Benford’s Law, A Growth Industry

Kenneth A. Ross

Abstract. Often data in the real world have the property that the first digit 1 appears about 30% of the time, the first digit 2 appears about 17% of the time, and so on with the first digit 9 appearing about 5% of the time. This phenomenon is known as Benford’s law. This paper provides a simple explanation, suitable for nonmathematicians, of why Benford’s law holds for data that have been growing (or shrinking) exponentially over time. Two theorems verify that Benford’s law holds if the initial values and rates of growth of the data appear at random.

INTRODUCTION. To get started, consider the sequence 2, 4, 8, 16, 32, 64, 128, \ldots, 2^{1000}. Each of these numbers has a first digit, starting out with 2, 4, 8, 1, 3, 6, 1. About how many of these numbers do you think have first digit 1? What about the other eight possible first digits? Consider the same questions for 9, 81, 729, \ldots, 9^{1000}, where the first few first digits are 9, 8, 7, 6, 5, 4, 4, 3, 3, 2, 2, 2, 2. We will return to these questions after Table 3.

Way back in 1881, Simon Newcomb observed that certain sets of data from the real world have the property that the first digits of the numbers do not appear uniformly often. In fact, the first digit 1 appears about 30% of the time, while the first digit 9 appears about 5% of the time. As often happens in science, this observation was forgotten and rediscovered later. Frank Benford observed this phenomenon in 1938 [2], and it has become known as Benford’s law or the “first-digit law.” He gave numerous examples of data from many sources, including newspaper items, areas of rivers, street addresses, cost data, and populations. We give another example in Table 1, namely data based on the populations of 117 cities in Indiana, from Wikipedia based on the 2000 census.1 For a good history up to 1975 with many references, see Raimi [18]. More references will be given at the end of the paper.

<table>
<thead>
<tr>
<th>first digit</th>
<th>1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
<th>6</th>
<th>7</th>
<th>8</th>
<th>9</th>
</tr>
</thead>
<tbody>
<tr>
<td>number of cities</td>
<td>34</td>
<td>17</td>
<td>16</td>
<td>8</td>
<td>16</td>
<td>12</td>
<td>5</td>
<td>3</td>
<td>6</td>
</tr>
<tr>
<td>percent of cities</td>
<td>29.1</td>
<td>14.5</td>
<td>13.7</td>
<td>6.8</td>
<td>13.7</td>
<td>10.3</td>
<td>4.3</td>
<td>2.6</td>
<td>5.1</td>
</tr>
</tbody>
</table>

Recently, Benford’s law has gone public. In fact, it was discussed recently on a public radio program, but listeners complained that it was not explained. That was because it isn’t possible to explain it in a sound bite. Since Benford’s law applies to many different sorts of data, it is natural to view it as a statistical phenomenon, so most serious attempts at understanding it use probability. The most successful seems to be Theodore Hill’s very nice, but sophisticated, analysis in [13]. I found his analysis challenging, and I am reasonably acquainted with probability. Thus, neither Hill’s result nor its proof is accessible to young students and others who have an interest in Benford’s law. In my view, what is needed is a collection of special cases that are

1 http://dx.doi.org/10.4169/amer.math.monthly.118.07.571
1 This is the first and only state that I examined for this purpose.
readily understandable. In fact, I think it might be easiest to start with nonstatistical, totally deterministic examples, which will be formalized in Theorem 1. Theorems 2 and 3 provide versions of Benford’s law.

This article is intended for four audiences. The discussion up to Figure 1 could be used to explain Benford’s law to people with essentially no mathematical background, while the discussion up to the proof of Theorem 1, just after Table 5, could be used to explain the law to any motivated person who understands logarithms. Upper-division undergraduate students, willing to assume the uniform distribution theorem, should be able to read the proof of Theorem 1. The remaining proofs involve a little measure theory and harmonic analysis.

MODELING BENFORD’S LAW. Here we will model populations of some old cities in some region, with the following assumptions. Though the populations of the cities probably started out as random, each of them grows (or shrinks) exponentially over time, so that a census is a snapshot of the population-growth curves of the various cities. We will see that this is why the observed populations years later will approximately satisfy Benford’s law.

We will focus at first on one city. The population of a city satisfies the equation \( P(n) = ar^n \). Here \( n \) represents units of time in integers (such as years or decades). We are only interested in the first digits of the populations. Since multiplying them by any power of 10 won’t change the mathematics, we henceforth assume that the initial “population” \( a \) satisfies \( 1 \leq a < 10 \) and the rate of growth \( r \) satisfies \( 1 < r < 10 \).

