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Author
Politis, D N

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Can the stock market be linearized?

Dimitris N. Politis*

* Depts. of Mathematics and Economics
University of California at San Diego
La Jolla, CA 92093-0112, USA
e-mail: dpolitis@ucsd.edu

Abstract: The evolution of financial markets is a complicated real-world phenomenon that ranks at the top in terms of difficulty of modeling and/or prediction. One reason for this difficulty is the well-documented nonlinearity that is inherently at work. The state-of-the-art on the nonlinear modeling of financial returns is given by the popular ARCH (Auto-Regressive Conditional Heteroscedasticity) models and their generalisations but they all have their short-comings. Foregoing the goal of finding the ‘best’ model, we propose an exploratory, model-free approach in trying to understand this difficult type of data. In particular, we propose to transform the problem into a more manageable setting such as the setting of linearity. The form and properties of such a transformation are given, and the issue of one-step-ahead prediction using the new approach is explicitly addressed.

Introduction

Consider data $X_1, \ldots, X_n$ arising as an observed stretch from a financial returns time series $\{X_t\}$ such as the percentage returns of a stock index, stock price or foreign exchange rate; the returns may be daily, weekly, or calculated at different (discrete) intervals. The returns $\{X_t\}$ are typically assumed to be strictly stationary having mean zero which—from a practical point of view—implies that trends and/or other nonstationarities have been
successfully removed.

At the turn of the 20th century, pioneering work of L. Bachelier [1] suggested the Gaussian random walk model for (the logarithm of) stock market prices. Because of the approximate equivalence of percentage returns to differences in the (logarithm of the) price series, the direct implication was that the returns series \( \{X_t\} \) can be modeled as independent, identically distributed (i.i.d.) random variables with Gaussian \( N(0,\sigma^2) \) distribution. Although Bachelier’s thesis was not so well-received by its examiners at the time, his work served as the foundation for financial modeling for a good part of the last century.

The Gaussian hypothesis was first challenged in the 1960s when it was noticed that the distribution of returns seemed to have fatter tails than the normal [11]. Recent work has empirically confirmed this fact, and has furthermore suggested that the degree of heavy tails is such that the distribution of returns has finite moments only up to order about two [14] [20].

Furthermore, in an early paper of B. Mandelbrot [16] the phenomenon of ‘volatility clustering’ was pointed out, i.e., the fact that high volatility days are clustered together and the same is true for low volatility days; this is effectively negating the assumption of independence of the returns in the implication that the absolute values (or squares) of the returns are positively correlated.

For example, Figure 1 depicts the daily returns of the S&P500 index from August 30, 1979 to August 30, 1991; the extreme values associated with the crash of October 1987 are very prominent in the plot. Figure 2 (a) is a ‘correlogram’ of the S&P500 returns, i.e., a plot of the estimated autocorrelation function (ACF); the plot is consistent with the hypothesis of uncorrelated returns. By contrast, the correlogram of the squared returns
Figure 1: Daily returns of the S&P500 index spanning the period 8-30-1979 to 8-30-1991.

of Figure 2 (b) shows some significant correlations thus lending support to the ‘volatility clustering’ hypothesis.

The celebrated ARCH (Auto-Regressive Conditional Heteroscedasticity) models of 2003 Nobel Laureate R. Engle [10] were designed to capture the phenomenon of volatility clustering by postulating a particular structure of dependence for the time series of squared returns \( \{ X_t^2 \} \). A typical ARCH\( (p) \) model is described by the equation:

\[
X_t = Z_t \sqrt{a + \sum_{i=1}^{p} a_i X_{t-i}^2}
\]  \( \text{(1)} \)

where \( a, a_1, a_2, \ldots \) are nonnegative real-valued parameters, \( p \) is a nonnegative integer indicating the ‘order’ of the model, and the series \( \{ Z_t \} \) is assumed to be i.i.d. \( N(0, \sigma^2) \). Bachelier’s model is a special case of the ARCH\( (p) \) model; just let \( a_i = 0 \) for all \( i \), effectively implying a model of order zero.
Figure 2: (a) Correlogram of the S&P500 returns. (b) Correlogram of the S&P500 squared returns.
Model (1) beautifully captures the phenomenon of volatility clustering in a simple equation at the same time implying a marginal distribution for the \{X_t\} returns that has heavier tails than the normal. Viewed differently, the ARCH(p) model may be considered an attempt to ‘normalize’ the returns, i.e., to reduce the problem to a model with normal residuals (the \(Z_t\)s). In that respect though the ARCH model (1) is only partially successful as empirical work suggests that the ARCH residuals often exhibit heavier tails than the normal; the same is true for the ARCH spin-off models such as GARCH, EGARCH, etc. [4] [24].

