# The Stationary Distribution of Wealth with Random Shocks

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Received:

# January 2002

Abstract. A convergence model with wealth accumulation subject to i.i.d. random shocks is examined. The transfer function shows what  $k_{t+1}$  - wealth at t + 1 - would be, given  $k_t$ , with no shock. It has a positive slope, but its concavity/convexity is indeterminate. The stationary distribution of wealth satis...es a Fredholm integral equation. This distribution can be examined by direct analysis of the wealth-accumulation stochastic process and via the Fredholm equation. The analysis resembles some econometric theory of time series. Economic theory forces consideration of a broad range of cases, including some which violate  $\bar{-}$ -convergence. "Twin peaks" in the stationary distribution cannot be excluded.

1980 AMS Subject Classi...cation: JEL Classi...cation: D3 E1 Keywords and Phrases: Convergence, stochastic process, wealth distribution

## 1. The Stationary Distribution of Wealth with Random Shocks

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## 1.1. Introduction

The Neoclassical convergence model has been in‡uential in recent years. See in particular Barro (1991), Barro and Sala-i-Martin (1992 and 1995). In its basic form the model says that all units within the relevant population<sup>2</sup> - which might be countries, regions, even individual families, as desired - tend to converge to a common level of capital and output per head. The theory leads to a relationship similar to:

$$\mathbf{k}_{t+1} = \mathbf{h}[\mathbf{k}_t] \tag{1}$$

where  $k_t$  is the logarithm of wealth (or income), and there is a unique stable value of  $k = k^n$ , such that  $k^n = h[k^n]$ . The function h[l] will be called the Transfer Function, because it indicates how the underlying dynamics of the system take it from  $k_t$  to  $k_{t+1}$ . This model refers to wealth, that is capital including human capital. And the economic theory to which the leading writers appeal applies to wealth accumulation. Empirical studies, however, typically use income rather than wealth, because income is far better measured. In what follows I shall always refer to wealth, even when discussing studies which use income. When it is precisely an income measure which is used, income may be interpreted as a proxy for wealth. Nothing essential in what follows is an ected by the income-wealth distinction.

Much of the convergence literature reads as if all units within the population will move closer and closer to k<sup>a</sup>. Or, in an approach, which Robert Barro has built and elaborated, individual countries have di<sup>a</sup>erent k<sup>a</sup> values, depending on numerous additional variables, such as democracy or the share of government expenditure in GDP. Barro (1997) is a convenient reference.

As an econometric model based on (1) has to include an error term, it is essential that random departures from the strict model be taken into account. Therefore (1) is modi...ed to become:

$$k_{t+1} = h[k_t] + {}^{2}_{t}$$
(2)

This is the equation of a non-linear stochastic process. Economic theory will impose various restrictions on its form, see below. An argument may arise as to whether the shocks

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Helpful comments from David Cox, Bent Nielsen, Derek Ritzman and Neil Shephard are gratefully acknowledged.

<sup>&</sup>lt;sup>2</sup>For countries many recent studies claim convergence to be conditional on qualifying properties, such as economic openess, or the protection of property rights.

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 $_{t}^{2}$  show true exogenous random shocks, or whether they represent the impact of missing variables, or both. Note that if shocks re‡ect missing variables they will not be i.i.d. The presence of unobserved exects for a particular country will surely be serially correlated. To follow that line would take the argument into the type of conditional convergence model promoted by Barro [see Barro (1997)]. Su¢ce it to say that in that model almost any stationary distribution is possible, depending upon the distribution of conditioning variables across the population of units observed.

Concentrating therefore on the case in which all units are essentially identical except for initial conditions, it will not be the case that all will converge to a common k¤ if random shocks constantly disturb the dynamic adjustment process which is represented in a simple form by (1). If random errors are important in their magnitudes they a¤ect the process of convergence. And if shocks are important, there will never be strict convergence. Then an interesting issue arises concerning which the current literature is notably silent. What will the stationary distribution of k values be when a convergent process such as (1) is modi...ed by random shocks? The stationary distribution is that which replicates itself under the combined e¤ects of the process (1) and the addition of random shocks. In particular one may ask, what are the respective contributions to the form of the stationary wealth distribution of:

- <sup>2</sup> the shape of the function h [<sup>¢</sup>]
- <sup>2</sup> the probability distribution of shocks

The problem addressed in this paper has been introduced in terms of the well-known convergence model of Barro and others. Formally, however, the same problem is encountered whenever an adjustment process of the form of (1) is shocked. Thus  $k_t$  might be the advertising budget of a ...rm with (1) representing slow adjustment of that budget to a long-run equilibrium level equal to  $k^{\mu}$ . If one can imagine such an adjustment process being regularly shocked, a stationary frequency distribution of k values shows the long-run probability that k will be found to be in any particular interval at a randomly chosen time.

Note that, as this last example indicates, treating  $k_t$  as being continuously distributed does not imply that k takes an uncountable in...nity of values at any time. The integral over an interval in the density of the distribution may measure the probability of ...nding k within that interval of values.

Returning to wealth accumulation, consider an empirical study based on a linearized and re-arranged version of the relation (2), viz:

$$\kappa_{t+1} i k_t = f[k_t] i k_t = R_i^{-1} k_t + L_t^{2}$$
(3)

With the normal ...nding 0 < - < 1, equation (3) says that on average poor units (countries) grow faster than rich units. This has been called -i convergence. This concept of

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convergence is not the same as  $\frac{3}{4}$  i convergence, which means that the variance of the population of k values declines over time<sup>3</sup>. Friedman (1992) claims that interpreting a negative coe¢cient on k<sub>t</sub> in a regression like (3) as convergence exhibits "Galton's Fallacy", on the ground that a negative coe¢cient is consistent with no tendency for the variance of k<sub>t</sub> to decrease with time: it may even increase<sup>4</sup>.

That point ...ts well with the argument of this paper. If  $k_t$  is distributed according to a stationary distribution, which replicates itself, there is plainly no ¾-convergence. Take a value of  $k_t$  far from  $k^{\underline{n}}$ . The expected value of  $k_{t+1}$  conditional on such a value of  $k_t$  will be closer to  $k^{\underline{n}}$ . It is ®  $i (\ i \ 1) \notin k_t$  when (3) applies. Here, while individual dispersed units tend to converge, their density is made good by units, including the less-dispersed, pushed outwards by random shocks. A stationary distribution of k values is invariant over time in the sense that it reproduces itself next period, although individual values will vary, partly systematically, showing  $\$ -convergence, and partly randomly, due to stochastic realisations of  ${}^2_t$ .

In an important contribution Quah (1993) considers income generation as a pure Markov process. See also Quah (1996a) and (1996b). He looks at observed transition patterns, without considering the stationary distribution. Also, he does not derive his Markov transitions from economic theory. The aim is to confront theoretically derived convergence models with the hard facts shown in the data. He ...nds a tendency for convergence within two groups - high and low income. He names this ...nding "twin peaks". There are non-negligible probabilities that a country will shift from one group to the other, but these transition probabilities are too low to iron out the twin peaks in the distribution. Quah's approach revives an older tradition in which wealth distribution, its change over time, and long-run equilibrium distributions, are modelled from some speci...cation of the process which transfers individuals from one wealth-state to another. See Champernowne (1953), Steindl (1972) and Wold and Whittle (1957).

The paper is organized as follows. Section 2 is concerned with clarifying terminology and explaining central concepts. These include asymptotic properties, stationarity of stochastic processes and what is here called long-memory. Sextion 3 de...nes the transfer function and notes how the presence of random shocks may a¤ect optimal adjustment. Section 4 examines the conditions under which a transfer function will exhibit <sup>-</sup>-convergence. Section 5, entitled Why do the Poor not save More?, mounts an all-out assault on the assumption of invariance of the elasticity of marginal utility. This assumption, which has become almost standard in the literature, is shown to be both highly restrictive and empirically questionable. Section 1.6 shows how the distribution of wealth is the sum of two o¤setting e¤ects, called convergence and scattering. For a stationary distribution these two e¤ects

<sup>&</sup>lt;sup>3</sup>For a clear exposition of the two concepts of convergence, and empirical discussion, see Sala-i-Martin (1996).