More generally, suppose you have data collected regularly over time, and the data describe a quantity which is increasing or decreasing exponentially. The data might be populations as in Table 1, values of portfolios, weights of vegetables at a county fair, etc. Suppose the length of time over which the data is collected is sufficiently large; this can be quantified and will depend on the rate of growth. Then Bedford’s law applies. Roughly, that’s because the time spent when the first digit is 1 is relatively flat and is considerably longer than the time spent when the curve is increasing more steeply. See Figure 1, showing how \( ar^t \) increases from \( 10^k \) to \( 10^{k+1} \). (If \( r < 1 \), the picture is the mirror image of the one in Figure 1.) The first digit of \( ar^t \) is 1 for the time interval \([c_1, c_2]\), it is 2 for the time interval \([c_2, c_3]\), etc. So clearly the first digit is 1 for much more time than it is 9.

We want to quantify these observations, so look again at Figure 1. For \( d = 1, 2, \ldots, 9 \), the time that the first digit is \( d \) is the length of the interval \([c_d, c_{d+1}]\). To determine these times, we observe that the graph of \( ar^t \) has height \( d \cdot 10^k \) at \( c_d \). Thus

\[
ar^{c_d} = d \cdot 10^k.
\]

Using the rules of logarithms, we obtain

\[
\log_{10} a + c_d \log_{10} r = \log_{10} d + k
\]

and similarly

\[
\log_{10} a + c_{d+1} \log_{10} r = \log_{10} (d+1) + k.
\]

Thus

\[
(c_{d+1} - c_d) \log_{10} r = \log_{10} (d+1) - \log_{10} d,
\]

\footnote{Several articles study applications of Benford’s law to the detection of accounting fraud; see, e.g., [7] and [15].}
so that
\[c_{d+1} - c_d = A[\log_{10}(d + 1) - \log_{10} d],\]
where \(A\) is the constant \(1/\log_{10} r\), which does not depend on \(k\). Moreover, the time that the curve takes to rise from \(10^k\) to \(10^{k+1}\) is exactly \(c_{10} - c_1 = A[\log_{10}(10) - \log_{10} 1] = A\). In other words, in this time interval, the fraction of time that the first digit of \(ar^k\) is \(d\) is
\[
\log_{10}(d + 1) - \log_{10} d = \log_{10}\left(\frac{d + 1}{d}\right).
\] (1)
Moreover, this observation applies to all such intervals \([10^k, 10^{k+1}], k = 0, \pm 1, \pm 2, \ldots\). The numbers in equation (1) are the numbers that arise in Benford’s law. Approximations of them are given in Table 2.

Table 2.

<table>
<thead>
<tr>
<th>(d)</th>
<th>1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
<th>6</th>
<th>7</th>
<th>8</th>
<th>9</th>
</tr>
</thead>
<tbody>
<tr>
<td>(\log_{10}\frac{d+1}{d})</td>
<td>0.30103</td>
<td>0.17609</td>
<td>0.12494</td>
<td>0.09691</td>
<td>0.07918</td>
<td>0.06695</td>
<td>0.05799</td>
<td>0.05115</td>
<td>0.04576</td>
</tr>
</tbody>
</table>

Let’s return to discrete time, i.e., return to \(P(n) = ar^n\) for integers \(n\), and visualize these values in Figure 1, i.e., visualize the points \((n, P(n))\) marching along the curve. From the statement involving (1), one would expect that the fraction of these points that would land in the interval \([c_d, c_{d+1}]\) would be about \(\log_{10}\left(\frac{d+1}{d}\right)\). This makes the following theorem very plausible, though its proof (given later) is somewhat more complicated.

For any \(x \geq 1\), we write \(D(x)\) for the first digit of \(x\).

**Theorem 1.** Suppose that \(1 \leq a < 10\), \(1 < r < 10\), and that \(r\) is rational. Then
\[
\lim_{n \to \infty} \frac{1}{n} \#\{k \leq n : D(ar^k) = d\} = \log_{10}(d + 1) - \log_{10} d,
\] (2)
for \(d = 1, 2, 3, 4, 5, 6, 7, 8, 9\).
The uninvited hypothesis, that \( r \) be rational, comes up in the proof. In fact, equation (2) holds for all but countably many real numbers \( r \) in \((1, 10)\). Theorem 1 appears in Raimi [18] and [19], and it essentially appears in Diaconis [6].

Theorem 1 suggests that, for reasonable choices of \( r \) and \( a \), the numbers \( \# \{k \leq 100: D(ar^k) = d\} \) will be close to \( 100[\log_{10}(d + 1) - \log_{10}d] \). It also suggests that the numbers \( \#\{11 \leq k \leq 110: D(ar^k) = d\} \) will be close to the same values, and I prefer to look at this because the first few values of \( D(ar^k) \) are likely to be predictable, especially if \( r \) is close to 1. Table 3 gives some examples of actual counts. For example, the third column indicates that the theorem predicts that about 30.1 out of 100 values of \( a \cdot r^k \) will have first digit equal to 1. The fourth row indicates that exactly 31 of the values \( \{9 \cdot 2^k: 11 \leq k \leq 110\} \) have first digit equal to 1.

