

ON TRANSIENT MARKOV PROCESSES WITH A COUNTABLE  
NUMBER OF STATES AND STATIONARY  
TRANSITION PROBABILITIES<sup>1</sup>

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**1. Summary.** We consider a Markov process  $x_0, x_1, \dots$  with a countable set  $S$  of states and stationary transition probabilities  $p(t | s) = P\{x_{n+1} = t | x_n = s\}$ . Call a set  $C$  of states *almost closed* if (a)  $P\{x_n \in C \text{ for an infinite number of } n\} > 0$  and (b)  $x_n \in C$  infinitely often implies  $x_n \in C$  for all sufficiently large  $n$ , with probability one. It is shown that there is a set  $(C_1, C_2, \dots)$  essentially unique, of disjoint almost closed sets such that (a) all except at most one of the  $C_i$  are atomic, that is,  $C_i$  does not contain two disjoint almost closed subsets, (b) the non-atomic  $C_i$ , if present, contains no atomic subsets, (c) the process is certain to enter and remain in some set  $C_i$ . A relation between the sets  $C_i$  and the bounded solutions of the system of equations

$$(1) \quad \alpha(s) = \sum_t \alpha(t)p(t | s)$$

is obtained; in particular there is only one atomic  $C_i$  and no non-atomic  $C_i$  if and only if the only bounded solution of (1) is  $\alpha(t) = \text{constant}$ . This condition is shown to hold if the process is the sum of independent identical (numerical or vector) variables; whence, for such a process, the probability of entering a set  $J$  infinitely often is zero or one. The results are new only if the process has transient components. The main tool is the martingale convergence theorem.

**2. The structure theorem.**

**THEOREM 1.** *Let  $x_0, x_1, \dots$  be a Markov process with a countable set  $S$  of states (we restrict  $S$  to those states with a positive probability of being entered) and stationary transition probabilities. For any subset  $I$  of  $S$ , denote by  $L(I), U(I)$  the events  $\liminf \{x_n \in I\}, \limsup \{x_n \in I\}$  respectively, by  $\mathfrak{N}$  the class of  $I$  with  $P(U(I)) = 0$ , and by  $\mathfrak{C}$  the class of  $I$  with  $L(I) = U(I)$  a.e. If  $C \in \mathfrak{C}$  and  $C \notin \mathfrak{N}$ ,  $C$  will be called almost closed.*

(1) *Call a Borel measurable function  $f$  on the space  $\Omega$  of all infinite sequences  $\omega = (x_0, x_1, \dots)$ ,  $x_n \in S$ , invariant if for every  $\omega$ ,  $f(\omega) = f(T\omega)$ , where  $T(x_0, x_1, \dots) = (x_1, x_2, \dots)$ , and call an event invariant if its characteristic function is invariant (so that, for any  $J \subset S$ ,  $L(J)$  and  $U(J)$  are invariant). For any invariant event  $V$  there is a  $C \in \mathfrak{C}$  such that  $U(C) = V$  a.e., so that the Borel field of invariant events is identical, up to events of probability zero, with the (Borel) field  $\mathfrak{D}$  of events of the form  $D = U(J)$  a.e.,  $J \subset S$ .*

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(2) There is a finite or countable collection  $\{C_1, C_2, \dots\}$  of disjoint almost closed sets with the following properties:

(a) every  $C_i$  except at most one is atomic, that is, does not contain two disjoint almost closed subsets,

(b) the non-atomic  $C_i$ , if present, contains no atomic subsets,

(c)  $\sum_i P(L(C_i)) = 1$ .

The collection  $\{C_1, C_2, \dots\}$  is essentially unique, that is, if  $\{C'_1, C'_2, \dots\}$  has properties (a), (b), (c), then each  $C'_i$  differs from some  $C_j$  by a set in  $\mathfrak{M}$ .

The  $C_i$  may be chosen so that all states in a given  $C_i$  are of the same type: either all states  $s$  in  $C_i$  are transient, that is,  $\{s\} \in \mathfrak{M}$  or all states in  $C_i$  are recurrent, that is, nontransient. The non-atomic  $C_i$ , if present, consists entirely of transient states.

