Foundations of Information Theory

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Preface

The intention of the author in writing this book is to provide a concise but rigorous exposition of the fundamentals of the mathematical theory of information. Since the basic work of C. E. Shannon appeared in 1948, a good deal of work has been done in this field. Nevertheless there exists at present no single reference to which one interested in this subject can turn for a presentation which is up to date and yet reasonably complete. It is hoped that the present work may help to fill this gap.

At the present stage of its development, information theory can be divided into two fairly distinct branches. The first of these deals with the general properties of channels of various types, proving various theorems concerning their suitability for transmitting information, while the second deals with the actual implementation of these general theorems. At present, the former is somewhat further developed than the latter. For example, while it is known that a noisy channel with nonzero capacity can be used to transmit information at a rate arbitrarily close to its capacity and with an arbitrarily small probability of reception error, no constructive method yet exists for achieving this, even in the simplest nontrivial channel. There are, nonetheless, a number of interesting and ingenious papers

devoted to the construction of specific codes in channels of various types. However, a detailed treatment of these papers would have increased the length of this book considerably beyond that contemplated. The closest approach to their subject matter occurs in Section 7.3.

The main tool used in Chapters 1 to 4 and 7 is the probability theory of finite systems, as treated, for example, by Feller [1]. The mathematical prerequisites for Chapters 5 and 6 are somewhat more stringent. What is needed is the notion of a general probability space; the mathematical apparatus required for this is the theory of measure and the general Lebesgue integral. While we have included a brief section which sketches those parts of the theory needed, for all details the reader must be referred to one of the several excellent treatises on the subject. Chapter 7 is entirely independent of Chapters 5 and 6, however, and may be read, if desired, following Section 4.1.

The Remarks consist of various material intended to complement the discussion in the text. Most of the references are given here, as well as the details involved in the more flagrant cases of results which are referred to in the text as "easily demonstrated." For this reason, they should be read along with, if not in advance of, the text.

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CHAPTER ONE

Introductory Concepts

The theory of information which we shall develop in this book originated in the work of C. E. Shannon in 1948. In his fundamental paper, Shannon set up a mathematical scheme in which the concepts of the production and the transmission of information could be defined quantitatively. He then formulated and proved a number of very general results which showed the importance and usefulness of these definitions. Since 1948 a number of papers have been published which simplify and extend Shannon's original work. At the same time, various others have attempted to apply information theory to physics, chemistry, and various branches of biology and psychology. These applications will not be touched upon here; in any case, there exist several books, dealing specifically with them, to which the interested reader may be referred.

Our subject matter is mainly deductive in nature, i.e., it is possible to start with a small number of definitions and derive everything else from them, and indeed this is the path we shall essentially take. At the same time, the theory deals with terms such as information content, information source, rate of transmission of information, i.e., terms which carry a certain amount of intuitive meaning. It will therefore be of interest, at each

step of the mathematical development, to compare our results with the dictates of intuition. We shall see that the basic concepts of the theory are readily interpreted in terms of intuitive notions. Of course we cannot expect the same of the more advanced results of the theory. These are basically limit theorems of an involved nature, which may serve as a guide to the intuition of anyone wishing to delve further into the field.

Of the difficulties which confront us when we attempt to construct a quantitative theory in which the concepts "production of information" and "transmission of information" are meaningful, two stand cut at once. First, we must construct a mathematical model in which we can speak of information being produced and transmitted. Second, we must assign a quantitative measure to the amount of information involved. At first glance, it might appear that the solution to the second problem would follow directly from that of the first. That this is actually not the case will be apparent from further consideration of the question.

Intuitively, we would agree that we receive information whenever we are informed of an event whose occurrence was previously not certain. Furthermore it is reasonable that, within certain limits, at least, the more likely an event is, the less information is conveyed us by the knowledge of its actual occurrence. Ignoring for the moment this last remark, we can already introduce a certain amount of formalism into the discussion. Let x represent an event (i.e., its occurrence) and x' its complement (i.e., its nonoccurrence), and let p_x and $p_{x'}$ denote the probabilities of these two events, so that $p_x + p_{x'} = 1$. Let I_x denote the amount of information conveyed to us by the knowledge of the occurrence of x. Since x is specified only by its probability px, we assume I_x to be a function of p_x , i.e., $I_x = I(p_x)$, where I() is a non-negative function defined on the range of values of p_z , namely $0 < p_x \le 1$, the value $p_x = 0$ being meaningless in the present discussion. Similarly we put $I_{x'} = I(p_{x'})$. Since the