Nonetheless, the goal of normalization is most worthwhile and it is indeed achievable as will be shown in the sequel where the connection with the issue of nonlinearity of stock market returns will also be brought forward.

**Linear and Gaussian time series**

Consider a mean zero, stationary time series \(\{Y_t\}\). The most basic tool for quantifying the inherent strength of dependence is given by the autocovariance function \(\gamma(k) = E Y_t Y_{t+k}\) and the corresponding Fourier series \(f(w) = (2\pi)^{-1} \sum_{k=-\infty}^{\infty} \gamma(k) e^{-iwk}\); the latter function is termed the spectral density. We can also define the autocorrelation function (ACF) as \(\rho(k) = \gamma(k)/\gamma(0)\). If \(\rho(k) = 0\) for all \(k > 0\), then the series \(\{Y_t\}\) is said to be a white noise, i.e., an uncorrelated sequence; the reason for the term ‘white’ is the constancy of the resulting spectral density function.

The ACF is the sequence of second order moments of the variables \(\{Y_t\}\); more technically, it represents the second order cumulants [6]. The third order cumulants are given by the function \(\Gamma(j,k) = E Y_t Y_{t+j} Y_{t+k}\) whose Fourier series \(F(w_1, w_2) = (2\pi)^{-2} \sum_{j=-\infty}^{\infty} \sum_{k=-\infty}^{\infty} \Gamma(j,k) e^{-iw_1 j-iw_2 k}\)
is termed the \textit{bispectral density}. We can similarly define the cumulants of higher order, and their corresponding Fourier series that constitute the so-called higher order spectra.

The set of cumulant functions of \textit{all} orders, or equivalently the set of all higher order spectral density functions, is a complete description of the dependence structure of the general time series \{Y_t\}. Of course, working with an infinity of functions is very cumbersome; a short-cut is desperately needed, and presented to us by the notion of linearity.

A time series \{Y_t\} is called \textit{linear} if it satisfies an equation of the type:

\begin{equation}
Y_t = \sum_{k=-\infty}^{\infty} \beta_k Z_{t-k} \tag{2}
\end{equation}

where the coefficients \(\beta_k\) are (at least) square-summable, and the series \{Z_t\} is i.i.d. with mean zero and variance \(\sigma^2\). Eq. (2) is not to be confused with the Wold decomposition that all purely nondeterministic time series possess. In the Wold decomposition the ‘error’ series \{Z_t\} is only assumed to be a white noise and not i.i.d.; the latter assumption is much stronger.

Linear time series are easy objects to work with since the totality of their dependence structure is perfectly captured by a single entity, namely the sequence of \(\beta_k\) coefficients. To elaborate, the autocovariance and spectral density of \{Y_t\} can be calculated to be \(\gamma(k) = \sigma^2 \sum_{s=-\infty}^{\infty} \beta_s \beta_{s+k}\) and \(f(w) = (2\pi)^{-1} \sigma^2 |\beta(w)|^2\) respectively where \(\beta(w)\) is the Fourier series of the \(\beta_k\) coefficients, i.e., \(\beta(w) = \sum_{k=-\infty}^{\infty} \beta_k e^{iwk}\). In addition, the bispectral density is simply given by

\begin{equation}
F(w_1, w_2) = (2\pi)^{-2} \mu_3 \beta(-w_1)\beta(-w_2)\beta(w_1 + w_2) \tag{3}
\end{equation}

where \(\mu_3 = EZ_t^3\) is the 3rd moment of the errors. Similarly, all higher order spectra can be calculated in terms of \(\beta(w)\).
The prime example of a linear, mean zero, time series is given by the Auto-Regressive (AR) family [28] in which the time series \( \{ Y_t \} \) has a linear representation with respect to its own lagged values, namely

\[
Y_t = \sum_{k=1}^{p} \theta_k Y_{t-k} + Z_t
\]

(4)

with the error process \( \{ Z_t \} \) being i.i.d. as in eq. (2). AR modeling lends itself ideally to the problem of prediction of future values of the time series.