<sup>&</sup>lt;sup>4</sup>See also Quah (1993) and Bliss (1999) and (2000).

must precisely o¤set each other. Therefore their e¤ects on each moment of the distribution must be of equal magnitude and opposite sign. Section 7 examines the joint e¤ects of transfer and scattering and shows that these imply the satisfaction of integral equations. In particular a stationary distribution satis...es a particular Fredholm integral equation in which the same function appears outside and under the integral. Section 8 is concerned with special cases in which one of the two grand e¤ects - transfer and scattering - is assumed to take a simple form, so that the in‡uence of important features in the other e¤ect can emerge clearly. Section 9 extends this last approach to elucidate two important features of a stationary distribution - asymmetry and single-peakedness. Section 10 concludes.

## 1.2. Asymptotic Properties and Stationary Distributions

This section is concerned with clarifying terminology and some concepts employed in the paper.

Definition 1. A stationary distribution is one for which the probability of ...nding k within any interval is the same next period as in the current period.

A stationary wealth distribution may be compared to a liquid in thermodynamic equilibrium. Its macro properties are invariant. Temperature and pressure in any part of the liquid do not alter with time. Even so, microscopic inspection will reveal local random changes (Brownian motion) involving a small number of molecules. In practice these movements always average out to no change at the macroscopic level. The implication that a thermodynamic equilibrium is only static in a statistical sense - the expected value of the system next observation period is the same as the current distribution - applies equally to a wealth distribution. Given the possibility of suciently large shocks, everyone may be twice as wealthy tomorrow as they are today. Unfortunately we may have to wait a very long time for this happy outcome to occur. The remaining age of the Earth may not be sucient to make the probability as high as 0.01%.

A stationary distribution requires stationarity of the stochastic process. This is one of the most basic and important properties of a stochastic process. Put simply a stochastic process is stationary if the probability density which attaches to any sub-sequence of values is independent of the dates at which those events are observed. Then (2) de...nes a stationary process if the  $^{2}_{t}$  values are always drawn independently from the same distribution (are i.i.d. in the usual terminology). For in that case the sequence:

$$fk_t; \dots; k_{t+n}g \tag{4}$$

requires the random variable <sup>2</sup> to take speci...c values in the ordered sequence:

Then those same <sup>2</sup> values ordered as in (5) will generate the equally probable sequence:

$$fk_T$$
; ::::;  $k_{T+n}g$  (6)

starting at a di¤erent date.

Stationarity must be distinguished from a property which will play an important role below.

Definition 2. A stochastic process which generates the values  $fx_1$ ; ...;  $x_t$ ; ...; will be said to have long-memory if the probability that  $x_T \ge I$ , where I is a closed interval of values of x, is not asymptotically independent of  $x_1$  as t ! 1.

Long-memory corresponds to what in economic theory is sometimes called path-dependence. This means roughly that the economic system may be in a subset of its state space which it will never leave, although when it is not in that subset it never enters it. David (1993) discusses the implications of path-dependence in economic history. In another famous paper, David (1985), the same author argues that the QWERTY arrangement of keys on a standard typewriter is the result of historical accident, but is now locked in because the costs of change prohibit its substitution by a layout more e $\mathcal{C}$ cient for modern machines. Further examples of this type are easily found.<sup>5</sup> However, where a stochastic process describes the economy, the conditions for long-memory are more demanding than is the case with a deterministic economy, as will be made clear by the example which follows shortly. The opposite of long-memory is short-memory, taken here to mean simply that the condition speci...ed by De...nition 2 is not satis...ed, so that we have asymptotic independence of  $x_1$ .

Bradford DeLong (1999) describes a short-memory system as ergodic. In this paper the

term ergodic has been avoided. Too many meanings have been attached to it, ranging from the precise but highly abstract, to the vague and hand-waving (as Granger and Terasvirta (1993), p.10). Probably all would agree that a long-memory system is not ergodic. For that case we now have another precise term to hand.

Some of the distinctions made above are illustrated well by an example which falls outside Barro's model of optimal economic growth, which is based on the Ramsey model. The objective function is the standard:

<sup>&</sup>lt;sup>5</sup>Consider the keywork on a standard modern ‡ute. Theobald Boehm, who is credited with the design of the precursor of today's ‡utes, preferred an open G# design. This respects the simple principle that putting down ...ngers produces lower notes. However because ‡ute players of that time were used to closed G# ...ngering, closed G# ‡utes were made by Boehm and his followers, and today this ...ngering is standard for mass-produced ‡utes and teaching manuals. This is so despite the fact that closed G# gives rise to formidable intonation di¢culties for E in the third-octave, requiring more keywork to solve the problem. While soloists, who have their ‡utes hand-made to their own speci...cation, often choose to play open G# instruments, the huge economies of scale which attach to the closed G# design allow it to predominate.

$$\mathbf{Z}_{1} \quad [c(t)] e^{i t} dt \quad (7)$$

In place of the usual concave production function substitute a non-concave function which has the property that the marginal product of capital goes almost to zero for medium levels of capital, but above a range where the marginal product is low, the marginal product rises sharply, although it later falls away when capital is extremely high. This economy is like two concave economies stitched together, yet su¢ciently separated to allow optimal growth paths to exist in two regions and to converge to di¤erent levels of k. Figure 1 illustrates the form of the h[¢] function. The dynamics of the system when no random shocks intervene is plain. There are three values of k such that k = h[k], labelled A, B and C. B is unstable, the other two points, A and C are stable.

Without random shocks the system has long memory. Ignoring the zero-measure point B, all initial values of k lie in one of two basins of attraction from which k asymptotes either to A or to C, according to which basin applies. Furthermore, as the curve through B is steep, k is thrown some distance to the left or to the right away from B by the transfer process by itself. Then if the distribution of <sup>2</sup> values is uniform on [i a; a] with a small, the random e<sup>a</sup>ects can never overcome the powerful centrifugal force of the transfer process, and the system can never move from one basin of attraction to the other. Even with random shocks in this case, there is long-memory.

The next theorem builds on the example illustrated in Figure 1 in the following sense. The properties on which it depends are precisely those which are violated by the example. First a regularity condition on the distribution of <sup>2</sup> values makes the proof of the theorem straightforward, and does not exclude any case likely to be of interest.

Definition 3. Let a probability density function for values of the random shocks <sup>2</sup> be de...ned on [ $_i$  1; +1], and denoted ¼ (<sup>2</sup>). Then the distribution of shocks will be said to be Regular if:

(i)  $\frac{1}{4}$  (<sup>2</sup>) is continuous in <sup>2</sup>; and

(ii) given any four values 
$${}^{2}_{1}, {}^{2}_{2}, {}^{2}_{3}$$
 and  ${}^{2}_{4}$ , with  ${}^{2}_{1} < {}^{2}_{2} < {}^{2}_{3} < {}^{2}_{4}$ :  
 $Z {}^{2}_{2} \qquad Z {}^{2}_{4}$   
 ${}^{4}_{4}({}^{2}) d^{2} > 0$  and  ${}^{4}_{2}({}^{2}) d^{2} > 0$  (8)

implies:

$$Z_{2_{3}} = \frac{1}{2} \frac{1}{2}$$

The de...nition excludes distributions with "holes" in them; that is ranges with zero probability density enclosed between ranges with positive probability density. Empty tails, such as are seen with the uniform distribution are not excluded.

## Theorem 1. If

$$\mathbf{k}_{t+1} = \mathbf{h}\left[\mathbf{k}_t\right] \tag{10}$$

has a unique stable solution for a constant k (k = h[k]), and the values  ${}^{2}_{t}$  come from a regular i.i.d. generator, then

$$k_{t+1} = h[k_t] + {}^{2}_{t}$$
(11)

de...nes a short-memory process.

<u>Proof</u>: Is in two parts. Let the constant solution be  $k^{\underline{n}}$ . First it is shown that given any interval centered on  $k^{\underline{n}}$ ,  $I_{\underline{C}} = [k^{\underline{n}}_{\underline{i}} \, \underline{C}; k^{\underline{n}} + \underline{C}]$ , and any initial value  $k_1$ , the probability that  $k_t$  will never be in  $I_{\underline{C}}$  tends to zero as t ! 1. Secondly, given any closed interval  $I = [k^{\underline{0}}; k^{\underline{0}}]$  and an initial value  $k_1$ , the probability that  $k_T$  will be in I is a continuous function of  $k_1$ .