Table 3.

<table>
<thead>
<tr>
<th>expected # first digits</th>
<th>30.1</th>
<th>17.6</th>
<th>12.5</th>
<th>9.7</th>
<th>7.9</th>
<th>6.7</th>
<th>5.8</th>
<th>5.1</th>
<th>4.6</th>
</tr>
</thead>
<tbody>
<tr>
<td>( a ) growth rate ( r )</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
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<td>4</td>
<td>5</td>
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<td>7</td>
<td>8</td>
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</tr>
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<td>13</td>
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<td>6</td>
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<td>4</td>
</tr>
<tr>
<td>1</td>
<td>3</td>
<td>32</td>
<td>16</td>
<td>14</td>
<td>10</td>
<td>8</td>
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<td>12</td>
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<td>7</td>
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<td>5</td>
</tr>
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<td>14</td>
<td>10</td>
<td>7</td>
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<td>4</td>
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<td>1</td>
<td>1.1</td>
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<td>6</td>
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<tr>
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<td>33</td>
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<tr>
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<tr>
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<tr>
<td>9</td>
<td>1.01</td>
<td>70</td>
<td>30</td>
<td>0</td>
<td>0</td>
<td>0</td>
<td>0</td>
<td>0</td>
<td>0</td>
</tr>
</tbody>
</table>

The answers to the questions at the beginning of the paper are similar. The nine answers to the first question about \( 2, 4, 8, 16, \ldots, 2^{1000} \) are 301, 176, 125, 97, 79, 69, 66, 52, 45. The answers to the question about \( 9, 81, 729, \ldots, 9^{1000} \) are 297, 177, 125, 98, 80, 67, 58, 52, 46.

Note that \( r = 1.01 \) is too close to 1 to illustrate the theorem only checking \( D(ar^k) = d \) for \( k \) up to 110. On the other hand, with \( a = 1 \) and \( r = 1.01 \), the counts for \( \#\{101 \leq k \leq 1100: D(ar^k) = d\} \), which should be close to \( 1000[\log_{10}(d + 1) - \log_{10}d] \), turn out to be 279, 173, 144, 112, 87, 61, 54, 48, and 42.

Theorem 1 is interesting, but it doesn’t explain the model involving old cities. It explains the model if we focused on one city and its various populations over a long period of time. With several cities in mind, it seems likely that their initial populations
and growth rates are random. Therefore, the following version of Benford’s law seems quite plausible.

**Theorem 2.** \( \lim_{n \to \infty} \Pr[D(ar^n) = d] = \log_{10}(d + 1) - \log_{10} d. \)

Probability? We assume that the initial population \( a \) is selected at random from \([1, 10)\) by any random process, but the reader won’t lose anything by assuming that the choices are uniformly distributed, i.e., for each interval \( I \) in \([1, 10)\), the values \( a \) fall into the interval \( I \) with probability length \( I/9 \). We assume that the growth rate is selected uniformly from a fixed interval \([s, u)\) in \([1, 10)\) so that, for each interval \( I \) in \([s, u)\), the values are selected with probability length \( I/(u - s) \).

Theorem 2 implies that if \( n \) is sufficiently large, and if we take a sufficiently large sample (say of size \( N \)) of values \((a, r)\) uniformly selected from the rectangle \([1, 10) \times [s, u)\), then

\[
\#\{(a, r) \text{ in the sample satisfying } D(ar^n) = d\} \approx N[\log_{10}(d + 1) - \log_{10} d],
\]

for \( d = 1, 2, 3, 4, 5, 6, 7, 8, 9 \). We illustrate counts, by taking random samples from \([1, 10) \times [1, 1.5)\) of size \( N = 1000 \). See Table 4.

<table>
<thead>
<tr>
<th>expected # first digits (\rightarrow)</th>
<th>301</th>
<th>176</th>
<th>125</th>
<th>97</th>
<th>79</th>
<th>67</th>
<th>58</th>
<th>51</th>
<th>46</th>
</tr>
</thead>
<tbody>
<tr>
<td>value of ( n ) (\downarrow)</td>
<td>1</td>
<td>2</td>
<td>3</td>
<td>4</td>
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<td>6</td>
<td>7</td>
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<td>10</td>
<td>273</td>
<td>190</td>
<td>140</td>
<td>129</td>
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<td>10</td>
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<td>92</td>
<td>70</td>
<td>66</td>
<td>61</td>
<td>46</td>
<td>43</td>
</tr>
<tr>
<td>20</td>
<td>305</td>
<td>176</td>
<td>127</td>
<td>89</td>
<td>74</td>
<td>73</td>
<td>58</td>
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<td>40</td>
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<tr>
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<td>154</td>
<td>143</td>
<td>85</td>
<td>75</td>
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<td>50</td>
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<td>174</td>
<td>127</td>
<td>103</td>
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<td>57</td>
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<tr>
<td>100</td>
<td>303</td>
<td>194</td>
<td>122</td>
<td>109</td>
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<td>50</td>
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<td>49</td>
</tr>
<tr>
<td>100</td>
<td>306</td>
<td>179</td>
<td>132</td>
<td>91</td>
<td>78</td>
<td>58</td>
<td>61</td>
<td>45</td>
<td>50</td>
</tr>
</tbody>
</table>

