{c}{2}\right) .
\end{aligned}
$$

Now, (1.7) and (1.8) will be equal if $H_{c}=$ $\cup_{i=1}^{\infty}\left\{\left(E_{i}+\frac{c}{2}\right) \times\left(F_{i}+\frac{c}{2}\right)\right\}$, which is true since the points in $H_{c}$ are obtained by shifting each point in $H$ over to the right by $\frac{c}{2}$ units and then up by $\frac{c}{2}$ units.

If $X$ and $Y$ are s.i., then unlike independence, $X$ and $\alpha Y$ are not necessarily s.i. for any real $\alpha \neq 1$. This demonstrates how weak is the concept of subindependence in comparison with that of independence. Please observe the following simple example.
Example 1.1. Let $X$ and $Y$ have the joint $c f$ given by

$$
\begin{aligned}
& \varphi_{X, Y}\left(t_{1}, t_{2}\right) \\
& =\exp \left\{-\left(t_{1}^{2}+t_{2}^{2}\right) / 2\right\} \\
& \times\left[1+\beta t_{1} t_{2}\left(t_{1}-t_{2}\right)^{2} \exp \left\{\left(t_{1}^{2}+t_{2}^{2}\right) / 4\right\}\right] \\
& \left(t_{1}, t_{2}\right) \in \mathbb{R}^{2}
\end{aligned}
$$

where $\beta$ is an appropriate constant. (The characteristic function is the Fourier transform of $p d f$ (probability density function), so the corresponding joint $p d f$ is given by

$$
\begin{aligned}
& f(x, y) \\
& =\frac{1}{2 \pi} \exp \left\{-\left(x^{2}+y^{2}\right) / 2\right\} \\
& \times\left[1-16 \beta p(x, y) \exp \left\{-\left(x^{2}+y^{2}\right) / 2\right\}\right] \\
& (x, y) \in \mathbb{R}^{2}
\end{aligned}
$$

where $p(x, y)=\left\{6 x y-2 x^{2}-2 y^{2}+4 x^{2} y^{2}-2 x^{3} y-\right.$ $2 x y 3+1$ ).

Then $X$ and $Y$ are s.i. standard normal $r v^{\prime} s$, and hence $X+Y$ is normal with mean 0 and variance 2 , but $X$ and $-Y$ are not s.i. and consequently $X-Y$ does not have a normal distribution.

The concept of sub-independence defined above can be extended to $n(>2) r v ' s$ as follows.
Definition 1.2. The $r v^{\prime} s X_{1}, X_{2}, \ldots, X_{n}$ are $s . i$. if for each subset $\left\{X_{\alpha_{1}}, X_{\alpha_{2}}, \ldots, X_{\alpha_{r}}\right\}$ of $\left\{X_{1}, X_{2}, \ldots, X_{n}\right\}$

$$
\varphi_{X_{\alpha_{1}}}, \ldots, X_{\alpha_{r}}(t, \ldots, t)=\Pi_{i=1}^{r} \varphi_{X_{\alpha_{i}}}(t), \text { for all } t \in \mathbb{R}
$$

## 2. Remarks

i) If the $r v^{\prime} s X$ and $Y$ are s.i.i.d. with common Quasi Lindley distribution with parameters $\alpha, \theta$, the characteristic function of $X+Y$ is

$$
\varphi_{X+Y}(t)=\frac{\theta^{4}[\alpha(\theta-i t)+\theta]^{4}}{(1+\alpha)^{4}(\theta+t)^{8}}, t \in \mathbb{R}
$$

The $c f$ of $X$ is

$$
\begin{aligned}
& \varphi_{X}(t)=\int_{0}^{\infty} e^{i t x} \frac{\theta(\alpha+\theta x)}{\alpha+1} e^{-\theta x} d x \\
& =\frac{\theta^{2}[\alpha(\theta-i t)+\theta]^{2}}{(1+\alpha)^{2}(\theta+t)^{4}}, t \in \mathbb{R},
\end{aligned}
$$