These two results together imply the theorem. The …rst result entails that for any two starting values of k,  $k_1^1$  and  $k_1^2$ , both the paths leading from these values will enter  $I_{\oplus}$  at some time (not necessarily the same time for each series) with limiting probability 1. Then we can reset the clocks for each realization of the process so that  $k_1^1$  and  $k_1^2$  will now both be in  $I_{\oplus}$  at t = 1. By further choosing  $\oplus$  suitably small, and now using the second part of the proof (continuity of the probability that  $k_T$  will be in I as a function of  $k_1$ ), we can make the probabilities that the two paths will each be in I at T (on the respective reset clocks) as close together as desired. This contradicts long-memory.

<u>Part 1</u> First suppose that  $k_1 > k = + C$ . Then for the transfer process not to take k into  $I_{C}$ , <u>either</u> the average of  ${}^2_t$  realizations must be positive, however large t, <u>or</u> any negative realizations of  ${}^2_t$  must be  $\cdot_i 2C$ . Therefore k starting from a value > k = +C not to enter  $I_C$  eventually requires a hole in the distribution of  ${}^2$  at least over the interval [i C; 0]. By symmetry, should k start from a value < k = +C, a hole in the distribution of  ${}^2$  must be found at least over the interval [0; +C]. In summary there must be a hole in the distribution of  ${}^2$  at least over the interval [i C; +C]. This contradicts the assumption that the distribution of  ${}^2$  is regular.

<u>Part 2</u> Take the interval  $I = [k^0; k^{00}]$  and an initial value  $k_1$ . Denote by Prhk<sub>t</sub> 2 I j  $k_1i$  the probability that  $k_t$  will be in I conditional on k taking the value  $k_1$  at t = 1. Then Prhk<sub>t</sub> 2 I j  $k_1i = {}_E \frac{1}{4} \binom{2}{}$ , where E is the set of all realizations of the <sup>2</sup> values such that  $k_1$  is transferred to  $k_t$  by the process (11). Now consider Prhk<sub>t</sub> 2 I j  $k_1 + \pm i$ , where  $\pm$  may be arbitrarily small, Then relative to any speci...c element of E, the realizations of <sup>2</sup> may be modi...ed so as to make k follow the sequence:

$$\frac{\frac{1}{2}}{k_1 + \pm; k_1 + \pm \frac{t_i}{t_i} \frac{2}{1}; k_1 + \pm \frac{t_i}{t_i} \frac{3}{1}; \dots; k_t}$$
 (12)

Notice that (12) establishes a one-to-one correspondence between paths from  $k_1$  to  $k_t$  and paths from  $k_1 + \pm$  to  $k_t$ . Plainly any path from  $k_1$  to  $k_t$  has a unique partner de...ned by

(12). Equally a path from  $k_1 + \pm$  to  $k_t$  has a unique partner given by:

$${}^{\frac{1}{2}}_{k_{1} j_{1} \pm; k_{1} j_{1} \pm \frac{t_{1} 2}{t_{1} 1}; k_{1} j_{1} \pm \frac{t_{1} 3}{t_{1} 1}; \dots; k_{t}$$
(13)

The importance of this one-to-one correspondence lies in the fact that in integrating over all paths to  $k_t$  to determine summed probabilities we will always be integrating over sets of equal measure regardless of starting point, so that only di¤erences in probabilities will matter.

For  $\pm$  su¢ciently small alterations to the <sup>2</sup> values can be made arbitrarily small, in which case the alteration to  $_{E}$  ¼ (<sup>2</sup>) may be made arbitrarily small. It follows that Prhk<sub>t</sub> 2 I j k<sub>1</sub>i is continuous in k<sub>1</sub> as required.¤

The proof is complete. ¤

If capital accumulation is a short-memory process we reach the type of striking conclusion, already noted by Quah, that if we only wait long enough, Bangladesh will be richer than the US in per capita terms at some t with probability 1. Barro's conditional convergence would apparently destroy that conclusion. However one might believe that the US will certainly become undemocratic, closed and without secure property rights, etc., as t ! 1. These speculations underline the curious nature of asymptotic results.

## 1.3. Convergence and the Transfer Function

It is convenient to work with the logarithm of wealth because it is not bounded below by zero, which makes possible in...nite-tail distributions, such as the normal. Obviously, were k to be normally distributed, wealth itself would be distributed as the log-normal distribution. Starting with the model (1), we add i.i.d. errors <sup>2</sup> with mean zero to obtain:

$$k_{t+1} = h[k_t] + {}^{2}_{t}$$
(14)

The equation (14) is a stochastic non-linear di¤erence equation in  $k_t$ . The study of the time series generated by di¤erence equations has become a central concern of modern econometric theory. Hamilton (1994) provides a wide-ranging account of this ...eld, and other texts, for instance Øksendal. (1998), treat the theory of stochastic processes at a more advanced mathematical level. Despite the existence of a substantial literature, some questions which the present economic theory brings to the foreground are not answered by the econometrics literature. Much of the latter concentrates on linear di¤erence equations, which are too limited for our needs here. Also while the stationary distribution is a natural long-run model for economic theory, econometricians have been more interested, understandably, in the estimation of equation parameters, or spectral analysis of time series generated by di¤erence equations, and even in models other than the i.i.d. case.

The transfer function shows how much capital would be held one period later, starting from a level  $k_t$ , were no random shock to arrive to throw the adjustment process o¤ its intended path. The agent starts with  $k_t$ , turns that into  $h[k_t]$  by saving, or dissaving as the case may be, and ends up with  $k_{t+1}$  after the shock has taken e¤ect. One could assume:

$$k_{t+1} = h[k_t + {}^2_t]$$
(15)

meaning that shocks a ect wealth before the adjustment decision is made. However (14) ...ts most simply with existing econometric approaches.

The model is altered more radically if one assumes:

$$k_{t+1} = h[k_t; {}^2_t]$$
(16)

which means that shocks a ect the adjustment process in a non-linear manner. This formulation represents a consequence of some problems examined Binder and Pesaran (1999), who consider what happens when a Solow, or an AK, model (but not an optimal saving model) includes stochastic technical progress or stochastic labour supply. These authors show that such changes undermine the short-memory of the process. In any case, (16) is a far more complex stochastic process than those that will be considered in the present paper.

It is important in interpreting (14) to understand what is implied by the i.i.d. assumption, and what is not implied by it. The additive i.i.d. shock entails that the value of  $h[\ell]$  is una<sup>x</sup>ected by the particular value taken by <sup>2</sup>. That does not imply that  $h[\ell]$  is una<sup>x</sup>ected by the distributional properties of <sup>2</sup>, in particular by the fact that <sup>2</sup> does not always take the value zero. As  $h[\ell]$  shows an optimal saving (adjustment) rule, that rule may be in‡uenced by the existence of uncertainty. Computing the properties of optimal saving rules under uncertainty is formidably di¢cult and will not be attempted below. As will shortly be seen, there is ambiguity for an important property of  $h[\ell]$  even when the saving decision takes no account of uncertainty. When  $h[\ell]$  shows an optimal saving rule which re‡ects the existence of uncertainty, it does not use information on the current value of <sup>2</sup>.

For stability one must have:

$$\frac{@h[k]}{@k} \Big|_{k=k^{n}} < 1$$
(17)

## 1.4. The Transfer Function and <sup>-</sup>-Convergence

Recall that <sup>-</sup>-convergence requires that the growth rate of k should decline as k increases. So with k being the logarithm of capital, we have strict <sup>-</sup>-convergence if:

$$\frac{\mathrm{dh}(k)}{\mathrm{dk}} \mathbf{i} \quad 1 < 0 \tag{18}$$

From (17) plus continuity we must have  $\bar{}$ -convergence close to  $k = k^{a}$ . Therefore if k tends asymptotically to  $k^{a}$ , we must have  $\bar{}$ -convergence in the limit. With stochastic shocks there will not be strict asymptotic convergence to  $k^{a}$ , so this last point is without force. In any case an asymptotic property cannot con...ne what may happen globally, see below.

Thinking on terms of the level of capital K, rather than its logarithm k, (18) can be written:

$$\frac{dK_{t+1}}{dK_t} \frac{1}{K_{t+1}} i \frac{1}{K_t} < 0$$
(19)

which is the same as requiring that the elasticity of next period's capital with respect to current capital shall be less than unity. Can this property, which would both entail <sup>-</sup>-convergence, and also help to detail features of the stationary distribution, be assumed? An examination of this seemingly technical question throws up points which the current literature on convergence has pushed aside.

