

pfsl12.tex

Lecture 12. 6.11.2012

Take expectations: as $Ey = \mu$, $Eg(y) \sim g(\mu)$. So

$$g(y) - g(\mu) \sim g(y) - Eg(y) \sim g'(\mu)(y - \mu).$$

Square both sides:

$$[g(y) - g(\mu)]^2 \sim [g'(\mu)]^2(y - \mu)^2.$$

Take expectations: as $Ey = \mu$ and $Eg(y) \sim g(\mu)$, this says

$$\text{var}(g(y)) \sim [g'(\mu)]^2 \text{var}(y).$$

Regression. So if

$$E(y_i|x_i) = \mu_i, \quad \text{var}(y_i|x_i) = \sigma_i^2,$$

we use EDA to try to find some link between the means μ_i and the variances σ_i^2 . Suppose we try $\sigma_i^2 = H(\mu_i)$, or

$$\sigma^2 = H(\mu).$$

Then by above,

$$\text{var}(g(y)) \sim [g'(\mu)]^2 \sigma^2 = [g'(\mu)]^2 H(\mu).$$

We want *constant variance*, c^2 say. So we want

$$[g'(\mu)]^2 H(\mu) = c^2, \quad g'(\mu) = \frac{c}{\sqrt{H(\mu)}}, \quad g(y) = c \int \frac{dy}{\sqrt{H(y)}}.$$

Note. The idea of variance-stabilising transformations (like so much else in Statistics) goes back to Fisher (R. A. (Sir Ronald) FISHER (1890-1962)). He found the density of the sample correlation coefficient r^2 in the bivariate normal distribution – a complicated function involving the population correlation coefficient ρ^2 , simplifying somewhat in the case $\rho = 0$ (see e.g. [KS1], S16.27, 28). But Fisher's z -transformation of 1921 ([KS1], S16.33)

$$r = \tanh z, \quad z = \frac{1}{2} \log\left(\frac{1+r}{1-r}\right), \quad \rho = \tanh \zeta, \quad \zeta = \frac{1}{2} \log\left(\frac{1+\rho}{1-\rho}\right)$$

gives z approximately normal, with variance almost independent of ρ :

$$z \sim N(0, 1/(n-1)).$$

4. Infinite divisibility; self-decomposability; stability: $I \supset SD \supset S$

In the CLT, the limit distribution is normal, $N(0, 1)$, CF $\exp\{-\frac{1}{2}t^2\}$. Note that for each $n = 1, 2, \dots$,

$$\exp\{-\frac{1}{2}t^2\} = [\exp\{-\frac{1}{2}t^2/n\}]^n$$

expresses the CF of the limit law $N(0, 1)$ as the n th power of the CF of another probability law, $N(0, 1/n)$. So $N(0, 1)$ is the n th convolution of $N(0, 1/n)$. We think of this as ‘splitting $N(0, 1)$ up into n independent parts’: $N(0, 1)$ is n times ‘divisible’. We can do this for each n , so $N(0, 1)$ is ‘infinitely divisible’.

Similarly for X Poisson $P(\lambda)$: the CF is

$$E[e^{itX}] = \sum_{n=0}^{\infty} e^{-\lambda} \lambda^n \cdot e^{itn}/n! = \exp\{-\lambda(1 - e^{it})\} = [\exp\{-(\lambda/n)(1 - e^{it})\}]^n,$$

so $P(\lambda)$ is the n -fold convolution of $P(\lambda/n)$, for each n . So the Poisson distributions are infinitely divisible (id).

We can extend this to the *compound Poisson distribution* $CP(\lambda, F)$, which is very important in the actuarial/insurance industry. Suppose that the number of claims is Poisson $P(\lambda)$, and that the claim sizes are iid, with distribution F and CF ϕ . Then conditional on the number of claims being n , the total claimed in the n th convolution F^{*n} , and the CF is ϕ^n . So the total X claimed has CF

$$E[e^{itX}] = \sum_{n=0}^{\infty} e^{-\lambda} \lambda^n \cdot \phi(t)^n/n! = \exp\{-\lambda(1 - \phi(t))\} = [\exp\{-(\lambda/n)(1 - \phi(t))\}]^n.$$

So $CP(\lambda, F)$ is the n -fold convolution of $CP(\lambda/n, F)$ for each n , so is id. But this holds much more generally.

Definition. We say that a random variable X , or its distribution F , is *infinitely divisible* (id) if for each $n = 1, 2, \dots$, X has the same distribution as the sum of n independent identically distributed random variables. We write I for the class of infinitely divisible distributions.