For concreteness, let us focus on the one-step-ahead prediction problem, i.e., predicting the value of \( Y_{n+1} \) on the basis of the observed data \( Y_1, \ldots, Y_n \), and denote by \( \hat{Y}_{n+1} \) the optimal (with respect to Mean Squared Error) predictor. In general, we can write \( \hat{Y}_{n+1} = g_n(Y_1, \ldots, Y_n) \) where \( g_n(\cdot) \) is an appropriate function. As can be shown [2], the function \( g_n(\cdot) \) that achieves this optimal prediction is given by the conditional expectation, i.e., \( \hat{Y}_{n+1} = E(Y_{n+1}|Y_1, \ldots, Y_n) \). Thus, to implement the one-step-ahead prediction in a general nonlinear setting requires knowledge (or accurate estimation) of the unknown function \( g_n(\cdot) \) which is far from trivial [13] [27] [26].

In the case of a causal [7] AR model however, it is easy to show that the function \( g_n(\cdot) \) is actually linear, and that \( \hat{Y}_{n+1} = \sum_{k=1}^{p} \theta_k Y_{n+1-k} \). Note furthermore the property of ‘finite memory’ in that the prediction function \( g_n(\cdot) \) is only sensitive to its last \( p \) arguments. Although the ‘finite memory’ property is specific to finite-order causal AR (and Markov) models, the linearity of the optimal prediction function \( g_n(\cdot) \) is a property shared by a large subset of the class of linear time series, namely the class of all causal and invertible, i.e., “minimum-phase” [22], ARMA models with i.i.d. innovations.

Gaussian series form another most interesting subset of the class of linear time series; they occur when the series \( \{ Z_t \} \) of eq. (2) is i.i.d. \( N(0, \sigma^2) \),
and they too exhibit the useful linearity of the optimal prediction function $g_n(\cdot)$. Furthermore, in the Gaussian case all spectra of order higher than two are identically zero; it follows that all dependence information is concentrated in the spectral density $f(w)$. Thus, the investigation of a Gaussian series’ dependence structure can focus on the simple study of second order properties, namely the ACF $\rho(k)$ and/or the spectral density $f(w)$. For example, an uncorrelated Gaussian series, i.e., one satisfying $\rho(k) = 0$ for all $k$, necessarily consists of independent random variables.

To some extent, this last remark can be generalised to the linear setting: if a linear time series is deemed to be uncorrelated, then practitioners typically infer that it is independent as well. Strictly speaking, however, this inference is only valid for the aforementioned class of causal and invertible ARMA models [5]. Note that to check/test whether an estimated ACF, denoted by $\hat{\rho}(k)$, is significantly different from zero, the Bartlett confidence limits are typically used—see e.g. the bands in Figure 2 (a); but those too are valid only for linear time series [17].

**Linearizing or normalizing the stock market?**

It should come as no surprise that a simple parametric model as (1) might not perfectly capture the behavior of a complicated real-world phenomenon such as the evolution of financial returns that—almost by definition of market ‘efficiency’—ranks at the top in terms of difficulty of modeling/prediction. As a consequence, researchers have recently been focusing on nonparametric methods for the analysis of financial time series; see [12] for a detailed overview. For example, consider the nonparametric ARCH
model defined by the equation:

$$X_t = g_p(X_{t-1}, \ldots, X_{t-p}) Z_t$$  \hspace{1cm} (5)

where $Z_t$ is i.i.d. $(0, \sigma^2)$, and $g_p$ is an unknown smooth function to be estimated from the data.

Despite its nonparametric character, eq. (5) is nevertheless just another model attempting to fully capture/describe the probabilistic characteristics of the time series in question which is perhaps an overly ambitious task. Foregoing the goal of finding the ‘best’ model, we may instead resort to an exploratory, model-free approach in trying to understand this difficult type of data. In particular, we may attempt to transform the problem into a more manageable setting such as the setting of linearity.