The table seems consistent with Theorem 2, though the counts based on random samples vary more than I had expected.

Theodore P. Hill has suggested that Theorem 2 would be a more realistic model if the rates of growth varied over time. Theorem 3 below was motivated by this observation.

**Theorem 3.** \( \lim_{n \to \infty} \Pr[D(ar_1 r_2 \ldots r_n) = d] = \log_{10}(d + 1) - \log_{10} d. \)

As with Theorem 2, \( a \) is selected at random from \([1, 10)\) by any random process, but now each \( r_i \) is selected uniformly from some fixed interval \([s, u)\) in \([1, 10)\) and the selections of \( r_i \) are assumed to be independent. Table 5 gives some results, each based on samples of size 1000 from \([1, 10) \times [1, 1.5)\) of size \( N = 1000 \). See Table 4.

The requirements of uniform selection in Theorems 2 and 3 can be substantially relaxed, as will be explained later.
Table 5.

<table>
<thead>
<tr>
<th>value of n ↓</th>
<th>1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
<th>6</th>
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<td>178</td>
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<tr>
<td>100</td>
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<td>73</td>
<td>69</td>
<td>54</td>
<td>49</td>
<td>57</td>
</tr>
</tbody>
</table>

**PROOF OF THEOREM 1.** We will use the function $g(x) := \log_{10} x \mod 1$, for $x$ in $[1, \infty)$. Here mod 1 means that we are taking the decimal or fractional portion of $\log_{10} x$. Thus $g$ is a function from $[1, \infty)$ onto $[0, 1)$, though we will sometimes restrict it to $[1, 10)$ where it is one-to-one. We actually prefer to identify the image set $[0, 1)$ with the circle of circumference 1 by using the wrapping function:

$$W(x) := \frac{1}{2\pi} (\cos 2\pi x, \sin 2\pi x) \text{ or, if you prefer, } W(x) := \frac{1}{2\pi} e^{2\pi i x}.$$ 

By laws of logarithms,

$$g(a r^k) = (\log_{10} a + k \log_{10} r) \mod 1 \text{ for all } k. \quad (4)$$

First we assume $a = 1$, so that $g(r^k) = k \log_{10} r \mod 1$ for all $k$. Observe that $D(x) = d$ if and only if $d \cdot 10^j \leq x < (d + 1) \cdot 10^j$ for some integer $j$, so that $D(x) = d$ if and only if $\log_{10} d \leq \log_{10} x \mod 1 < \log_{10} (d + 1)$. Hence

$$D(x) = d \text{ if and only if } \log_{10} d \leq g(x) < \log_{10} (d + 1), \quad (5)$$

and

$$D(r^k) = d \text{ if and only if } \log_{10} d \leq k \log_{10} r \mod 1 < \log_{10} (d + 1). \quad (6)$$

Figure 2 shows a picture of the circle of circumference 1 and a picture of the regions corresponding to where $D = d$ for $d = 1, 2, \ldots, 9$. Old-timers may be interested that Raimi [18] included a picture similar to our second picture in Figure 2 and that he discussed its relationship to a circular slide rule.

We need the following well-known fact.

**Uniform Distribution Theorem.** For any irrational $x$ in $[0, 1)$, the sequence $(k x \mod 1)$ is uniformly distributed. That is, for any interval $I$ in $[0, 1)$, we have

$$\lim_{n \to \infty} \frac{1}{n} \# \{ k \leq n : k x \mod 1 \text{ is in } I \} = \text{length}(I).$$

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This result is false if \( x \) is rational, because then the sequence \((kx \mod 1)\) takes on only finitely many values. Before we discuss and prove this theorem, we apply it to \( x = \log_{10} r \), so we need \( \log_{10} r \) to be irrational; but this is true by [17, Theorem 2.11] since we’re assuming \( r \) is rational and \( 1 < r < 10 \). We conclude that for any integer \( d \), \( 1 \leq d \leq 9 \),

\[
\lim_{n \to \infty} \frac{1}{n} \# \{ k \leq n : \log_{10} d \leq k \log_{10} r \mod 1 < \log_{10} (d + 1) \} = \log_{10} (d + 1) - \log_{10} d.
\]

Equation (2) for \( a = 1 \) follows from this and the observation in (6). The proof works even if \( a > 1 \), because the uniform distribution theorem remains true if \( kx \mod 1 \) is replaced by \((c + kx) \mod 1 \) for some constant \( c \), since the sequence \(( (c + kx) \mod 1 ) \) is just the sequence \(( kx \mod 1 ) \) rotated a constant angle around the circle. This completes the proof of Theorem 1.