PROOF. We represent  $x_0, x_1, \dots$ , as usual, as coordinate variables on the space  $\Omega$ . For any bounded measurable  $f$  on  $\Omega$ , we have

$$(4) \quad E(T^n f | x_0 = s_0, \dots, x_n = s_n) = E(f | x_0 = s_n),$$

where  $T^n f$  denotes the function  $g$  defined by  $g(\omega) = f(T^n \omega)$ . For, since  $T^n f$  depends only on  $x_n, x_{n+1}, \dots$ , the left side of (4) equals  $E(T^n f | x_n = s_n)$  which, because of the stationarity of transition probabilities, is easily shown to equal  $E(f | x_0 = s_n)$ . Thus if  $f$  is invariant,

$$(5) \quad E(f | x_0 = s_0, \dots, x_n = s_n) = E(f | x_0 = s_n).$$

For any invariant event  $V$  and any state  $s$ , define  $\phi(s) = E(v | x_0 = s)$ , where  $v$  is the characteristic function of  $V$ . From the forward martingale convergence theorem ([2], p. 319) and (5),  $\phi(x_n) \rightarrow v$  with probability one as  $n \rightarrow \infty$ . Thus if  $I$  is the set of all states  $s$  with  $\phi(s) > \frac{1}{2}$ ,  $\omega \in V$  implies  $\omega \in L(I)$  with probability one, while  $\omega \notin V$  implies  $\omega \notin U(I)$  with probability one, that is,  $U(I) \subset V \subset L(I)$  a.e. Since always  $U(I) \supset L(I)$ , we have  $I \in \mathcal{C}$  and  $U(I) = V$  a.e. This establishes part (1) of the theorem.

For part (2), we decompose the measure  $P$  on the Borel field of invariant sets into atoms and a completely non-atomic part (see for instance [1], p. 565). Let  $V_1, V_2, \dots$  be the sets of this decomposition, and choose  $I_n \in \mathcal{C}$  such that  $U(I_n) = V_n$  a.e. Since  $V_i, V_j$  are disjoint for  $i \neq j$ ,  $I_i \cap I_j \in \mathfrak{M}$  for  $i \neq j$ , so that, with  $C_n = I_n - \bigcup_{j < n} I_j$ ,  $C_n \in \mathcal{C}$   $U(C_n) = V_n$  a.e., and  $C_1, C_2, \dots$  are disjoint. Properties (a), (b), (c) of part (2) and the essential uniqueness are immediate. The final assertion of part (1) is a consequence of the known ([3], p. 322) facts that if  $\{s\} \notin \mathfrak{M}$ , that is, if it is possible for the process to enter  $s$  infinitely often, then if the system ever enters  $s$  it is certain to enter infinitely often  $s$  and all states which can be entered from  $s$ . The latter class  $C$  of states is almost closed (in fact is closed, that is, if the system ever enters  $C$  it remains in  $C$ ), consists entirely of recurrent states, and has no almost closed proper subsets. This completes the proof of the theorem.

Any collection of sets  $\{C_1, C_2, \dots\}$  of the form described in part (2) of Theorem 1 will be called a *decomposition* of the Markov process. A process will

be called *simple* if its decomposition consists of a single set  $C_1$ , and a simple process will be called non-atomic or atomic according to the type of  $C_1$ .

**3. The equation**  $\alpha(s) = \sum \alpha(t)p(t|s)$ . For a process of the type considered in Section 1, write  $p(t|s) = P\{x_{n+1} = t | x_n = s\}$ , so that  $p(t|s) \geq 0$ ,  $\sum_t p(t|s) = 1$  for all  $s$ . The structure of the process is closely related to the nature of the bounded solutions of the equation

$$(6) \quad \alpha(s) = \sum_t \alpha(t)p(t|s).$$

This relation is most simply stated in terms of the bounded invariant functions of the process, as follows:

**THEOREM 2.** *For any bounded invariant function  $f$ ,  $\alpha(s) = E(f | x_0 = s)$  satisfies (6), and every bounded solution of (6) may be represented in this form.*

**PROOF.** That any  $\alpha(s) = E(f | x_0 = s)$  satisfies (6) is clear, since  $E(f | x_0 = s, x_1 = t) = \alpha(t)$  from (5), and (6) then follows from the formula  $E(f | x_0) = E(E(f | x_0, x_1) | x_0)$ . Conversely if  $\alpha(t)$  is any bounded solution of (6), the sequence  $z_n = \alpha(x_n)$  is a bounded martingale and hence converges to a limiting bounded  $f$ . Since  $z_n(T\omega) = z_{n+1}(\omega)$ ,  $f(T\omega) = (f\omega)$  a.e., that is,  $f$  is invariant. The martingale convergence theorem also yields  $E(f | x_0) = z_0 = \alpha(x_0)$ , so that the solution  $\alpha(t)$  has the required form.