Consider again the financial returns data $\underline{X_n} = (X_1, \ldots, X_n)$, and a transformation of the type $\underline{V_n} = H(\underline{X_n})$ where $\underline{V_n}$ is also $n$-dimensional. Ideally, we would like the transformed series $\underline{V_n} = (V_1, \ldots, V_n)$ to be linear since, as mentioned before, such time series are easy to work with.

However, just asking for linearity of the transformed series is not enough. For example, the naive transformation $V_t = sign(X_t)$ may be thought of as a linearizing transformation since, by the efficient market hypothesis, $sign(X_t)$ is i.i.d. (taking the values +1 and -1 with equal probability), and therefore linear. Nevertheless, in spite of the successful linearization, the $sign$ transformation is not at all useful as the passage from $\underline{X_n}$ to $\underline{V_n}$ is associated with a profound loss of information.

To avoid such information loss “due to processing” [9], we should further require that the transformation $H$ be in some suitable sense invertible, allowing us to work with the linear series $V_t$ but then being able to recapture the original series by the inverse transformation $H^{-1}(\underline{V_n})$. Interestingly, the key
to finding such a transformation is asking for more: look for a normalizing (instead of just linearizing) information preserving transformation.

We now show how this quest may indeed be fruitful using the ARCH equation (1) as a stepping stone. To this end, define \( V_t = X_t / s_t \) for \( t = 1, 2, \ldots, p \), and

\[
V_t = \frac{X_t}{\sqrt{\alpha s_{t-1}^2 + a_0 X_t^2 + \sum_{i=1}^{p} a_i X_{t-i}^2}} \quad \text{for} \quad t = p + 1, p + 2, \ldots, n; \quad (6)
\]

in the above, \( \alpha, a_0, a_1, \ldots, a_p \) are nonnegative real-valued parameters, and \( s_{t-1}^2 \) is an estimator of \( \sigma_X^2 = Var(X_1) \) based on the data up to (but not including) time \( t \). Under the zero mean assumption for \( X_t \), the natural estimator is \( s_{t-1}^2 = (t - 1)^{-1} \sum_{k=1}^{t-1} X_k^2 \).

The invertibility of the above transformation is manifested by solving eq. (6) for \( X_t \), thus obtaining:

\[
X_t = \frac{V_t}{\sqrt{1 - a_0 V_t^2}} \sqrt{\alpha s_{t-1}^2 + \sum_{i=1}^{p} a_i X_{t-i}^2} \quad \text{for} \quad t = p + 1, p + 2, \ldots, n. \quad (7)
\]

Given the initial conditions \( X_1, \ldots, X_p \), the information set \( \mathcal{F}_n^X = \{X_t, 1 \leq t \leq n\} \) is equivalent to the information set \( \mathcal{F}_n^V = \{V_t, 1 \leq t \leq n\} \). To see this, note that with eq. (7) we can recursively re-generate \( X_t \) for \( t = p + 1, p + 2, \ldots, n \) using just \( \mathcal{F}_n^V \) and the initial conditions; conversely, eq. (6) defines \( V_t \) in terms of \( \mathcal{F}_n^X \).

Equation (6) describes the proposed normalizing (and therefore also linearizing) transformation, i.e., the operator \( H \) in \( V_n = H(X_n) \). Note that formally the main difference between eq. (6) and the ARCH eq. (1) is the presence of the term \( X_t^2 \) paired with the coefficient \( a_0 \) inside the square root; this is a small but crucial difference without which the normalization goal may not always be feasible [18] [19].
Despite its similarity to model (1), eq. (6) is not to be interpreted as a “model” for the \( \{X_t\} \) series. In a modeling situation, the characteristics of the model are pre-specified (e.g., errors that are i.i.d. \( N(0, \sigma^2) \), etc.), and standard methods such as Maximum Likelihood or Least Squares are used to fit the model to the data. By contrast, eq. (6) does not aspire to fully describe the probabilistic behavior of the \( \{X_t\} \) series. The order \( p \) and the vector of nonnegative parameters \((\alpha, a_0, \ldots, a_p)\) are chosen by the practitioner with just the normalization goal in mind, i.e., in trying to render the transformed series \( \{V_t\} \) as close to normal as possible; here, ‘closeness’ to normality can be conveniently measured by the Shapiro-Wilk (SW) test statistic \([23]\) or its corresponding \( P \)-value.