As noted in Raimi [19], the conclusion

\[
\lim_{n \to \infty} \frac{1}{n} \# \{ k \leq n : D(FIB(k)) = d \} = \log_{10} (d + 1) - \log_{10} d,
\]

for the Fibonacci sequence \((FIB(k))\), has been noted by D. Singmaster and others since the 1960s. This is true because this sequence is asymptotic with the golden-mean geometric sequence \(((1 + \sqrt{5})/2)^k\).

The proof of the uniform distribution theorem is somewhat complicated, but the result is reasonable and is what one would hope would be true. The theorem was proved, independently, by P. Bohl, W. Sierpinski, and H. Weyl in 1909–1910. For two proofs, see [17, Theorem 6.3]. One proof uses continued fractions; the other uses Fourier series. An intricate proof, using only limits and epsilon-delta arguments, is given in [16, Chapter 2, Theorem 8.7]. In [16, Section 5.9], the uniform distribution theorem is discussed as an application of Birkhoff’s individual ergodic theorem. We will give another, and more modern, proof.

First, we collect together some standard facts in the setting of the compact group \( G = [0, 1) \); here 0 and 1 are identified. For a measure \( \mu \) on \( G \), its Fourier-Stieltjes transform \( \hat{\mu} \) is defined by \( \hat{\mu}(m) := \int_G \chi_m \, d\mu \) for integers \( m \), where \( \chi_m \) is defined by \( \chi_m(x) := \exp(2\pi imx) \) for \( x \in G \). A good reference is [14, Chapter 1, §7.4, §7.12].
I thank my friend Robert B. Burckel for pointing out how the use of this transform simplifies our proofs.

**Lemma 1.** Given a sequence \((\mu_n)\) of probability measures on \(G = [0, 1)\) and a probability measure \(\mu\) on \(G\), statements (a)–(d) below are equivalent.

(a) \(\mu_n \to \mu\) in the weak-* (or vague) topology, i.e., \(\lim_n \int_G f \, d\mu_n = \int_G f \, d\mu\) for all continuous functions on \(G\). \(^3\)
(b) If \(F \subseteq G\) is closed, then \(\limsup_n \mu_n(F) \leq \mu(F)\).
(c) If \(V \subseteq G\) is open, then \(\liminf_n \mu_n(V) \geq \mu(V)\).
(d) \(\lim_n \hat{\mu}_n(m) = \hat{\mu}(m)\) for all integers \(m\).

Also, if (a)–(d) are true, we have:

(e) If \(\mu(\{x\}) = 0\) for all \(x \in [0, 1)\), then
\[
\lim_n \mu_n(I) = \mu(I) \quad \text{for each interval } I \text{ in } G.
\] (7)

**Proof.** The equivalence of (a)–(c) is an elementary fact that holds in all metric spaces; see, for example, [25, Theorem 3.1.5].

We have (a) \(\implies\) (d) because each function \(\chi_m\) is a continuous function on \(G\). We have (d) \(\implies\) (a) because the space of linear combinations of the functions \(\chi_m\) is uniformly dense in the space of continuous functions on \(G\), by an application of the Stone-Weierstrass theorem.

For assertion (e), note that there exist \(a\) and \(b\) so that \((a, b) \subseteq I \subseteq [a, b]\). Then (7) follows from (b) and (c) via the inequalities:
\[
\limsup_n \mu_n(I) \leq \limsup_n \mu_n([a, b]) \leq \mu([a, b]) \leq \liminf_n \mu_n((a, b)) \leq \liminf_n \mu_n(I).
\]

**PROOF OF THE UNIFORM DISTRIBUTION THEOREM.** Let \(P_n\) be the probability \(\frac{1}{n} \sum_{k=1}^n \delta_{kx}\) on \(G = [0, 1)\), where \(\delta_{kx}\) is the point mass at \(kx \mod 1\). The conclusion of the theorem is equivalent to
\[
\lim_{n} P_n(I) = \lambda(I) \quad \text{for each interval } I \text{ in } G,
\]
where \(\lambda\) is Lebesgue measure on \([0, 1)\). Since \(\hat{\lambda}(0) = \hat{P}_n(0) = 1\) and \(\hat{\lambda}(m) = 0\) for \(m \neq 0\), Lemma 1 shows that it suffices to confirm
\[
\lim_{n} \hat{P}_n(m) = 0 \quad \text{for } m \neq 0.
\]

Since \(\hat{\delta}_{kx}(m) = \exp(-2\pi imkx)\), we need to prove that
\[
\lim_{n} \frac{1}{n} \sum_{k=1}^n \alpha^k = 0 \quad \text{for } \alpha := \exp(-2\pi ix).
\] (8)

But, since \(mx\) is not an integer, we have \(\alpha \neq 1\) and so this geometric sum is equal to
\[
\frac{\alpha^n - 1}{\alpha - 1},
\]
which is clearly bounded in \(n\). Hence (8) holds, completing the proof.