Suppose that (1) shows the outcome of the optimal wealth accumulation of a Ramseystyle agent, who maximizes:

$$\mathbf{X}_{t=1}^{t_{i}} U [f(K_{t}) + K_{t} K_{t+1}]$$
(20)

where K is wealth itself, not its logarithm k.  $K_1$  is given;  $\pm < 1$ ; U [¢] is a strictly concave utility function; f (¢), which is the production function in per capita terms, is strictly concave, and the argument of U [¢], f (K<sub>t</sub>) + K<sub>t i</sub> K<sub>t+1</sub> shows consumption at t.

The maximization of (20) with respect to K<sub>t</sub> requires:

$$\pm^{t_{i}} U_{1}[f(K_{t}) + K_{t_{i}} K_{t+1}] ff_{1}(K_{t}) + 1g_{i} \pm^{t_{i}} U_{1}[f(K_{t_{i}}) + K_{t_{i}} K_{t}] = 0$$
(21)

where subscripts, apart from those indicating time, denote di¤erentiation. Rearranging and showing consumption levels explicitly gives:

$$\frac{U_1[C_t]}{U_1[C_{t_i,1}]} = \frac{1}{\pm ff_1(K_t) + 1g}$$
(22)

This is the equivalent to the optimality condition for the continuous case (see Barro and Sala-i-Martin (1995), p.63, equation (2.8)) and Blanchard and Fischer (1989), p.40 equation (7b4).

Taking logarithms of both sides of (22) gives:

$$\ln U_1[C_t]_i \ \ln U_1[C_{t_i 1}] = i \ \ln t_i \ \ln f_1(K_t) + 1g$$
(23)

Or,

$$i C^{m} \frac{U_{11}[C^{m}]}{U_{1}[C^{m}]} i \frac{C_{t} i C_{t}}{C^{m}} = i \ln t \ln f_{1}(K_{t}) + 1g$$
(24)

where  $C^m$  is a value chosen to satisfy the requirement of the mean-value theorem. Denote i  $C^m \frac{U_{11}[C^m]}{U_1[C^m]}$ , the elasticity of marginal utility, by <sup>3</sup>. Equation (24) can be interpreted as follows:

[Elasticity of marginal utility = 3][-Growth rate of consumption]

 $= i \ln \pm i$  In fMarginal product of capitalg (25)

where the growth rate of consumption in (25) is measured on the base of a mean-value theorem level, and  $^{3}$  is evaluated at the same point.

The larger is the value of  $K_t$  the smaller is  $f_1(K_t)$ , then the larger is the right-hand side of (25). For the left-hand side of (25) to be correspondingly larger, either the growth rate of consumption must be smaller, or <sup>3</sup> must be larger. If we follow Barro and Sala-i-Martin (1995, p.64) and assume <sup>3</sup> to be a constant, we arrive at a standard type of result. The growth rate of consumption must be smaller when  $K_t$  is larger. In that particular sense we have <sup>-</sup>-convergence.

To move from this  $\bar{}$ -convergence result to a conclusion concerning how the growth rate of K depends upon the level of K involves some complex calculation. See Barro and Sala-i-Martin (1995), Appendix 2C to Chapter 2, for analysis of the parallel. problem for the continuous case. Here there is no need to pursue the issue further. As will be argued vigorously below, in the next section, the supposition that the elasticity of marginal utility might be a constant is, despite its convenience and popularity, highly dubious. Therefore, if we cannot be sure that consumption will grow more rapidly when K is small, we surely cannot require that the growth rate of K should be smaller when K is large.

## 1.5. Why do the Poor not Save More?

Empirical cross-section studies of economic growth show the hypothesis of <sup>-</sup>-convergence to be well supported for fairly homogenous cross-sections (OECD countries or US States from about 1960) and poorly supported for broader cross sections (all countries in the Summers-Heston data set). In the latter case in particular poor countries have not grown as fast as the <sup>-</sup>-convergence hypothesis would have one suppose. In a similar manner the wealth of poor families in national samples does not typically grow as rapidly as an optimistic convergence view would have one expect. Why is this?

The argument could easily become complex, because so many factors and in‡uences might be involved. To keep the discussion as simple as possible, the focus here will be on the saving rate of the poor, de...ned as the share of saving in total income. When the saving rate is independent of wealth, as in the Stigltiz (1969) model, we obtain <sup>-</sup>-convergence, and

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the wealth of the poor grows more rapidly than that of the rich. Under the approach of Barro and Sala-i-Martin, where growth is determined by consumer optimization, we again arrive at <sup>-</sup>-convergence, as long as <sup>3</sup> is constant. Then these authors, and Barro (1997), allow numerous extra variables, from the share of government expenditure in output to religion, to explain why poor countries do not grow as fast as a simple <sup>-</sup>-convergence model would predict.

Everything here hangs on the idea that the wealth of poor units will grow rapidly because:

- <sup>2</sup> they will enjoy a higher return to saving, having less capital relative to their labour; and
- <sup>2</sup> their utility discount rate will be the same as everyone else's

Given these presumptions and equation (18), the poor must have their consumption growing faster than the rich. Yet the pain of saving when one is poor is not only measured by the utility discount rate and the rate at which consumption grows over time. Inspection of equation (18) reveals that if capital holding is low, when the marginal product of capital will be high; and given the value of  $\pm$ , the growth rate of consumption may yet be low if <sup>3</sup> should be large.

This is not an unreasonable case. The value of <sup>3</sup> could be said to measure the ease with which the intertemporal substitution of consumption takes place. A high value of <sup>3</sup> corresponds to intertemporal substitution being di¢cult. Imagine an agent poor and hungry. He can eat a little less and is guaranteed a high return on any such saving. In comparing his current marginal utility with its future value he is no more subject to myopic discounting than any one else (he shares the value of ± that applies to richer individuals). Yet the issue remains of how fast his consumption has to rise over time to equate the ratio of his discounted future marginal utility to his current marginal utility with the high net return on his saving. That growth rate of consumption may be low if <sup>3</sup> is low.

In this connection the argument of Barro and Sala-i-Martin (1995), which claims correctly that <sup>3</sup> must be constant asymptotically, is irrelevant. As all units in this type of model must converge in the limit to the same level of consumption, <sup>3</sup> must similarly converge. That proposition says nothing about the global constancy of <sup>3</sup>. Without assuming <sup>3</sup> to be invariant, we lose any guarantee of <sup>-</sup>-convergence, and we have to abandon the imposition of simple restrictions on the form of h(k); beyond its root property of being an increasing function.

To conclude this part of the argument, the description of a stationary distribution for wealth needs to take into account wider possibilities for the transfer function than are embraced by a model which gives <sup>-</sup>-convergence. If the transfer function for the logarithm

of wealth should be linear, <sup>-</sup>-convergence follows [see equation (18) above]. However nonlinearities in the transfer function greatly enrich the range of possible stationary wealth distributions and they have to be taken into account. In a similar spirit we may note that non-normality and asymmetry in the generating process for the shocks <sup>2</sup> is a feature which has to be taken seriously, and this in turn may shape a stationary wealth distribution which a "well-behaved" case would not allow. All this will become clearer below.

## 1.6. Convergence and Scattering

The exect on the distribution of wealth in moving from one period to the next is the sum of two separate transformations. First each k value maps to h[k]. This is optimal adjustment without any shock; called h-transfer. Next all values are scattered by the addition of random shocks  ${}^{2}_{t}$ . We call this scattering.

Applying the mathematical expectation operator E to (3) gives:

$$\mathsf{E}\mathsf{k} = \mathsf{E}\mathsf{h}[\mathsf{k}] \tag{26}$$

The relation of Eh[k] to h[Ek] depends upon the concavity/convexity of h[l], which is ambiguous. Subtracting (3) from (25) and rearranging gives:

$$E [k_{t i} Ek]^{2} = E [h[k_{t}]_{i} Eh[k] + {}^{2}_{t}]^{2}$$
  
= E [h[k]]^{a} + E [Eh[k]g]^{b} + E [{}^{2}\_{t}]^{a}\_{i} 2E fh[k] Eh[k]g (27)  
= E [{}^{2}\_{2}]^{a} + E [h[k]]^{2}\_{i} fEh[k]g^{2}

where time subscripts have been dropped, because they are irrelevant when a stationary distribution is under consideration. On account of  $_{t}^{2}$  being i.i.d., expectations of products involving  $_{t}^{2}$  have been equated to zero.