and since $X$ and $Y$ are s.i., we have

$$
\begin{aligned}
& \varphi_{X+Y}(t)=\varphi_{X}(t) \varphi_{Y}(t)=\left\{\frac{\theta^{2}[\alpha(\theta-i t)+\theta]^{2}}{(1+\alpha)^{2}(\theta+t)^{4}}\right\}^{2} \\
& =\frac{\theta^{4}[\alpha(\theta-i t)+\theta]^{4}}{(1+\alpha)^{4}(\theta+t)^{8}}, t \in \mathbb{R}
\end{aligned}
$$

ii) If the $r v^{\prime} s X$ and $Y$ are i.d. with common Quasi Lindley distribution with parameters $\alpha, \theta$, and if $X$ and $-Y$ are s. i., the characteristic function of $X-Y$ is

$$
\varphi_{X-Y}(t)=\frac{\theta^{4}\left[\theta^{4}(\alpha+1)^{4}+\alpha^{4} t^{2}\right]}{(1+\alpha)^{4}\left(\theta^{2}+t^{2}\right)^{4}}, t \in \mathbb{R}
$$

The $c f$ of $X-Y$, under the assumption of s.i. of $X$ and $-Y$, is

$$
\begin{aligned}
& \varphi_{X-Y}(t) \\
& =\varphi_{X}(t) \varphi_{-Y}(t)=\varphi_{X}(t) \varphi_{Y}(-t) \\
& =\left\{\frac{\theta^{2}\left[\theta^{2}(\alpha+1)^{2}-\alpha^{2} i t\right]}{(1+\alpha)^{2}\left(\theta^{2}+t^{2}\right)^{2}}\right\}\left\{\frac{\theta^{2}\left[\theta^{2}(\alpha+1)^{2}+\alpha^{2} i t\right]}{(1+\alpha)^{2}\left(\theta^{2}+t^{2}\right)^{2}}\right\} \\
& =\frac{\theta^{4}\left[\theta^{4}(\alpha+1)^{4}+\alpha^{4} t^{2}\right]}{(1+\alpha)^{4}\left(\theta^{2}+t^{2}\right)^{4}}, t \in \mathbb{R} .
\end{aligned}
$$

iii) In view of $\mathbf{i}$ ) and $\mathbf{i i}$ ), the assumption of "independence" in Amer et al. paper can be replaced with that of "sub-independence".
iv) Equation (2.3) of Wang and Ma holds for s.i. gamma r.v. s.
v) In Theorem 3.2 of Wang and Ma, the distribution of $\sum_{i=1}^{n} X_{i}$ is a gamma distribution with parameters $\sum_{i=1}^{n} \alpha_{i}$, $\beta$ under the assumption that $X_{i}^{\prime} s$ are $s$. $i$. and $\beta_{i}^{\prime} s$ are equal to $\beta$.
vi) In Theorem 5.1 of Wang and Ma, the distribution of $\sum_{i=1}^{m} X_{i}$ is a Chi-square distribution with parameter $\sum_{i=1}^{m} n_{i}$ under the assumption that $X_{i}^{\prime} s$ are s.i..
vii) For a detailed treatment of the concept of sub-independence, we refer the interested reader to Hamedani [3].

## 3. Characterizations of the 2S-Lindley and 2D-Lindley Distributions

Amer et al. [1] introduced the distributions of the sum and differences of two i.i.d.(now, s.i.i.d.) Quasi Lindley random variables with parameters $\alpha>-1, \theta>0$ (called 2SQLindley and 2DQLindley) with their respective $p d f s$ given by

$$
\begin{align*}
& f_{\text {2SQLindley }}(x) \\
& =\frac{\theta^{2} x}{(\alpha+1)^{2}}\left(\alpha^{2}+\theta \alpha x+\frac{\theta^{2} x^{2}}{6}\right) e^{-\theta x}, x>0 \tag{3.1}
\end{align*}
$$

and

$$
\begin{align*}
& f_{2 D Q L i n d l e y}(x) \\
& =\frac{\theta e^{-\theta|x|}}{4(\alpha+1)^{2}}\left[\theta(2 \alpha+1)|x|+2 \alpha^{2}+2 \alpha+1\right], x \in \mathbb{R} \tag{3.2}
\end{align*}
$$