\(^3\)Since 0 and 1 are identified, continuity of \(f\) includes the requirement \(\lim_{x \to 1^-} f(x) = f(0)\).
**INDUCED MEASURES.** If \( \phi \) is a measurable function from a measurable space \((X, \mathcal{S})\) to a measurable space \((Y, \mathcal{T})\), then every measure \( \nu \) on \( X \) induces a measure \( \nu \circ \phi^{-1} \) on \( Y \) defined by

\[
\nu \circ \phi^{-1}(E) = \nu(\phi^{-1}(E)) \quad \text{for measurable sets } E \subseteq Y.
\]

(9)

Observe that

\[
\int f(y) d(\nu \circ \phi^{-1})(y) = \int f(\phi(x)) d\nu(x)
\]

for bounded measurable functions \( f \) on \( Y \).

**PROOF OF THEOREM 2.** In this theorem, the \( r \)'s are selected at random from a fixed interval \([s, u]\) in \([1, 10]\). Let \( U \) be the uniform probability on \([s, u]\), so that

\[
U(I) := \frac{\text{length}(I)}{u - s} \quad \text{for each interval } I \subseteq [s, u].
\]

Also, the \( a \)'s are selected, independently of the \( r \)'s, using any probability measure \( \mu \) on \([1, 10]\). We again use the function \( g : [1, 10) \to [0, 1) \) given by \( g(r) = \log_{10} r \). We have \( D(ar^n) = d \) if and only if \( d \cdot 10^k \leq ar^n < (d + 1)10^k \) for some \( k \) and hence if and only if

\[
\log_{10} d \leq g(a) + ng(r) < \log_{10}(d + 1);
\]

here addition is in the group \([0, 1)\). Thus we have

\[
\Pr[D(ar^n) = d] = \int_1^{10} \int_1^{10} \chi_{I_d}(g(a) + ng(r)) d\mu(a) dU(r),
\]

(11)

where \( \chi_{I_d} \) is the indicator function of the interval \( I_d = [\log_{10} d, \log_{10}(d + 1)) \). Our goal is to show that the quantity in (11) converges to \( \log_{10}(d + 1) - \log_{10} d = \lambda(I_d) \).

Obviously \( dU(r) = h(r) dr \) for a bounded measurable function \( h \) on \([1, 10] \), i.e.,

\[
\int_1^{10} f(r) dU(r) = \int_1^{10} f(r) h(r) dr
\]

(12)

for bounded measurable functions \( f \) on \([1, 10]\). We mention this here because we will use only (12) in the proof, so that it will be easy to see how to generalize this result. Let \( P \) and \( Q \) be the induced measures \( U \circ g^{-1} \) and \( \mu \circ g^{-1} \) on \([0, 1)\). Then from (11) and (10), with \( X = [1, 10] \) and \( Y = [0, 1) \), we obtain

\[
\Pr[D(ar^n) = d] = \int_0^{1} \int_0^{1} \chi_{I_d}(b + nx) dQ(b) dP(x).
\]

For each \( n \), let \( P_n \) be the induced measure \( P \circ \phi^{-1}_n \) on \([0, 1)\), where \( \phi_n(x) = nx \text{ mod } 1 \) for \( x \) in \([0, 1)\). Then, using (10) again, with \( X = Y = [0, 1) \), we obtain

\[
\Pr[D(ar^n) = d] = \int_0^{1} \int_0^{1} \chi_{I_d}(b + \phi_n(x)) dP(x) dQ(b)
\]

(13)

\[
= \int_0^{1} \int_0^{1} \chi_{I_d}(b + y) dQ(b) dP_n(y).
\]
This is exactly the integral with respect to the convolution $Q * P_n$; see [14, Chapter I, §7.10] or [10, (19.10)]. Hence

$$\Pr[D(\ar^n) = d] = \int_0^1 \chi_{I_d} d(Q * P_n) = (Q * P_n)(I_d),$$

and we want to show $\lim_n (Q * P_n)(I_d) = \lambda(I_d)$. As in the proof of the uniform distribution theorem, it suffices to show

$$\lim_n \widehat{Q * P_n}(m) = 0 \text{ for } m \neq 0.$$

Since $\widehat{Q * P_n}(m) = \widehat{Q}(m)\widehat{P_n}(m)$ for all $m$, it suffices to show

$$\lim_n \widehat{P_n}(m) = 0 \text{ for } m \neq 0.$$

Since $\chi_m(\phi_n(x)) = \exp(2\pi im\phi_n(x)) = \exp(2\pi imnx) = \chi_{mn}(x)$, we have

$$\widehat{P_n}(m) = \int_0^1 \chi_m(x) dP_n(x) = \int_0^1 \chi_m(\phi_n(x)) dP(x) = \int_0^1 \chi_{mn}(x) dP(x) = \widehat{P}(mn).$$