Notice that the variance of h(k) is given by:

$$E fh(k)_{i} Eh(k)g^{2} = E^{(k)}h(k)^{2} + [Eh(k)]^{2}_{i} 2h(k)Eh(k)^{a}$$
$$= E^{(k)}h(k)^{a}_{i} fEh(k)g^{2}$$
(28)

Now equations (27) and (28) together can be interpreted in a very natural result.

Theorem 2. An h<sub>i</sub> transformation always subtracts variance from the distribution of k. For a stationary distribution it subtracts precisely the amount of variance that is added by scattering.

<u>Proof</u>: Notice that the result is not trivial. While  $h_i$  transformation obviously moves every k closer to  $k^{\mu}$ , there is no immediate guarantee that it moves every k closer to the mean of the k values. However from (27) the second moment of the distribution of k in general,

and hence the same moment in a stationary distribution, is the sum of the variances of h(k) and of <sup>2</sup>. In that case, k itself must have a larger variance than h(k). Evidently scattering restores equality, as required.<sup>x</sup>

Similar calculations for the third moment of the distribution of k, assuming E  $f^{23}g = 0$ , produce:

$$E fk_{i} Ekg^{3} = E^{(k)}h(k)^{3}_{i} 3E fh(k)g E^{(k)}h(k)^{2} + 2[E fh(k)g]^{3}$$
 (29)

Theorem 3. If the distribution of shocks  $_{t}^{2}$  has a third moment about its mean equal to zero; hence in particular if it is symmetrical about zero; an h<sub>i</sub> transformation applied to a stationary distribution does not a<sup>x</sup>ect the third moment of k about its mean.

<u>Proof</u>: The right-hand side of (29) is the third moment h(k) about its mean. Therefore the result follows immediately.<sup>x</sup>

## 1.7. The Shape of the Stationary Distribution: Transfer Plus Scattering

In this section, like the last, the joint exects of h-transfer and scattering are taken into account without making simplifying assumptions on either side. Thus confers the bene...t of great generality. The cost is that one is then confronted with the complex product of the two exects, in manner which will be made fully precise. In that case it is not always easy to see how the separate intuences of convergence and scattering axect the shape of a stationary wealth distribution. In later sections simpler cases will be displayed which make the separate exect of one or other of the two intuences more transparent.

In analysing the distribution of k values, it is sometimes convenient to work in terms of the cumulative distribution. Hence  $\mathfrak{C}(k)$  is the proportion of the population with wealth not greater than k. Clearly  $\mathfrak{C}(\mathfrak{i} \ 1) = 0$  and  $\mathfrak{C}(1) = 1$ .

Recall that the exect on the distribution of wealth in moving from one period to the next is the sum of two separate transformations. First each k value maps to h[k]; that is h-transformation. Next all values are scattered by the addition of random shocks  $^{2}_{t}$ . Consider the ...rst step. Before h-transformation the probability density of k, x [k], is given by:

$$\approx [k] = \frac{d\Phi[k]}{dk}$$
(30)

Whereas the cumulative distribution of k after h-transformation, j [k], satis...es:

$$\mathbf{i} [\mathbf{k}] = \mathbf{C} \mathbf{\hat{h}}^{i 1} [\mathbf{k}]^{\mathbf{m}}$$
(31)

Then:

$$\frac{d_{i}[k]}{dk} = \frac{d\mathfrak{C}[h^{i}][k]}{d[h^{i}[k]]} \frac{d[h^{i}[k]]}{dk} = \frac{\mathfrak{m}[h^{i}[k]]}{\frac{dh[k]}{dk}}$$
(32)

is the density of wealth distribution after  $h_i$  transformation. Equation (32) de...nes how the adjustment function a<sup>x</sup>ects the distribution of wealth in the absence of random e<sup>x</sup>ects.

Denote the transformed distribution by  $\mathbb{O}(k)$ . So:

$$^{\mathbb{C}}(\mathbf{k}) = \frac{\mathbb{E}\left[\mathbf{h}^{i} \ {}^{1}\left[\mathbf{k}\right]\right]}{\frac{\mathrm{d}\mathbf{h}\left[\mathbf{k}\right]}{\mathrm{d}\mathbf{k}}} \tag{33}$$

Consider a maximum of  $^{\odot}(k)$  at  $k = k^0$ . Then:

$$\frac{\frac{1}{2}}{\frac{d^{\frac{1}{2}}[h^{i^{-1}}[k]]}{dk}}_{k} = \frac{\mathbf{f}}{h^{i^{-1}}[k]} \frac{\mathbf{f}}{h^{i^{-1}}[k]} \frac{\frac{d^{2}h[k]}{dk^{2}}}{\frac{1}{k^{2}}} \frac{\mathbf{f}}{k^{2}} = 0$$
(34)

Equation (34) is useful when locating a maximum, including a mode, of a wealth distribution after  $h_i$  transformation when the location of a maximum of  $\alpha$  [k] is known. Suppose, for instance, that h[k] is so nearly linear in the relevant range that  $\frac{d^2h[k]}{dk^2}$  may be replaced by zero. Then (34) says that one should look for a maximum of  $\alpha$  (k) to the left (right) of a maximum of  $\alpha$  [k] according as k is less than (greater than) k $\alpha$ .

The sequential exects of h<sub>i</sub> transfer and scattering in a stationary case can be exhibited mathematically as follows. Take any value of k. Suppose k < kx . A symmetrical argument works for the other side. For any level of wealth between h<sup>i 1</sup> [k] and k, h<sub>i</sub> transfer will carry wealth across the border marked by k from lower to higher values of wealth. Next, after h<sub>i</sub> transformation, scattering will carry a certain mass of wealth across the same border, travelling in the same direction, while scattering will also carry another mass of wealth across the border in the opposite direction. It is an evident equilibrium condition for a stationary distribution that the net movement of wealth across the border shall be zero. That condition is expressed in the following equation.

where  $| (\mathfrak{c})$  is the cumulative distribution of i.i.d. shocks; that is the probability that  $_{t}^{2}$  will be  $\cdot$  the argument of  $| (\mathfrak{c})$ .

As (35) holds as an identity in k, we may dimerentiate it with respect to k to obtain:

Notice that if the distribution of shocks is symmetrical about 0, then  $| (0) = \frac{1}{2}$ . However the argument does not use that property. Simplifying (36) gives:

Changing variable in (37) -  $_{s} = h^{i-1}(\cdot)$  - provides another integral equation description of a stationary equilibrium.

$$\mathbf{Z}_{+1} \\ \mathbf{x} [\mathbf{k}] = \int_{\mathbf{i}}^{\mathbf{i}} \mathbf{k} [\mathbf{k}_{\mathbf{i}} \ \mathbf{h} [\mathbf{j}]] \mathbf{k} \mathbf{x} [\mathbf{j}] \mathbf{d} \mathbf{j}$$
(38)

Or, denoting the variable of integration again by  $\cdot$  :

The integral on the right-hand side of (39) is the sum of all transitions from  $\cdot$  to k weighted by the probability that the initial value is  $\cdot$ , which is  $\approx [\cdot]$ , and the probability of a transition to k, which is the probability that  $^{2}_{t}$  takes the value k i h[ $\cdot$ ]. Placing the same function  $\approx [c]$  on both sides of (39) identi...es a stationary ...xed point outcome.

Equation (39) is a Fredholm Equation of the second kind<sup>6</sup>. This derivation is somewhat similar to the so-called Theory of Breakage which leads to the equation:

$$F_{j}(\mathbf{x}) = \bigcup_{u} H_{j} \frac{\mathbf{x}}{u} dF_{j_{i} 1}[u]$$
(40)

for which see Aitchison and Brown (1957), pp.26-7.

The process:

$$k_{t+1} = h[k_t + {}^2_t]$$
(41)

generates another Fredholm Equation, viz:

which is quite similar.

To keep things simple, we concentrate below on the Fredholm Equation (39).