probability of receiving the amount I_x of information is p_x , and that of receiving $I_{x'}$ is $p_{x'}$, the average (or expected) amount of information received is given by $p_x I(p_x) + p_{x'} I(p_{x'})$. Similarly, if we have a set x_1, \ldots, x_n of mutually exclusive events such that $p_{x_1} + \cdots + p_{x_n} = 1$, then it is reasonable to consider $p_{x_1} I(p_{x_1}) + \cdots + p_{x_n} I(p_{x_n})$ as the average amount of information conveyed by the knowledge of which x_i actually occurred. If $p_{x_i} = 0$, evidently the corresponding term should simply be omitted from consideration. The choice of the function $I(\cdot)$ is as yet in no way indicated.

It is, however, possible by continuing this type of reasoning to obtain strong constraints upon the form of I(). Consider three mutually exclusive events x, y, z, such that $p_x + p_y + p_z =$ Then $H(x,y,z) = p_x I(p_x) + p_y I(p_y) + p_z I(p_z)$ represents the average amount of information conveyed by the knowledge of which among x, y, z actually occurred. Now, to determine which among x, y, z actually occurred, it is, for example, sufficient to determine whether or not x occurred, and in the event that x did not occur, to determine which of y, z did occur. The amount of information conveyed by the first determination is evidently given by $p_x I(p_x) + (1 - p_x) I(1 - p_x)$, which we may denote by If x did not occur, then the (conditional) probabilities of y and z are given by $\frac{p_y}{p_{x'}}$ and $\frac{p_z}{p_{x'}}$ respectively. The amount of information conveyed by the second determination is therefore given by $\frac{p_y}{p_{x'}} I\left(\frac{p_y}{p_{x'}}\right) + \frac{p_z}{p_{x'}} I\left(\frac{p_z}{p_{x'}}\right)$, which we may denote by H(y,z|x'). But this latter amount of information is conveyed only when the event x' occurs, and so we may take H(x,x') + $p_{x'}H(y,z|x')$ as the total amount of information conveyed, on the average, by both determinations. It follows that a reasonable requirement on the function $I(\cdot)$ is that the relation H(x,x') + $p_{x'}H(y,z|x') = H(x,y,z)$ be satisfied for all allowable values of pz, py, pz greater than zero.

Let us express the H-functions in terms of p_x , p_y , p_z ; i.e., we put $H(x,x') = H(p_x, 1 - p_x)$, $H(x,y,z) = H(p_x,p_y,p_z)$, and $H(y,z|x') = H\left(\frac{p_y}{1-p_x},\frac{p_z}{1-p_z}\right)$. We require, then, that

$$H(p_x, p_y, p_z) = H(p_x, 1 - p_x) + (1 - p_x)H\left(\frac{p_y}{1 - p_x}, \frac{p_z}{1 - p_z}\right)$$

Now the identical reasoning which led to this relation applies if x is considered to be a composite event; i.e., we may consider x to consist of the mutually exclusive events u_1, \ldots, u_{n-1} whose probabilities we denote by p_1, \ldots, p_{n-1} . We then obtain $H(u_1, \ldots, u_{n-1}, y, z) = H(u_1, \ldots, u_{n-1}, x') + p_{x'}H(y, z|x')$, or, letting $q_1 = p_y$, $q_2 = p_z$, and $p_n = p_{x'}$, we obtain

$$H(p_1, \ldots, p_{n-1}, q_1, q_2) = H(p_1, \ldots, p_n) + p_n H\left(\frac{q_1}{p_n}, \frac{q_2}{p_n}\right)$$

This is actually a very strong condition; in fact, as we shall shortly see, it practically suffices to determine the form of $H(p_1, \ldots, p_n)$ without regard for the definition of H in terms of $I(\cdot)$. However, one further condition is suggested by the fact that terms $p_iI(p_i)$ were to be dropped when $p_i=0$, namely that $H(p_1, \ldots, p_n)$ be defined even when some of the p_i vanish, but that it be continuous in the domain defined by $p_i \geq 0$, $i=1,\ldots,n$, $p_1+\cdots+p_n=1$. We can now state:

Theorem. The following three conditions determine the function $H(p_1, \ldots, p_n)$ up to a multiplicative constant, whose value serves only to determine the size of the unit of information.