Using (12) and (10), it can be shown that $dP = H d\lambda$ for a bounded measurable function $H$ on $[0, 1)$. In other words,

$$\int_0^1 f(x) dP(x) = \int_0^1 f(x)H(x) dx$$

for bounded measurable functions $f$ on $[0, 1)$. In fact, $H(x) = \log_\epsilon(10) h(10^x)10^x$. Therefore, $\widehat{P} = \widehat{H}$ and, since $H$ is integrable, $\lim_{n \to \infty} \widehat{H}(m) = 0$ by the Riemann-Lebesgue lemma; see, e.g., [14, Chapter I, Theorem 2.8]. We conclude that for $m \neq 0$:

$$\lim_n \widehat{P_n}(m) = \lim_n \widehat{P}(mn) = \lim_n \widehat{H}(mn) = 0,$$

as desired.

**PROOF OF THEOREM 3.** We use the following notation from the proof of Theorem 2: $g$, $U$, $\mu$, $I_d$, $P$, and $Q$. This time, each $r_i$ is selected at random from $[1, 10)$ using $U$, and the $a$’s and $r_i$’s are selected independently. Now, all we need about $U$ is the probability $U$ is not supported on a finite subset of $[1, 10)$. \hspace{1cm} (14)

We have $D(\ar_1 r_2 \cdots r_n) = d$ if and only if

$$\log_{10} d \leq g(a) + g(r_1) + g(r_2) + \cdots + g(r_n) < \log_{10}(d + 1).$$

Thus we have

$$\Pr[D(\ar_1 r_2 \cdots r_n) = d] = \int_1^{10} \int_1^{10} \int_1^{10} \cdots \int_1^{10} \chi_{I_d}(g(a) + g(r_1) + g(r_2) + \cdots + g(r_n)) d\mu(a) dU(r_1) dU(r_2) \cdots dU(r_n), \hspace{1cm} (15)$$

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and we need to show that the quantity in (15) converges to \( \log_{10}(d + 1) - \log_{10} d \). Our tool will be \( Q \ast P^n \), where \( P^n \) is the convolution power \( P^n = P \ast P \ast \cdots \ast P \). In this setting, (10) tells us that

\[
\int_1^{10} f(g(r)) dv(r) = \int_0^1 f(x) d(v \circ g^{-1})(x)
\]

for measures \( v \) on \([1, 10] \) and for bounded measurable functions \( f \) on \([0, 1] \). Using this \( n + 1 \) times, the iterated integral in (15) becomes

\[
\int_0^1 \int_0^1 \cdots \int_0^1 \chi_{I_d}(x + x_1 + x_2 + \cdots + x_n) dQ(x) dP(x_1) dP(x_2) \cdots dP(x_n).
\]

This is the integral with respect to the convolution \( Q \ast P^n \), so that

\[
\Pr[D(a_1 r_2 \cdots r_n) = d] = \int_0^1 \chi_{I_d}(Q \ast P^n) = (Q \ast P^n)(I_d),
\]

and our goal is reduced to showing that

\[
\lim_{n \to \infty} (Q \ast P^n)(I_d) = \log_{10}(d + 1) - \log_{10} d = \lambda(I_d).
\]

Assumption (14) implies that \( P = U \circ g^{-1} \) is not supported on any finite subset of \([0, 1] \). As in the proof of Theorem 2, it suffices to prove

\[
\lim_n \widehat{P^n}(m) = 0 \text{ for all integers } m \neq 0.
\] (16)

This will be verified in the proof of the following lemma, which is an easy consequence of Lemma 1.

**Lemma 2.** For a probability measure \( P \) on \( G = [0, 1] \), we have

\[
P^n \to \lambda \text{ in the weak-* (or vague) topology}
\]

provided that the support of \( P \) is not a subset of any coset of a finite subgroup of \( G \).