It is advantageous (and parsimonious) in practice to assign a simple structure of decay for the \( a_k \) coefficients. The most popular such structure—shared by the popular GARCH(1,1) model \([3]\)—is associated with an exponential rate of decay, i.e., to postulate that \( a_k = Ce^{-dk} \) for some positive constants \( d \) and \( C \) which—together with the parameter \( \alpha \)—are to be chosen by the practitioner.

Taking into account the convexity requirement \( \alpha + \sum_{k=0}^{p} a_k = 1 \), the exponential coefficients scheme effectively has only two free parameters that can be chosen with the normalization goal in mind, i.e., chosen to maximise the SW statistic calculated on the transformed series \( V_t \) or linear combinations thereof—the latter to ensure normality of joint distributions.

As it turns out, the normalization goal can typically be achieved by a great number of combinations of these two free parameters, yielding an equally great number of possible normalizing transformations. Among those equally valid normalizing transformations the simplest one corresponds to letting \( \alpha = 0 \). Alternatively, the value of \( \alpha \) may be chosen by an additional
optimisation criterion driven by an application of interest such as predictive ability.

For illustration, let us revisit the S&P500 returns dataset. The normalising transformation with \( a_k = Ce^{-dk} \) and the simple choice \( \alpha = 0 \) is achieved with \( d = 0.0675 \); the resulting transformed \( V \)-series is plotted in Figure 3 which should be compared to Figure 1. Not only is the phenomenon of volatility clustering totally absent in the transformed series but the outliers corresponding to the crash of October 1987 are hardly (if at all) discernible.

**Quantifying the degree of nonlinearity and nonnormality**

There are many indications pointing to the nonlinearity of financial returns. For instance, the fact that returns are uncorrelated but not independent is a good indicator; see e.g. Figure 2 (a) and (b). Notably, the ARCH model
and its generalisations are all models for nonlinear series.

To quantify nonlinearity, it is useful to define the new function

$$K(w_1, w_2) = \frac{|F(w_1, w_2)|^2}{f(w_1)f(w_2)f(w_1 + w_2)}.$$  \hspace{1cm} (8)

From eq. (3), it is apparent that if the time series is linear, then \(K(w_1, w_2)\) is the constant function—equal to \(\mu_3^2/(2\pi \sigma^6)\) for all \(w_1, w_2\); this observation can be used in order to test a time series for linearity [25] [15]. In the Gaussian case we have \(\mu_3 = 0\) and therefore \(F(w_1, w_2) = 0\) and \(K(w_1, w_2) = 0\) as well.

Let \(\hat{K}(w_1, w_2)\) denote a data-based nonparametric estimator of the quantity \(K(w_1, w_2)\). For our purposes, \(\hat{K}(w_1, w_2)\) will be a kernel smoothed estimator based on infinite-order flat-top kernels that lead to improved accuracy [21]. Figure 4 shows a plot of \(\hat{K}(w_1, w_2)\) for the S&P500 returns; its nonconstancy is direct evidence of nonlinearity. By contrast, Figure 5 shows a plot of \(\hat{K}(w_1, w_2)\) for the normalized S&P500 returns, i.e., the \(V\)-series. Note that, in order to show some nontrivial pattern, the scale on the vertical axis of Figure 5 is 500 times smaller than that of Figure 4. Thus, it should come as no surprise that the function \(\hat{K}(w_1, w_2)\) for the normalized S&P500 returns is not statistically different from the zero function, lending support to the fact that the transformed series is linear with distribution symmetric about zero; the normal is such a distribution but it is not the only one.