Results from the Fredholm Equation. To be able to write down an equation showing a stationary solution as (39) is encouraging. Unfortunately this equation cannot be solved for x [k]. However it yields two useful results.

Theorem 4. The set of functions satisfying (39) is convex<sup>7</sup>. <u>Proof</u>: Is immediate. If  $x^{1}[k]$  and  $x^{2}[k]$  both satisfy (39), then:

$$\mathbf{Z}_{i+1} = \mathbf{X}_{i} [k] = \mathbf{X}_{i} [k_{i} \ h_{i} [\cdot]] \ \mathbf{x}_{i} \ \mathbf{x}^{i} [\cdot] d.$$
(43)

<sup>&</sup>lt;sup>6</sup>See Hildebrand (1961) p. 381-2. In section 4.5 of the same chapter the author explains the connection between this type of equation and the joint exect of many causes.

<sup>&</sup>lt;sup>7</sup>To say that the set of functions is convex is not, of course, to say that the functions are convex functions.

for i = 1 or 2, and for any value of  $\int$ . Hence:

$$Z_{+1} = \frac{[k] + (1_{i}]) (x^{2} [k]}{[k_{i} h[\cdot]] (x^{2} [k] + (1_{i}]) (x^{2} [k])}$$

$$= \frac{[k_{i} h[\cdot]] (x^{2} (x^{2} [k] + (1_{i}]) (x^{2} [k]) (k^{2} [k]) ($$

¤

The next theorem uses the Fredholm equation to establish continuity of  $\approx [k]$  with respect to k. It assumes that  $\frac{1}{4}$  [¢] is uniformly continuous. For a probability density function this is a mild condition. For note that, because any continuous function is uniformly continuous on a compact support; if there is any problem with uniform continuity of  $\frac{1}{4}$  [¢], it can only arise from extraordinary behaviour of the function in its tails.

Theorem 5. If  $\frac{1}{2}$  [ $\beta$ ] is uniformly continuous, a stationary distribution value for x [k] is continuous in k.

<u>Proof</u>: From (33) it will be seen that for a stationary distribution of k, the continuity of x [k] is implied by the continuity of:

Take a sequence of values  $k_1; k_2; ...; k_n; ...$  with limit k. Then the sequence of values

are given by terms of the form:

which are less than or equal to:

From uniform continuity it follows that for any  $\mu > 0$ , jk i k<sub>n</sub>j su¢ciently small implies that (48) is less than: 7

$$\mu \prod_{i=1}^{\mu} \mathbb{E}\left[\cdot\right] d \cdot = \mu$$
(49)

From which the continuity of *¤* [k] follows. *¤* 

Continuity of x [k] does not rule out the possibility that an stationary density might split into two or more disjoint segments: say a high wealth segment with positive density; a low wealth segment with positive density; and a region between these two where x [k] = 0.

We have encountered such a possibility with the long-memory example of page 7 above. There in the context of long/short memory we saw that in the limit k will be found in one of two (but in general it could be three or more) zones of attraction, and once inside such a zone k will always remain in that zone. It is immediate that this situation can only arise if the distribution of shocks <sup>2</sup> has one or more bounded tails. Otherwise a shock greater than or equal in absolute value to that required to throw the system form one zone to another will occur with probability 1 in the limit, and the system will have short-memory.

Definition 4. The distribution of shocks <sup>2</sup> will be said to be regular if the set of closed intervals over which the probability measure for <sup>2</sup> is positive is a convex set.

The de...nition will surely be satis...ed by any distribution likely to be of interest. Here speci...cally the de...nition excludes the possibility that zero probability will attach to a closed interval of values of <sup>2</sup>, while positive probability attaches to closed intervals for values of <sup>2</sup> both above and below that interval.

Theorem 6. If the distribution of shocks <sup>2</sup> is regular, and

$$k_{t+1} = h[k_t]$$
 (50)

is a short-memory process, a stationary density cannot be disjoint.

<u>Proof</u>: If the distribution is disjoint there will exist at least two open intervals of values of k such that  $\approx [k] > 0$  for values of k in those intervals. Call the said intervals  $k_1$  and  $I^{\circ\circ}$ . Between  $I^{\circ}$  and  $I^{\circ\circ}$  will be found an interval such that  $\approx [k] = 0$  for all values of k in that interval. Call this last interval <sup>100</sup>. Because the process (50) is short-memory, k will transit between  $I^{\circ}$  and  $I^{\circ\circ}$  with probability 1 during any in...nite history. Consider such a transit which will satisfy:

$$k_2 = h[k_1] + {}^2_1 \tag{51}$$

for  $k_1 \ge 1^\circ$  and  $k_2 \ge 1^\circ$ . Then a transit from  $k_1 \ge 1^\circ$  to  $k_3 \ge 1^\circ\circ$  could occur, with  $k_3$  closer to  $k_1$  than is  $k_2$ , if  $^2$  were to take a value closer in absolute value to zero than is  $^2_1$ . As the mean of  $^2$  is zero, the support of  $^2$  must include both positive and negative intervals. Hence because the distribution of shocks is regular, positive probability must attach to transits to  $1^\circ\circ\circ$ . Then there cannot be zero probability density on  $1^\circ\circ\circ$ , not in the limit of any history, hence not in particular for a stationary distribution. a

1.8. The Shape of the Stationary Distribution: Special Cases

The general analysis of Section 1.7 above denies us simple insights into the way in which the stationary distribution is shaped by the separate forces operating on it. For that understanding it is helpful to look at models which are designed to isolate particular

exects, while other intuences which combine to make up the Fredholm integral are muted by simplifying special assumptions.

In the following section the focus is on the form of the transfer function, in particular the exect of non-linearities in the transfer function; so only the most simple speci...cations for the distribution of shocks are admitted. In section after the next the focus is on the intuence of the shocks themselves. In that case a linear transfer function is ideal, and the outcome with that assumption will be examined.

The Shape of the Stationary Distribution: The form of the Transfer Function. To elucidate the stationary distributional properties of a variable generated by the stochastic process (2), the following strange, yet understandable, hydraulic model may be helpful.

A Hydraulic Model. In the centre is a rift valley, running due north-south, and viewed in cross section. Rivers  $\ddagger$  work down from highlands on the east side and from the west. Position is measured by a variable k which runs from  $\ddagger 1$  (inde...nitely far west) to  $\pm 1$  (inde...nitely far east).

These are not normal rivers, fed by springs, and rainfall originating outside the river system. The system is completely closed. All rainfall originates from water in the rivers themselves. Evaporation constantly redistributes water within the system. The amount of water evaporated depends on the volume at a point. One molecule of water may travel any distance, east or west. The probability of any such journey depends upon the absolute distance travelled, and it decreases monotonically with absolute distance. Elevation as such has no exect on precipitation. Indeed the high highlands are dry, because they are far from the great mass of water. Finally water runs down hill and it runs faster the steeper the absolute gradient.

The bottom of the rift valley is at  $k^{\mu}$ . The ‡ow of river water towards the valley represents non-stochastic transformation of values of k through the function h [k]; which is to say that it represents neoclassical convergence. Evaporation and the random redistribution of water represent the e<sup> $\mu$ </sup> ect of i.i.d. shocks, which called scattering as above. The depth of water at any point k represents the density of wealth at that point. When this hydraulic system is in a stationary state, depth is constant at any point. The rivers ‡ow always towards  $k^{\mu}$ . However evaporation and the random redistribution of water frustrate that process. A deep lake may build up around  $k^{\mu}$ . Yet if redistribution is signi...cant, the lake can never contain all the water in the system, because redistribution will always throw some water back into the highlands.