- 1. H(p, 1-p) is a continuous function of p for $0 \le p \le 1$.
- 2. $H(p_1, \ldots, p_n)$ is a symmetric function of all its variables.
- 3. If $p_n = q_1 + q_2 > 0$, then

$$H(p_1, \ldots, p_{n-1}, q_1, q_2) = H(p_1, \ldots, p_n) + p_n H\left(\frac{q_1}{p_n}, \frac{q_2}{p_n}\right)$$

It is understood throughout that $H(p_1, \ldots, p_n)$ is defined only for a complete set of probabilities, i.e., a set of non-negative numbers whose sum equals one.

The proof is carried out in a series of lemmas.

Lemma 1. We have H(1,0) = 0.

Proof. Using condition 3, $H(\frac{1}{2},\frac{1}{2},0) = H(\frac{1}{2},\frac{1}{2}) + \frac{1}{2}H(1,0)$. But, using condition 2 and then 3, we have: $H(\frac{1}{2},\frac{1}{2},0) = H(0,\frac{1}{2},\frac{1}{2}) = H(0,1) + H(\frac{1}{2},\frac{1}{2})$, which implies H(1,0) = 0, since H(1,0) = H(0,1).

Lemma 2. We have $H(p_1, \ldots, p_n, 0) = H(p_1, \ldots, p_n)$. Proof. By 2 we can assume that $p_n > 0$, and the result then follows by using 3 and then applying Lemma 1.

Lemma 3. We have

$$H(p_1, \ldots, p_{n-1}, q_1, \ldots, q_m) = H(p_1, \ldots, p_n) +$$

 $p_n H\left(rac{q_1}{p_n}, \cdots, rac{q_m}{p_n}
ight)$

where $p_n = q_1 + \cdots + q_m > 0$.

Proof. For m=2 this is precisely condition 3. We proceed by induction on m; from Lemma 2 it is clear that we need consider only the case when $q_i > 0$, $i = 1, \ldots, m$. Suppose that there is an m such that the assertion is true for all n. Then using condition 3, we have

$$H(p_{1}, \ldots, p_{n-1}, q_{1}, \ldots, q_{m+1})$$

$$= H(p_{1}, \ldots, p_{n-1}, q_{1}, p') + p'H\left(\frac{q_{2}}{p'}, \ldots, \frac{q_{m+1}}{p'}\right)$$

$$= H(p_{1}, \ldots, p_{n}) + p_{n}H\left(\frac{q_{1}}{p_{n}}, \frac{p'}{p_{n}}\right) + p'H\left(\frac{q_{2}}{p'}, \ldots, \frac{q_{m+1}}{p'}\right)$$

where $p' = q_2 + \cdots + q_{m+1}$.

But further,

$$H\left(\frac{q_1}{p_n}, \cdots, \frac{q_{m+1}}{p_n}\right) = H\left(\frac{q_1}{p_n}, \frac{p'}{p_n}\right) + \frac{p'}{p_n} H\left(\frac{q_2}{p'}, \cdots, \frac{q_{m+1}}{p'}\right)$$

Substituting this into the preceding equation, we obtain the assertion of the lemma for m+1. Note that the induction

must proceed first along m for n = 2, and then along m for general n.

Lemma 4. We have $H(q_{11}, \ldots, q_{1m_1}; \ldots; q_{n1}, \ldots, q_{nm_n}) = H(p_1, \ldots, p_n) + \sum_{i=1}^n p_i H\left(\frac{q_{i1}}{p_i}, \ldots, \frac{q_{im_i}}{p_i}\right)$, where $p_i = q_{i1} + \cdots + q_{im_i} > 0$.