The inequality

$$(6') \quad \alpha(s) \geq \sum_t \alpha(t)p(t|s)$$

has been studied by Kendall [7] and Foster [4], who related the existence of solutions of (6') such that  $\alpha(j) \rightarrow \infty$  as  $j \rightarrow \infty$  (enumerating the states by the positive integers) to the existence of a finite closed set of states.

**COROLLARY.** *The process is simple and atomic if and only if the only bounded solution of (6) is  $\alpha(t) = \text{constant}$ .*

**PROOF.** If the only bounded solution of (6) is constant, then any bounded invariant  $f$  is constant, since, with  $\alpha(s) = E(f | x_0 = s)$ ,  $\alpha(x_n) \rightarrow f$  a.e. as  $n \rightarrow \infty$ . Thus any invariant set has probability zero or one, and the process is simple and atomic. Conversely, if the process is simple and atomic, every invariant set has probability zero or one, every bounded invariant function  $f$ , being measurable with respect to the class of invariant sets, is constant a.e., so that every solution of (6), having the form  $E(f | x_0 = s)$ , is constant.

As an application of the corollary, we consider processes which are sums of independent identically distributed variables.

**THEOREM 3.** *Let  $V$  be a finite or countable set of vectors  $v_1, v_2, \dots$  in  $N$ -space, and let  $p_1, p_2, \dots$  be positive numbers with sum one. Let  $S$  consist of the origin and all vectors representable as  $w_1 + \dots + w_n$ ,  $n = 1, 2, \dots$ ,  $w_j \in V$ . The only bounded solution of the equation*

$$(7) \quad \alpha(S) = \sum_j \alpha(s + v_j)p_j, \quad s \in S$$

is  $\alpha(s) = \text{constant}$ . Consequently if  $y_1, y_2, \dots$  are independent variables with  $P\{y_n = v_j\} = p_j$ , the Markov process  $x_k = y_1 + \dots + y_k, x_0 = (0, \dots, 0)$  is simple and atomic.

PROOF. Repeated use of (7) yields

$$\begin{aligned} \alpha(s) &= \sum_{j_1, \dots, j_k} \alpha(s + v_{j_1} + \dots + v_{j_k}) p_{j_1} \dots p_{j_k} \\ (8) \qquad &= \sum_{r \in R_k} \alpha(s + r_1 v_1 + r_2 v_2 + \dots) q_k(r), \end{aligned}$$

where  $R_k$  consists of all sequences  $r = (r_1, r_2, \dots)$  of nonnegative integers whose sum is  $k$ , and

$$q_k(r) = k! \prod_j (p_j^{r_j} / r_j!).$$

Replacing  $s$  by  $s + v_1$  and  $k$  by  $k - 1$  in (8) yields

$$\begin{aligned} \alpha(s + v_1) &= \sum_{r \in R_{k-1}} \alpha(s + v_1 + r_1 v_1 + \dots) q_{k-1}(r) \\ (9) \qquad &= \sum_{r \in R_k} \alpha(s + r_1 v_1 + r_2 v_2 + \dots) q_k(r) (r_1 / k p_1). \end{aligned}$$

Subtracting (9) from (8) yields

$$(10) \qquad \alpha(s) - \alpha(s + v_1) = \sum_{r \in R_k} f(r) q_k(r) (1 - (r_1 / k p_1)),$$

where  $f(r)$  is for fixed  $s$  and  $k$  a bounded function of  $r$  and is uniformly bounded in  $s, k$  as well, say  $|f(r)| \leq M$  for all  $s, k, r$ . For fixed  $\epsilon > 0$ , let  $T_k$  denote that subset of  $R_k$  for which  $|1 - (r_1 / k p_1)| < \epsilon$ . Then

$$\sum_{r \notin T_k} |f(r)| q_k(r) |1 - (r_1 / k p_1)| \leq M p_1 \sum_{r \notin T_k} q_k(r),$$