The hydraulic model o¤ers a helpful mental picture of how the type of non-linear stochastic process under consideration in this paper might appear in a stationary equilibrium. It is plain that non-linearities in the slopes of the valley walls will shape the stationary distribution, almost as a potter's hands shape the ...nal pot. More detail on how non-linearities have their e¤ects will follow below. The Shape of the Stationary Distribution: The Intuence of Shocks. Take an arbitrary value of k,  $k_0$ , and follow its random path as it is repeatedly transformed by the process (3). This can be written as:

$$k_{t+1} i k^{\mu} = h[k_t] i h[k^{\mu}] + {}^{2}_{t}$$
 (52)

Or,

$$\boldsymbol{\Re}_{t+1} = \boldsymbol{\Re}_{t} \mathbf{h}^{0} \stackrel{\mathbf{f}}{=} \mathbf{k}_{t}^{M} \stackrel{\mathbf{a}}{=} \mathbf{k}_{t}^{2}$$
(53)

where  $\mathbf{\hat{R}}_t = k_t \mathbf{i} \mathbf{k}^{\mathbf{x}}$ , the prime <sup>1</sup> denotes dimerentiation, and  $k_t^{\mathbf{M}}$  is the value of k which makes (53) correct. The mean-value theorem says that such a value lying between k and km always exists. From (53):

$$\mathbf{\hat{R}}_{t+1} = \mathbf{\hat{R}}_{t_{i} 1} \mathbf{h}^{0} \mathbf{\hat{E}}_{t_{i} 1}^{M} + {}^{2}_{t_{i} 1} \mathbf{h}^{0} \mathbf{\hat{E}}_{t_{i} 1}^{M} + {}^{2}_{t_{i}}$$
(54)

Which simpli...es to:

$$\mathbf{R}_{t+1} = \mathbf{R}_{t_{i}\ 1} \mathbf{k}_{s=t_{i}\ 1}^{t} \mathbf{h}^{0} \mathbf{k}_{s}^{M} \mathbf{k}^{*} + \mathbf{k}_{t_{i}\ 1}^{t} \mathbf{h}^{0} \mathbf{k}_{t}^{M} \mathbf{k}^{*} + \mathbf{k}_{t}^{2}$$
(55)

Similarly, repeated substitutions give:

$$\mathbf{\hat{R}}_{t+1} = \mathbf{\hat{R}}_0 \overset{\dagger}{}_{s=1}^{t} \mathbf{h}^0 \overset{\mathbf{f}}{\mathbf{k}}_s^{\mathsf{M}} \overset{\mathbf{a}}{\mathbf{k}} + \mathbf{S}$$
(56)

where S is equal to:

$${}^{2}_{t} + {}^{2}_{t_{i}} {}^{1}_{1}h^{0} {}^{\mathbf{f}}_{k_{t}} {}^{\mathbf{m}}_{k_{t}} + {}^{2}_{t_{i}} {}^{2}_{2}h^{0} {}^{\mathbf{f}}_{k_{t}} {}^{\mathbf{m}}_{k_{t}} h^{0} {}^{\mathbf{f}}_{k_{i}} {}^{\mathbf{m}}_{1} + :: + {}^{2}_{1} {}^{+}_{q} {}^{2}_{q} {}^{2}_{2}h^{0} {}^{\mathbf{f}}_{k_{q}} {}^{\mathbf{m}}_{q}$$
(57)

If the unshocked system is globally stable, the ...rst term on the right-hand side of (56) will go to zero as t ! 1. A su¢cient condition for that property is:

$$h^{0}[k] \cdot {}^{3} < 1 \text{ all } k$$
 (58)

While (5) ensures that (58) is satis...ed at  $k = k^{\mu}$ , the fact that  $h[\ell]$  may have any concavity/convexity, shown above, implies that is not necessarily satis...ed everywhere.

Given global stability and that the ...rst term of the right-hand side of (56) goes to zero, the limit of (56) as t ! 1 gives the frequency distribution of k. It would be incorrect to conclude that because this distribution is independent of  $k_0$ , the process must be short-memory. So long as  $h[\ell]$  is non-linear, the in‡uence of  $k_0$  may make itself felt via the particular values of  $h^{\ell} k_s^{M}$  that appear in (57). Should  $h[\ell]$  be linear, which possibility cannot be excluded, the limiting distribution of k depends only upon the distribution of the shock values <sup>2</sup> and upon powers of the constant slope coe¢cient.

Equation (57) can be employed to give insight into how non-linearity in the adjustment function h[k] translates into asymmetry in the stationary distribution. Suppose that h[k] is linear to the left of k<sup>a</sup>, and also to the right of k<sup>a</sup>, but the two slopes di<sup>a</sup>er. Let the slope to the left, h<sup>l</sup>, be larger than the slope to the right, h<sup>r</sup>. As negative realizations of the random variable <sup>2</sup> will on average be associated with k < k<sup>a</sup>, and positive realizations of the random variable <sup>2</sup> will on average be associated with k > k<sup>a</sup>, (57) can be read to say that negative values of <sup>2</sup> will be more heavily weighted, and the stationary distribution will have greater density to the left of k<sup>a</sup> than to the right. The fact that convergence to k<sup>a</sup> is more rapid from the right than from the left accounts for this asymmetry.

Figure 2 illustrates this case in which non-linearity of h(k) combines with symmetrical shocks to produce an asymmetrical distribution. If capital is at the point A, and ignoring shocks at this point, it is seen that increase in k is relatively modest. Starting from point B, however, the decline in capital is far larger. These exects will cause the stationary distribution to bunch to the right of kx and to spread to the left of kx, as required.

Equation (57) provides further insight if we transfer attention to another simple case. Now h [k] is linear, but the distribution of <sup>2</sup> values is asymmetric about zero. It is clear that the asymmetry of the distribution ¼ [<sup>2</sup>] is re‡ected in a similar asymmetry in the stationary distribution of k values. In the context of the theory of wealth accumulation this possibility is intriguing. We can always make the expected value of <sup>2</sup> equal zero, as adding or subtracting a constant to the <sup>2</sup> values, and subtracting or adding the same value to h [k] makes no di¤erence. That point does not dispose of the possibility that higher probability density may attach to large negative shocks (bad growth set-backs) than to large positive shocks of similar size. In the type of case just described the stationary distribution will be fat to the left of k¤ relative to its density at a similar distance from k¤ to the right.

Theorem 7. If the stochastic process is standard and h[k] is linear, a stationary distribution is symmetric with its centre at k<sup>x</sup>.

<u>Proof</u>: If h(k) is linear, it follows that  $h^{\circ}(k)$  in (57) will be a constant, denoted H (0 < H < 1), and that this expression will become:

$${}^{2}_{t} + {}^{2}_{t_{i}} {}^{1}_{1}H + {}^{2}_{t_{i}} {}^{2}_{2}H^{2}_{t} + :: + {}^{2}_{1}H^{t_{i}} {}^{1}$$
(59)

In this case the process is short-memory, in the sense that the limiting probability that k lies in any interval is independent of its initial value. Suppose that the limiting probability that k lies in the closed interval of positive values  $[k^i; k^+]$  is  $p_0$ . This is equivalent to:

$$\operatorname{Lim}_{t! \ 1} \sum_{p}^{\mathfrak{C}} {}_{t \ + \ 2_{t_{i} \ 1}} H \ + \ :: \ + \ 2_{1} H^{t_{i} \ 1} \sum_{\mu=1}^{\mathfrak{a}} {}_{\mu} [2_{\mu}] = p_{0}$$
(60)

where the integration over P in (60) is over all values of  ${}^{2}_{\mu}$  such that the integral, and hence its limit, is equal to  $p_{0}$ . In that case it seems that (60) says nothing. It being the case that

the integration on the left-hand side is over all values of <sup>2</sup> such that the integral takes the value  $p_0$ , it conveys no information to state that all the integrals, and hence their limit as t ! 1, take the value  $p_0$ . For our present purposes however, what matters is not that (60) is satis...ed, but rather the manner, exhibited in (60), in which the various values of <sup>2</sup> and their probability densities combine to produce a value  $p_0$ . Consider the closed interval of negative values [ $i k^+$ ;  $i k^i$ ], and the limit of a sequence of integrals:

$$Z_{e} = \frac{\mathbf{z}_{t} + 2}{\mathbf{z}_{t} + 2} + \frac{2}{t_{i}} + \frac{$$

where integration over P<sub>i</sub> is over all values of  ${}^{2}_{\mu}$  such that k lies in the closed interval of negative values [k<sup>i</sup>; k<sup>+</sup>]: By symmetry of the density function  ${}^{4}$  [<sup>2</sup>], (61) takes the value p<sub>0</sub>, which is the result required.¤

# 1.9. The Shape of the Stationary Distribution: Asymmetry and Single-Peakedness

Definition 5. The wealth distribution will be said to be single-peaked if all its local maxima are attained on one convex set of values of k.