Proof. Using Lemma 3, we have

$$H(q_{11}, \ldots, q_{1m_1}; \ldots; q_{n1}, \ldots, q_{nm_n}) = p_n H\left(\frac{q_{n1}}{p_n}, \cdots, \frac{q_{nm_n}}{p_n}\right) + H(q_{11}, \ldots, q_{1m_1}; \ldots; q_{n-1,1}, \ldots, q_{n-1,m_{n-1}}; p_n)$$

Shifting p_n to the extreme left, we continue with the reduction, finally obtaining the desired result after n steps.

To continue, let us put $F(n) = H\left(\frac{1}{n}, \dots, \frac{1}{n}\right)$ for $n \geq 2$, and F(1) = 0. Applying the preceding lemma to the case $m_1 = \dots = m_n = m$, $q_{ij} = \frac{1}{mn}$, we obtain F(mn) = F(m) + F(n). For either m = 1 or n = 1, this relation is trivially satisfied. Further, applying Lemma 3 to $H\left(\frac{1}{n}, \dots, \frac{1}{n}\right)$, we obtain

$$H\left(\frac{1}{n}, \dots, \frac{1}{n}\right) = H\left(\frac{1}{n}, \frac{n-1}{n}\right) + \frac{n-1}{n}H\left(\frac{1}{n-1}, \dots, \frac{1}{n-1}\right)$$

from which follows

$$\eta_n \equiv H\left(\frac{1}{n}, \frac{n-1}{n}\right) = F(n) - \frac{n-1}{n}F(n-1)$$

We now show:

Lemma 5. As
$$n \to \infty$$
, $\mu_n \equiv \frac{F(n)}{n} \to 0$, and also $\lambda_n \equiv F(n) - F(n-1) \to 0$.

Proof. From the continuity of H(p, 1-p) follows $\eta_n \to H(0,1) = 0$ as $n \to \infty$. Further,

$$n\eta_n = nF(n) - (n-1)F(n-1)$$

from which follows $nF(n) = \sum_{k=1}^{n} k\eta_k$, or

$$\frac{F(n)}{n} = \frac{1}{n^2} \sum_{k=1}^{n} k \eta_k = \frac{n+1}{2n} \frac{2}{n(n+1)} \sum_{k=1}^{n} k \eta_k$$

But $\frac{2}{n(n+1)}\sum_{k=1}^{n}k\eta_{k}$ is simply the arithmetic mean of the first

 $\frac{n(n+1)}{2}$ terms of the sequence η_1 , η_2 , η_2 , η_3 , η_3 , η_3 , . . . , whose

limit, we have seen, is zero. Thus $\frac{2}{n(n+1)}\sum_{k=1}^{n}k\eta_{k}\to 0$ as

 $n \to \infty$, from which follows $\lim_{n \to \infty} \frac{F(n)}{n} = 0$. Finally, we have

$$\lambda_n \equiv F(n) - F(n-1) = \eta_n - \frac{1}{n}F(n-1) \to 0 \text{ as } n \to \infty.$$

We now proceed to the problem of determining the form of F(n). It is clear from F(mn) = F(m) + F(n) that we only need to know the value of F(n) for prime n. Indeed, for arbitrary n let $n = p_1^{\alpha_1} \cdot \cdot \cdot \cdot p_s^{\alpha_s}$ be the prime factorization of n; then we see that $F(n) = \alpha_1 F(p_1) + \cdot \cdot \cdot + \alpha_s F(p_s)$. We now put, for all prime $p, F(p) = c_p \ln p$, where \ln represents, as usual, the natural logarithm. Then $F(n) = \alpha_1 c_{p_1} \ln p_1 + \cdot \cdot \cdot + \alpha_s c_{p_s} \ln p_s$.

Lemma 6. The sequence c_p , p = 2, 3, 5, 7, 11, ..., contains a largest member.