To further delve into the issue of normality, recall the aforementioned Shapiro-Wilk (SW) test which effectively measures the lack-of-fit of the
quantile-quantile plot (QQ-plot) to a straight line. Figure 6 shows the QQ-plot of the S&P500 returns; it is apparent that a straight line is not a good fit. As a matter of fact the SW test yields a P-value that is zero to several decimal points—the strongest evidence of nonnormality of stock returns. By contrast, the QQ-plot of the normalized S&P500 returns can be very well approximated by a straight line: the $R^2$ associated with the plot in Figure 7 is 0.9992, and the SW test yields a P-value of 0.153 lending strong support to the fact that the transformed series is indistinguishable from a Gaussian series.

The proposed transformation technique has been applied to a host of different financial datasets including returns from several stock indices, stock prices and foreign exchange rates. Invariably, it was proven successful in its dual goal of normalization and linearization.
APPENDIX: Prediction using the transformation technique

For concreteness, we focus on the problem of one-step ahead prediction, i.e., prediction of a function of the unobserved return $X_{n+1}$, say $h(X_{n+1})$, based on the observed data $\mathcal{F}_n^X = \{X_t, 1 \leq t \leq n\}$. Our normalizing transformation affords us the opportunity to carry out the prediction in the $V$-domain where the prediction problem is easiest since the problem of optimal prediction reduces to linear prediction in a Gaussian setting.

The prediction algorithm is outlined as follows:

- Calculate the transformed series $V_1, \ldots, V_n$ using eq. (6).
- Calculate the optimal predictor of $V_{n+1}$, denoted by $\hat{V}_{n+1}$, given $\mathcal{F}_n^V$. This predictor would have the general form $\hat{V}_{n+1} = \sum_{i=0}^{q-1} c_i V_{n-i}$. The $c_i$ coefficients can be found by Hilbert space projection techniques, or by simply fitting the causal AR model

$$V_{t+1} = \sum_{i=0}^{q-1} c_i V_{t-i} + \epsilon_{t+1}. \quad (9)$$

to the data where $\epsilon_t$ is i.i.d. $N(0, \sigma^2)$. The order $q$ can be chosen by an information criterion such as AIC or BIC [8].
Note that eq. (7) suggests that $h(X_{n+1}) = u_n(V_{n+1})$ where $u_n$ is given by

$$u_n(V) = h \left( \frac{V}{\sqrt{1 - a_0 V^2}} - \frac{\alpha_1 t - 1}{\sum_{i=1}^{p} a_i X_{n+1-i}^2} \right).$$

Thus, a quick-and-easy predictor of $h(X_{n+1})$ could then be given by $u_n(\hat{V}_{n+1})$.

A better predictor, however, is given by the center of location of the distribution of $u_n(V_{n+1})$ conditionally on $F_{V_n}$. Formally, to obtain an optimal predictor, the optimality criterion must first be specified, and correspondingly the form of the predictor is obtained based on the distribution of the quantity in question. Typical optimality criteria are $L_2$, $L_1$ and 0/1 losses with corresponding optimal predictors the (conditional) mean, median and mode of the distribution. For reasons of robustness, let us focus on the median of the distribution of $u_n(V_{n+1})$ as such a center of location.

Using eq. (9) it follows that the distribution of $V_{n+1}$ conditionally on $F_{V_n}$ is approximately $N(\hat{V}_{n+1}, \hat{\sigma}^2)$ where $\hat{\sigma}^2$ is an estimate of $\sigma^2$ in (9). Thus, the median-optimal one-step-ahead predictor of $h(X_{n+1})$ is the median of the distribution of $u_n(V)$ where $V$ has the normal distribution $N(\hat{V}_{n+1}, \hat{\sigma}^2)$ truncated to the values $\pm 1/\sqrt{a_0}$; this median is easily computable by Monte-Carlo simulation.

The above Monte-Carlo simulation actually creates a predictive distribution for the quantity $h(X_{n+1})$. Thus, we can go a step further from the notion of a point-predictor: clipping the left and right tail of this predictive distribution, say $\delta \cdot 100\%$ on each side, a $(1 - 2\delta)100\%$ prediction interval for $h(X_{n+1})$ is obtained.

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References


Figure 4: Plot of $\hat{K}(\omega_1, \omega_2)$ vs. $(\omega_1, \omega_2)$ for the S&P500 returns.

Figure 5: Plot of $\hat{K}(\omega_1, \omega_2)$ vs. $(\omega_1, \omega_2)$ for the normalized S&P500 returns.