The de...nition allows a "table mountain" case in which the maximum value is attained over a connected range of values of k. That case apart, the de...nition rules out multiple local maxima as distinct peaks. In the standard case, a linear h[k] function produces a symmetric stationary distribution. That does not by itself imply single-peakedness, as a symmetric distribution might have many isolated local maxima. However an argument similar to the proof of Theorem 7 shows that if the distribution of <sup>2</sup> values is symmetric around zero with probability density a declining function of the absolute distance from zero, a stationary distribution of k values has the same qualitative form.

Without excluding any case likely to be of interest, we may con...ne attention to distributions of <sup>2</sup> values which are centred on zero and with density monotonically decreasing in the absolute distance from zero, but not necessarily at the same rate for positive and negative values of <sup>2</sup>. Within that class of cases, we may, by collecting implications of results derived above, throw considerable light on the question of asymmetry in a stationary distribution of k values. Asymmetry can come about only from at least one of the following features:

- <sup>2</sup> non-linearity of the adjustment function h[k];
- <sup>2</sup> asymmetry of the density ¼ [<sup>2</sup>].

The hydraulic system described in above can also throw light on single-peakedness. Suppose that the steepness of the rift valley walls on both sides increases monotonically with absolute distance from k<sup>x</sup>, and is symmetrical on the two sides. Far from k<sup>x</sup> water is moved quickly towards  $k^{a}$ . All the water that far out has been transported a long distance. There cannot be much of it, and swift running rivers must be shallow. As one moves closer to the ‡oor of the valley, the absolute gradient becomes lower and rainfall rises, because the total water not too far away increases. Now rivers ‡ow slowly and are deep. Therefore the depth of water rises until it reaches its maximum at  $k^{a}$ . The density of water is symmetric around  $k^{a}$ .

Now modify the model just described. On the west side insert a range of values of k along which the gradient is quite ‡at. Follow it by a very steep interval closer to k<sup>a</sup>, and then return to a similar gradient to that prevailing on the opposite valley wall. The amount of water above these ranges will hardly be a<sup>a</sup>ected if redistribution is strong and the ranges described cover short intervals. Therefore water will move through the intervals ...rst slowly, next rapidly, then it will slow down. Depth will be high, then lower, then high again. The depth of water, which is to say the density of wealth, will exhibit twin peaks.

For the accumulation of wealth the model just described corresponds to the following state of a¤airs. For a range of low levels of wealth, is accumulated towards k¤, but at a slow rate. Then, when wealth gets a bit higher, the pace of accumulation picks up sharply. Later it moderates. If we allow the elasticity of marginal utility to vary with wealth, economic theory cannot exclude such a case. The only way to avoid a twin-peak outcome in such a case is to have a high density in the steep (fast-‡owing) section. That will never be a stationary equilibrium because that high density would be rapidly dissipated by ‡ow towards k¤ which rainfall could not replace.

An informal mathematical version of this pictorial argument runs as follows. Take a regular model with no twin-peaks in a stationary wealth distribution. Over a range of values of  $k < k^{\mu}$  which is small relative to a range which contains much of the density of  $\frac{1}{4}$  [¢], distort the h[¢] function so as to make its derivative large. Figure 3 sketches this case. Recall that with k the logarithm of wealth, the value h(k) i k measures the rate of growth of capital. In the ...gure this growth rate is the height of the thick non-linear curve above the linear 45° line through the origin. It will be seen that starting from the lowest levels of capital the growth rate of capital is successively rapid, slow, rapid slow. Ranges with these respective qualities are shown on the ...gure by the letters R and S. Given a suitable distribution.

An example of this type depends upon the magnitude of  $\frac{dh[k]}{dk}$  varying considerably over a narrow range: ...rst rising then falling. It has been argued above that such severe nonlinearities cannot be assumed away. Within the family of standard Ramsey growth models are to found examples in which growth proceeds slowly for very poor units, then rapidly for medium-income units, and the again slowly closer to k¤. To say that such an example contradicts the constancy of the elasticity of marginal utility is like saying that total sales revenue for a market falling as total sales increase contradicts the assumption that the elasticity of demand for that market is unity. The statement is correct but uninformative.

When Quah<sup>8</sup> published his empirical evidence showing the twin-peak pattern in international cross section per capita income data, I read it as evidence against the simple convergence model, as no doubt did many other readers. It is interesting to note that Quah himself advances no such claim. First he is very clear that he is describing the development of income distribution over a short period of time. Secondly, Quah is aware of the possibility that even the apparently disconnected distributions he observes may be generated by a process which in the long run is ergodic. Now the theoretical investigation of a stationary wealth distribution has shown that it may have twin (indeed multiple) local peaks. So it seems that even such a surprising feature may be completely consistent with a standard convergence model.

That is not a good way of looking at matters. To be worthy of study, an economic model has not only to be true in some high abstract sense; it has to be useful. A twin peak case can only arise when h[k] is severely non-linear. The estimation of such a model presents many di¢culties. In a way the strength of the Baumol-Barro convergence model is its crude simplicity. If it has to be rescued by re...ned mathematical argument, it loses its appeal. Also, in the example, and more generally, twin peaks in a stationary distribution can only happen over a range within the reach of a single-period realization of the random shock. Therefore if twin peaks are an important feature of the distribution, it must be the case that shocks are large in absolute value. This is another way of saying that the explanatory power of the model is weak.

# 1.10. Concluding Remarks

The long history of the analysis of income or wealth distributions, going back to Pareto, includes di¤erent approaches. One is purely empirical. The shape of the distribution is examined and the …tness of simple mathematical speci...cations is investigated. Another approach is to start with postulates concerning the process which generates the distribution and then to investigate mathematically what is the limiting distribution which results. Yet the limiting distribution does not have to be the object of concern. The shorter term conditional transfer process can itself be the focus of investigation. Indeed for some neoclassical convergence theorists that is all that can be done, because for them the limiting distribution is trivial, being a state in which all countries - or individuals in the case of a personal distribution - are at the common limit point k¤. When the adjustment process is taken to include random e¤ects there are wider possibilities than when it is modelled using non-stochastic economic theory.

The present paper marries two di¤erent traditions. They are the pure neoclassical approach, according to which wealth accumulation is systematic and deliberate; and the

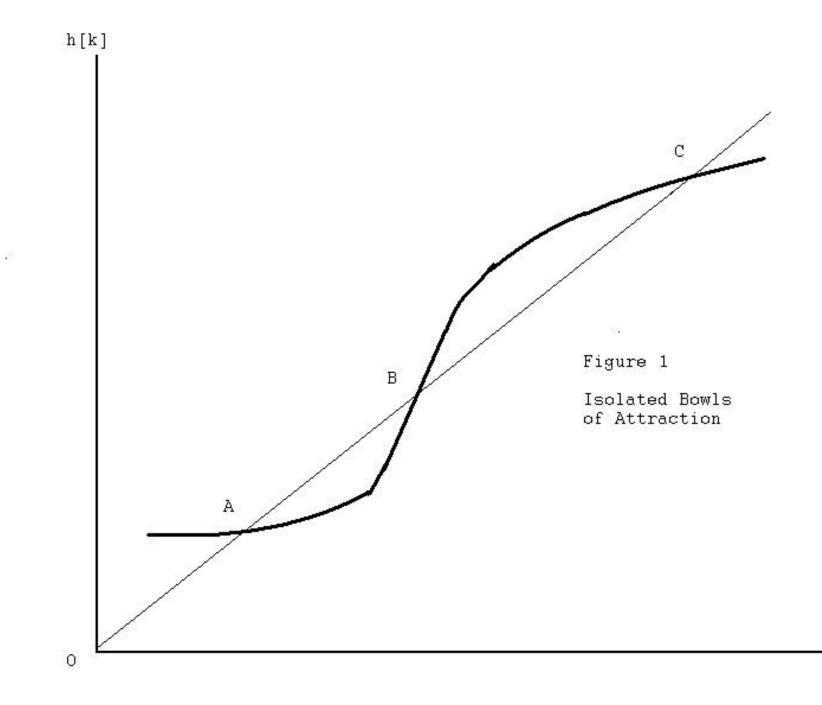
<sup>&</sup>lt;sup>8</sup>See Quah (1996a) and (1996b).

random shocks approach, according to which wealth accumulation is purely haphazard. As would be expected, such a model is complicated, and direct mathematical solution is hardly possible. Even so, we have been able to obtain a series of results which together reveal many features of a stationary distribution of wealth levels.

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