Following our [3] work, to understand the behavior of the data obtained through a given process, we need to be able to describe this behavior via its approximate probability law. This, however, requires to establish conditions which govern the required probability law. In other words we need to have certain conditions under which we may be able to recover the probability law of the data. So, characterization of a distribution is important in applied sciences, where an investigator is vitally interested to find out if their model follows the selected distribution. Therefore, the investigator relies on conditions under which their model would follow a specified distribution. A probability distribution can be characterized in different directions one of which is based on the truncated moments. This type of characterization initiated by Galambos and Kotz [5] and followed by other authors such as Kotz and Shanbhag [6], Glänzel et al. [7], Glänzel [8], Glänzel and Hamedani [9] and Kim and Jeon [10], to name a few. For example, Kim and Jeon [10] proposed a credibility theory based on the truncation of the loss data to estimate conditional mean loss for a given risk function. It should also be mentioned that characterization results are mathematically challenging and elegant. In this section, we present characterizations 2S-Lindley and 2D-Lindley distributions based on the conditional expectation (truncated moments) of certain function of the random variable.

We will employ Theorem 1 of Glänzel [8] given in the Appendix A. As shown in Glänzel [11], this characterization is stable in the sense of weak convergence. Proposition 3.1. Let $X$ be a continuous random variable and let $q_{1}(x)=x^{-1}\left(\alpha^{2}+\theta \alpha x+\frac{\theta^{2} x^{2}}{6}\right)^{-1}$ and $q_{2}(x)=$
$q_{1}(x) e^{-\theta x}$ for $x>0$. Then $X$ has $p d f$ (3.1) if and only if the function $\xi$ defined in Theorem 1 is of the form

$$
\xi(x)=\frac{1}{2} e^{-\theta x}, x>0
$$

Proof. If $X$ has $p d f$ (3.1), then

$$
\left(1-F_{2 S-L}(x)\right) E\left[q_{1}(X) \mid X \geq x\right]=\frac{\theta}{(\alpha+1)^{2}} e^{-\theta x}, x>0
$$

and

$$
\begin{aligned}
& \left(1-F_{2 S-L}(x)\right) E\left[q_{2}(X) \mid X \geq x\right] \\
& =\frac{\theta}{2(\alpha+1)^{2}} e^{-2 \theta x}, x>0,
\end{aligned}
$$

and hence

$$
\xi(x)=\frac{1}{2} e^{-\theta x}, x>0
$$

We also have

$$
\xi(x) q_{1}(x)-q_{2}(x)=-\frac{1}{2} q_{1}(x) e^{-\theta x}<0, \text { for } x>0
$$

Conversely, if $\xi$ is of the above form, then

$$
s^{\prime}(x)=\frac{\xi^{\prime}(x) q_{1}(x)}{\xi(x) q_{1}(x)-q_{2}(x)}=\theta, x>0 .
$$

and

$$
s(x)=\theta x .
$$

Now, according to Theorem $1, X$ has density (3.1).
Corollary 3.1. Suppose $X$ is a continuous random variable. Let $q_{1}(x)$ be as in Proposition 3.1. Then $X$ has density (3.1) if and only if there exist functions $q_{2}$ and $\xi$ defined in Theorem 1 for which the following first order differential equation holds

$$
\frac{\xi^{\prime}(x) q_{1}(x)}{\xi(x) q_{1}(x)-q_{2}(x)}=\theta, x>0
$$

Corollary 3.2. The differential equation in Corollary 3.1 has the following general solution

$$
\xi(x)=e^{\theta x}\left[-\int \theta e^{-\theta x}\left(q_{1}(x)\right)^{-1} q_{2}(x)+D\right]
$$

where $D$ is a constant.
Proof. If $X$ has pdf (3.1), then clearly the differential equation holds. Now, if the differential equation holds, then

$$
\xi^{\prime}(x)=\xi(x) \theta-\theta\left(q_{1}(x)\right)^{-1} q_{2}(x)
$$

or

$$
\xi^{\prime}(x) e^{-\theta x}-\xi(x) \theta e^{-\theta x}=-\theta e^{-\theta x}\left(q_{1}(x)\right)^{-1} q_{2}(x)
$$

or

$$
\frac{d}{d x}\left\{e^{-\theta x} \xi(x)\right\}=-\theta e^{-\theta x}\left(q_{1}(x)\right)^{-1} q_{2}(x)
$$

from which we arrive at

$$
\xi(x)=e^{\theta x}\left[-\int \theta e^{-\theta x}\left(q_{1}(x)\right)^{-1} q_{2}(x)+D\right],
$$

A set of functions satisfying the above differential equation is given in Proposition 3.1 with $D=0$. Clearly, there are other triplets $\left(q_{1}, q_{2}, \xi\right)$ satisfying the conditions of Theorem 1.
Remark 3.1. Similar results can be stated for the 2DQLindley distribution as well.