It turns out that I is also the class of limit laws of row-sums of triangular arrays, as follows. We say that $\{x_{nk}\}$ ($k = 1, \dots, k_n$, $n = 1, 2, \dots$) is a *triangular array* if for each n , the X_{nk} are independent; we say that the array is *uniformly asymptotically negligible* (uan, more briefly *negligible*, if for all $\epsilon > 0$,

$$P\left(\max_{1 \leq k \leq k_n} |X_{nk}| > \epsilon\right) \rightarrow 0 \quad (n \rightarrow \infty).$$

The following are equivalent:

- (i) X is infinitely divisible, $X \in I$;
- (ii) X is the limit law of the row-sums $\sum_k X_{nk}$ of some negligible triangular array.

The classic reference for this material is Gnedenko and Kolmogorov [GnK].

It turns out also that the CFs of distributions in I can be characterised explicitly: they are those of the form

$$E[e^{itX}] = \exp\left\{iat - \frac{1}{2}\sigma^2 t^2 + \int_{-\infty}^{\infty} (e^{ixt} - 1 - ixtI_{(-1,1)})d\nu(t)\right\}, \quad (LK)$$

where the (positive) measure ν , the *Lévy measure*, satisfies

$$\int \min(1, |x|^2)d\nu(x) < \infty$$

(here we omit 0 from the range of the integration – or, we can include it, perhaps at the cost of changing σ), a , the *drift*, is real, and σ , the *Gaussian component*, is ≥ 0 ; (a, σ, ν) is called the *characteristic triplet* of X .

Equation (LK) above is called the *Lévy-Khintchine formula* (Lévy in 1934, Khintchine¹ in 1937, following work of de Finetti in 1929 and 1930, Kolmogorov in 1932). We return to it (Ch. VI) in connection with stochastic processes – *Lévy processes*. It gives a *semi-parametric* representation – think of (a, σ) as the parametric part and ν as the non-parametric part.²

Note. 1. In the integrand, we need three terms near the origin, but only two terms away from the origin. As we shall see later, the Lévy measure ν governs the *jumps* of the relevant Lévy process. We distinguish between the ‘big’ jumps (only finitely many of these in finite time), and the ‘little’ jumps

¹Khintchine as he wrote here in French; Khinchin is the usual transliteration of his name into English

²Here we follow the British usage of regarding a parameter as finite-dimensional. In Russian usage, the triplet would be a parametric description.

(there may be infinitely many in finite time!) We ‘compensate’ the little jumps by subtracting the mean – hence the $I_{(-1,1)}$. Actually, the ‘1’ here is arbitrary: any $c \in (0, \infty)$ would do, but $c = 1$ is customary and convenient. 2. The a in the triplet corresponds to a *deterministic* part, *at* (t is the time), called the *drift*; the σ part corresponds to a Gaussian component (Brownian motion – see VI.1). Any of these three components may be absent.

A triangular array is a *two-suffix* entity (needing a matrix of distributions). If we specialise to the *one-suffix* case (needing a sequence of distributions), then in each row, all the X_{nk} have the same distribution. This restricts the class I of infinitely divisible distributions, and we obtain now the class SD of *self-decomposable* distributions. These have CFs of the more restricted form, where

$$\nu(dx) = k(x)dx/|x|, \quad k \text{ increasing on } (-\infty, 0), \text{ decreasing on } (0, \infty).$$

Again, this is a semi-parametric description.

We can specialise even further, and have an array depending on only *one* distribution, F say. We have X_1, X_2, \dots iid with law F , and form the sequence of partial sums

$$S_n := X_1 + \dots + X_n;$$

then $S := \{S_n\}$ is called a *random walk* with *step-length* distribution F , or *generated by* F , $\{S_n\} \sim F$. Just as in the CLT, we seek to centre and scale so as to get a non-degenerate limit law. we ask for a non-degenerate limit of

$$(S_n - a_n)/b_n,$$

with a_n real, $b_n > 0$ (in the CLT $a_n = n\mu$ and $b_n = \sigma\sqrt{n}$ with μ the mean and σ^2 the variance, but here we need not have a mean or variance). So we get a *parametric* description, with four parameters – two essential, two not. *Type: location and scale.*

In one dimension, the mean μ gives us a natural measure of *location* for a distribution. The variance σ^2 , or standard deviation (SD) σ , give us a natural measure of *scale*.

Note. The variance has much better mathematical properties (e.g., it adds over independent, or even uncorrelated, summands). But the SD has the *dimensions* of the random variable, which is better from a physical point of view. As moving between them is mathematically trivial, we do so at will, without further comment.