Proof. Assume the contrary; then it is possible to construct an infinite sequence of primes $p_1 < p_2 < p_3 < \cdots$ such that $p_1 = 2$, and p_{i+1} is the first prime greater than p_i for which $c_{p_{i+1}} > c_{p_i}$. It follows from this construction that if q is a prime

less than p_i , then $c_i < c_{p_i}$. For i > 1, let $p_i - 1 = q_1^{\alpha_1} \cdot \cdot \cdot q_i^{\alpha_i}$ be the prime factorization of $p_i - 1$. Now

$$\lambda_{p_i} = F(p_i) - F(p_i - 1)$$

$$= F(p_i) - \frac{F(p_i)}{\ln p_i} \ln (p_i - 1) + c_{p_i} \ln (p_i - 1) - F(p_i - 1)$$

$$= \frac{F(p_i)}{p_i} \frac{p_i}{\ln p_i} \ln \frac{p_i}{p_i - 1} + \sum_{j=1}^s \alpha_j (c_{p_i} - c_{q_j}) \ln q_j$$

Since $p_i - 1$ is necessarily even, one of the q_i must take on the value 2. Since furthermore $c_{p_i} > c_{q_i}$ for $j = 1, \ldots, s$, we have $\sum_{j=1}^{s} \alpha_j (c_{p_i} - c_{q_j}) \ln q_j \geq (c_{p_i} - c_2) \ln 2 \geq (c_{p_2} - c_2) \ln 2.$ Now as $i \to \infty$, $p_i \to \infty$; by Lemma 5, $\lambda_{p_i} \to 0$ and $\frac{F(p_i)}{p_i} \to 0$, while it is

easily shown that $\frac{p_i}{\ln p_i} \ln \frac{p_i}{p_i - 1} \to 0$. Therefore we must have $(c_{p_2} - c_2) \ln 2 \le 0$, or $c_{p_2} \le c_2$, which contradicts the definition of p_2 .

In precisely the same way we can show the existence of a smallest member of c_p , $p=2, 3, 5, \ldots$

Lemma 7. $F(n) = c \ln n$, where c is a constant.

Proof. It suffices to show that all the c_p are equal. Suppose there is a prime p' such that $c_{p'} > c_2$. Let p be that prime for which c_p is a maximum; then $c_p > c_2$. Let m be a positive integer and $q_1^{\alpha_1} \cdot \cdot \cdot q_s^{\alpha_s}$ be the prime factorization of $p^m - 1$. From F(mn) = F(m) + F(n) follows $\frac{F(p^m)}{\ln p^m} = c_p$; then just as in the proof of Lemma 6 we obtain $\lambda_{p^m} \ge \frac{F(p^m)}{p^m} \frac{p^m}{\ln p^m} \ln \frac{p^m}{p^m - 1} + (c_p - c_2) \ln 2$. Letting $m \to \infty$, we obtain $(c_p - c_2) \ln 2 \le 0$, which contradicts $c_p > c_2$. In precisely the same way we can show the nonexistence of any prime q for which $c_p < c_2$; thus

all the c_p are equal. We can now complete the proof of the theorem.

Let $p = \frac{r}{s}$ for integer r, s. By Lemma 4 we have

$$H\left(\frac{1}{s}, \dots, \frac{1}{s}\right) = H\left(\frac{r}{s}, \frac{s-r}{s}\right) + \frac{r}{s}H\left(\frac{1}{r}, \dots, \frac{1}{r}\right) + \frac{s-r}{s}H\left(\frac{1}{s-r}, \dots, \frac{1}{s-r}\right)$$

from which follows

$$H(p, 1 - p) = F(s) - pF(r) - (1 - p) F(s - r)$$

$$= c \ln s - pc \ln r - (1 - p)c \ln (s - r)$$

$$= c \left(p \ln \frac{s}{r} + (1 - p) \ln \frac{s}{s - r} \right)$$

$$= c \left(p \ln \frac{1}{p} + (1 - p) \ln \frac{1}{1 - p} \right)$$

By continuity, this result immediately extends to all irrational p. Using condition 3, it follows at once by induction on n that

$$H(p_1, \ldots, p_n) = c \sum_{i=1}^n p_i \ln \frac{1}{p_i}$$

We notice that $H(p_1, \ldots, p_n)$ is of the form $\sum_{i=1}^n p_i I(p_i)$, and that if we take c > 0, then I(p) is an increasing function of 1 - p.

Remarks

i. C. E. Shannon's original paper (Shannon [1]) is reprinted in Shannon and Weaver [1]. The application of information theoretic concepts to physics is thoroughly discussed by Brillouin [1]; a collection of various papers in which possible applications to chemistry, biology, and psychology are investigated has been edited by Quastler [1, 2].