and the sum on the right, being the probability that, in  $k$  independent trials with an event of probability  $p_1$ , the actual success ratio differs from  $p_1$  by at least  $\epsilon p_1$ , approaches zero as  $k \rightarrow \infty$  by the law of large numbers. Since

$$\sum_{r \in T_k} |f(r)| q_k(r) |1 - (r_1 / k p_1)| \leq M \epsilon,$$

we find from (10), summing separately for  $r \in T_k$  and  $r \notin T_k$  and letting  $k \rightarrow \infty$ , that  $|\alpha(s) - \alpha(s + v_1)| \leq M \epsilon$ . Since  $\epsilon$  is arbitrary,  $\alpha(s) = \alpha(s + v_1)$ . Clearly the same proof yields  $\alpha(s) = \alpha(s + v_j)$  for any  $j$ , and  $\alpha(s) = \text{constant}$ .

COROLLARY. If  $y_1, y_2, \dots$  are independent identically distributed (vector or scalar) variables with a finite or countable set of values, for any set  $J$  the probability that an infinite number of sums  $y_1 + \dots + y_k = x_k$  are in  $J$  is zero or one.

An example of a simple nonatomic process is  $x_n = \frac{1}{2} + \sum_{y=1}^n \epsilon_y / 2^{y+1}$ , where  $\epsilon_1, \epsilon_2, \dots$  are independent and assume the values  $\pm 1$  with probability  $\frac{1}{2}$  each. The states are all rational numbers  $m/2^n, m = 1, 3, \dots, 2^n - 1, n = 1, 2, \dots$ . Writing  $m/2^n = (m, n)$ , the system (7) becomes

$$(11) \qquad \alpha(m, n) = \frac{1}{2}(\alpha(2m - 1, n + 1) + \alpha(2m + 1, n + 1)).$$

If  $x^* = \lim x_n$ ,  $x^*$  is a.e. invariant and has a uniform distribution on  $(0, 1)$ . Each  $x_n$  is a function of  $x^*$  a.e., so that every function of  $x_0, x_1, x_2, \dots$  is a.e. a function of  $x^*$  and hence invariant a.e. For any bounded function  $f(x^*)$ ,

$$E(f | x_0 = (m, n)) = 2^{n-1} \int_{(m-1)/2^n}^{(m+1)/2^n} f(t) dt = \alpha(m, n).$$

It is easily verified by substitution that  $\alpha(m, n)$  satisfies (11).

The Corollary actually holds without the restriction to random variables with a finite or countable set of values; this is an immediate consequence of an unpublished theorem of Edwin Hewitt and L. J. Savage, which they communicated to the writer. The Hewitt-Savage Theorem, an improved version of the zero-one law, asserts that any event depending on a sequence of independent identically distributed random variables which is invariant under all permutations of every finite set of the random variables has probability zero or one; the event that an infinite number of sums  $y_1 + \dots + y_k$  are in  $J$  is clearly of this type. The conclusion of the Corollary, under different hypotheses on the  $y_i$ , has also been obtained by Chung and Derman in an unpublished manuscript which they communicated to the writer.

As a second application of the Corollary of Theorem 2, we obtain an interesting result of Foster [4], Harris [5], and Hodges and Rosenblatt [6] concerning the random walk on the nonnegative integers, with  $p(0 | 0) = 1$ ,  $p(i + 1 | i) = p_i$ ,  $p(i - 1 | i) = q_i = 1 - p_i$ ,  $0 < p_i < 1$ ,  $i > 0$ . The equation (6) becomes

$$(12) \quad \alpha(s) = p_s \alpha(s + 1) + q_s \alpha(s - 1), \quad s > 0.$$

The general solution of (12) is  $\alpha(s) = A + Bz_s$ , where  $z_0 = 0$ ,  $z_1 = 1$ ,  $z_s = c_1 + \dots + c_{s-1}$  for  $s > 1$ , where  $c_1 = 1/p_1$ ,  $c_j = (q_1 \dots q_j / p_1 \dots p_j)$  for  $j > 1$ .

Thus (12) has a bounded nonconstant solution if and only if the series

$$\sum_j (q_1 \dots q_j) / (p_1 \dots p_j)$$

converges, which is the condition obtained in [4], [5], and [6] for passage to the origin to be uncertain.

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