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## Appendix A

Theorem 1. Let $(\Omega, \mathcal{F}, \mathrm{P})$ be a given probability space and let $H=[a, b]$ be an interval for some $a<b$ ( $a=-\infty, b=$ $\infty$ might as well be allowed). Let $X: \Omega \rightarrow H$ be a continuous random variable with the distribution function $F$ and let $q_{1}$ and $q_{2}$ be two real functions defined on $H$ such that

$$
\mathrm{E}\left[q_{2}(X) \mid X \geq x\right]=\mathrm{E}\left[q_{1}(X) \mid X \geq x\right] \xi(x), x \in H
$$

is defined with some real function $\xi$. Assume that $q_{1}, q_{2} \in C^{1}(H), \xi \in C^{2}(H)$ and $F$ is twice continuously differentiable and strictly monotone function on the set $H$. Finally, assume that the equation $\xi q_{1}=q_{2}$ has no real solution in the interior of $H$. Then $F$ is uniquely determined by the functions $q_{1}, q_{2}$ and $\xi$, particularly

$$
F(x)=\int_{a}^{x} C\left|\frac{\xi^{\prime}(u)}{\xi(u) q_{1}(u)-q_{2}(u)}\right| \exp (-s(u)) d u
$$

where the function $s$ is a solution of the differential equation $s^{\prime}=\frac{\xi^{\prime} q_{1}}{\xi q_{1}-q_{2}}$ and $C$ is the normalization constant, such that $\int_{H} d F=1$.
Note: The goal is to have the function $\xi(x)$ as simple as possible.
We like to mention that this kind of characterization based on the ratio of truncated moments is stable in the sense of weak convergence (see, [11]), in particular, let us assume that there is a sequence $\left\{X_{n}\right\}$ of random variables with distribution functions $\left\{F_{n}\right\}$ such that the functions $q_{1 n}, q_{2 n}$ and $\xi_{n}(n \in \mathbb{N})$ satisfy the conditions of Theorem 1 and let $q_{1 n} \rightarrow q_{1}, q_{2 n} \rightarrow q_{2}$ for some continuously differentiable real functions $q_{1}$ and $q_{2}$. Let, finally, $X$ be a random variable with distribution $F$. Under the condition that $q_{1 n}(X)$ and $q_{2 n}(X)$ are uniformly integrable and the family $\left\{F_{n}\right\}$ is relatively compact, the sequence $X_{n}$ converges to $X$ in distribution if and only if $\xi_{n}$ converges to $\xi$, where

$$
\xi(x)=\frac{E\left[q_{2}(X) \mid X \geq x\right]}{E\left[q_{1}(X) \mid X \geq x\right]}
$$

This stability theorem makes sure that the convergence of distribution functions is reflected by corresponding convergence of the functions $q_{1}, q_{2}$ and $\xi$, respectively. It guarantees, for instance, the 'convergence' of characterization of the Wald distribution to that of the Lévy-Smirnov distribution if $\alpha \rightarrow \infty$.

A further consequence of the stability property of Theorem 1 is the application of this theorem to special tasks in statistical practice such as the estimation of the parameters of discrete distributions. For such purpose, the functions $q_{1}, q_{2}$ and, specially, $\xi$ should be as simple as possible. Since the function triplet is not uniquely determined it is often possible to choose $\xi$ as a linear function. Therefore, it is worth analyzing some special cases which helps to find new characterizations reflecting the relationship between individual continuous univariate distributions and appropriate in other areas of statistics.

In some cases, one can take $q_{1}(x) \equiv 1$, which reduces the condition of Theorem 1 to $\mathrm{E}\left[q_{2}(X) \mid X \geq x\right]=\xi(x), x \in H$. We, however, believe that employing three functions $q_{1}, q_{2}$ and $\xi$ will enhance the domain of applicability of Theorem 1 .
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