The essentially statistical nature of most communication systems had been recognized before Shannon's work. In particular, it had been strongly emphasized by Wiener [1], who was the

first to use it in attacking problems of prediction and filtering. However, it seems fair to say that the concepts of channel and information content, as well as the formulation of the basic coding theorems, are due uniquely to Shannon.

ii. If, to the three conditions which we have imposed on $H(p_1, \ldots, p_n)$, one adds the requirement that

$$F(n) = H\left(\frac{1}{n}, \cdots, \frac{1}{n}\right)$$

be an increasing function of n, then the derivation of the form of $H(p_1, \ldots, p_n)$ is much simpler (cf. Appendix 2, Shannon and Weaver [1]). Khintchine [1] assumes $H(p_1, \ldots, p_n) \leq F(n)$ and the assertion of Lemmas 2 and 4, in addition to conditions 1 and 2, to obtain the form of $H(p_1, \ldots, p_n)$. The neat treatment given here is due to Fadiev [1]. The result used, to the

effect that $\lim_{i\to\infty} a_i = a$ implies $\lim_{n\to\infty} \frac{1}{n} \sum_{i=1}^n a_i = a$, is easily proved by standard "epsilon-delta" arguments (cf. Theorem B, p. 202 of Halmos [1]).

CHAPTER TWO

Basic Properties of H(X)

2.1. Viewpoint

We have shown in the previous chapter how a few simple and very reasonable requirements on the behavior of $H(p_1, \ldots, p_n)$ suffice to determine it uniquely up to a multiplicative constant. Interesting as this result is, it in no way implies that the quantity $H(p_1, \ldots, p_n)$ lends itself to any useful or interesting applications. What we wish to stress here is that the usefulness of $H(p_1, \ldots, p_n)$ lies not in the fact that it satisfies conditions 1 to 3 of Chap. 1, but rather in that it appears in a fundamental role in many basic problems of coding and communication. Our point of view henceforth will be to assume the form $H(p_1, \ldots, p_n) = \sum_{n=0}^{\infty} P(p_n, \ldots, p_n)$

 $-c\sum_{i=1}^{n} p_{i} \ln p_{i}$ without further comment, leaving its justification

to the various theorems whose proof is our main objective.

2.2. Fundamental Inequalities

Let X be an abstract set consisting of a finite number of elements x. Let p() be a probability distribution defined over X, i.e., p(Q) is a non-negative number defined for each subset Q of X,

with the properties that p(X) = 1 and $p(Q_1 \cup Q_2) = p(Q_1) + p(Q_2)$ if Q_1 and Q_2 are disjoint. The totality of objects (X,x) and $p(\cdot)$ is called a finite probability space.

In terms of our discussion in Chap. 1, we see that any finite probability space can be considered an information source. We will define the information content H(X) of such a source to be the non-negative quantity $-\sum_{X} p(x) \log p(x)$, where here and

henceforth the base of the logarithm is 2. The choice of base is clearly a choice of units only; with the conventional choice which we have made, the unit of information is called the "bit." To avoid difficulties later on, we will define the indeterminate form $0 \cdot \log 0$ to have the value zero.

Suppose that (X,x) and (Y,y) are two finite abstract spaces. We denote by $X \otimes Y$ the finite abstract space consisting of all pairs (x,y), and by $p(\cdot,\cdot)$ a probability distribution over $X \otimes Y$. The information content of this source we write as $H(X,Y) = -\sum_X \sum_Y p(x,y) \log p(x,y)$. The distribution $p(\cdot,\cdot)$ gives rise to a distribution $p(x) = \sum_X p(x,y)$ over (X,x), and similarly to a distribution $p(y) = \sum_X p(x,y)$ over (Y,y). The information contents of the two sources thus defined are again denoted by H(X) and H(Y) respectively, without regard to the origin of p(x) or p(y). Further, the conditional probability $p(x|y) = \frac{p(x,y)}{p(y)}$ is, for each y satisfying p(y) > 0, a probability distribution over X. We may therefore define a conditional information content $H(X|y) = -\sum_X p(x|y) \log p(x|y)$, and also an average condi-

$$H(X|Y) = \sum_{Y} p(y)H(X|y) = -\sum_{X} \sum_{X} p(x,y) \log p(x|y)$$

tional information content