**Proof.** As before, it suffices to verify (16). Since \( \widehat{P^n}(m) = \widehat{P}(m)^n \) for all positive integers \( n \), it suffices to show that \( |\widehat{P}(m)| < 1 \) for \( m \neq 0 \). If \( |\widehat{P}(m)| = 1 \), then \( \int_G \chi_m dP = e^{i\theta} \) for some real number \( \theta \). This implies that the support of \( P \) is a subset of \( A = \{ x \in G : \chi_m(x) = e^{i\theta} \} \). Given a fixed \( x_0 \) in \( A \), it follows that \( A - x_0 = \chi_m^{-1}(1) \), which is a closed proper subgroup of \( G \) since \( \chi_m \neq 1 \). This violates the hypothesis regarding the support of \( P \), since all closed proper subgroups of \( G \) are finite. \( \square \)

**NOTES FOR THE SPECIALIST.** Since (14) is the only property of \( U \) used in the proof of Theorem 3, Theorem 3 holds so long as the \( r_i \)'s are selected independently at random using a probability on \([1, 10] \) that is not supported on any finite subset of \([1, 10] \). In particular, Theorem 3 holds if the \( r_i \)'s are selected from \([1, 10] \) using Benford’s probability measure \( \beta \) where \( \beta((a, b)) = \log_{10}b - \log_{10}a \) for all \( (a, b) \subseteq [1, 10] \). However, in this case, \( P := \beta \circ g^{-1} \) and all its convolution powers \( P^n \) are equal to \( \lambda \), so Lemma 2 is not needed.

The equivalence (a) \( \iff \) (d) in Lemma 1 holds for any compact abelian group \( G \), where assertion (d) becomes: \( \lim_n \widehat{\mu_n}(\chi) = \widehat{\mu}(\chi) \) for all characters of \( G \).
Using essentially the same proof, Lemma 2 holds for any compact abelian group $G$, provided the support of $P$ is not a subset of any coset of a proper closed subgroup of $G$. For general compact groups, Lemma 2 holds provided the support of $P$ is not a subset of any coset of a proper closed normal subgroup of $G$, but the proof is more complicated than that of Lemma 2. The result goes back to Yukiyosi Kawada and Kiyosi Itô (1940) and Karl Stromberg (1960) [24]. This generality is interesting, even for finite groups, because this theorem implies that repeated selections based on a fixed random process, such as shuffling cards, lead to a well-mixed selection, i.e., have the effect of a uniform selection. In the case of shuffling cards, the underlying group is a large finite nonabelian group of permutations. For much more about card shuffling, see [1].

Under very general assumptions, Peter Schatte [23, §5] shows that

$$\lim_{N \to \infty} \frac{1}{N} \sum_{n=1}^{N} \Pr[D(a_1 r_2 \cdots r_n) = d] = \log_{10}(d + 1) - \log_{10} d.$$ 

In particular, for each $i$, the $r_i$’s can be selected using different probabilities $U_i$. His conclusion is weaker than ours, but his hypotheses are also weaker.

MORE REFERENCES. At this time, the website Benford Online Bibliography [3] has a little over 600 entries regarding Benford’s law, mostly articles but some website and video items as well. Here are further comments on some papers published after Raimi’s fine summary [19]. In [11], Theodore P. Hill gives a formulation of the significant-digit problem based on the natural assumption of base invariance, and he proves a theorem about the probabilities of digits in different positions. He also notes, in [12], that these probabilities are dependent. For example, the (unconditional) probability that the second digit is 2 is $\approx 0.109$, whereas the conditional probability that the second digit is 2, given that the first digit is 1, is $\approx 0.115$. Other interesting mathematical treatments of Benford’s law appear in Raimi [20], Cohen and Katz [5], and the new book by Rodney Nillsen [16]. For a survey of applications of Benford’s law to the natural sciences, including observations from the fields of physics, astronomy, geophysics, chemistry, and engineering, see [22].

SUGGESTED STUDY. For what sequences does Theorem 1 hold? That is, for what sequences $(a_k)$ do the limits

$$\lim_{n \to \infty} \frac{1}{n} \#\{k \leq n : D(a_k) = d\}$$

exist and, if they do, equal $\log_{10}(d + 1) - \log_{10} d$ for $d = 1, 2, \ldots, 9$? It is shown in [21] that the limit in (17) does not exist for the sequences $(k^N)_{k=1}^\infty$ for integers $N \geq 2$. However, an interesting subsequence converges to $\left(\sqrt[d]{d + 1} - \sqrt[d]{d}\right)/(\sqrt[10]{10} - 1)$. Perhaps the problem with $(k^N)_{k=1}^\infty$ is that it diverges too slowly to infinity. Consider any $r > 1$, that is not a power of 10, and any integer $N \geq 2$. Experimental evidence suggests that the limits in (17) exist for the sequence $(a_k)$ defined by $a_0 = r$ and $a_{k+1} = a_k^N$ for $k \geq 1$, and equal $\log_{10}(d + 1) - \log_{10} d$, for $d = 1, 2, \ldots, 9$. The case $r = N = 2$ is interesting. But, we have no proofs.

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