## DISSERTATION

## SOME TWO-STEP SAMPLING PROCEDURES

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## IT IS RECOMMENDED THAT THE DISSERTATION PREPARED BY

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## I. THE SAMPLE SIZE PROBLEM

One of the most important statistical problems is estimating the value of an unknown parameter in a given frequency function. If a point estimate is desired and the sample size is not fixed in advance then the experimenter must decide how large a sample should be taken. For most problems the cost of an experiment increases with the sample size. This increase in cost may be in financial terms or perhaps in terms of time or effort. On the other hand, a decrease in sample size may increase the variance of the estimate (loss of precision) or decrease the "closeness" of the estimate to the true value of the parameter, $i_{\text {. }}$, , the probability that the estimate is within a given distance of the true value decreases. Thus the problem is to devise some procedure for determining the smallest sample size which still allows the experimenter to obtain an estimate of the parameter with certain restrictions on precision, closeness, on some other criterion.

An experimenter may prefer to obtain an interval estimate of the unknown parameter. The desirability of small sample size is the same as in the point estimation case. For interval estimation a reduction in sample size will in general increase the width of the interval or decrease the confidence coefficient. The problem is to obtain the smallest sample size possible with certain restrictions, determined by the experimenter, on the width or confidence coefficient of the interval estimate.

When the experimenter has limited financial resources or time, it is possible that he might be forced to relax desired restrictions to reduce sample size. In this case the experimenter must evaluate relative losses between increased sample size and decreased usefulness of the results. This indicates that the quality of results in experiments is often related to financial resources and time restrictions.

Dantzig [12] in 1940 was the first to show that not all sample size problems can be solved by a one-step procedure. In particular, Dantzigis results showed that a one-step procedure cannot be devised to test "student's" hypothesis such that its power function is independent of the variance.

For the class of all distributions for which the mean exists, Bahadur and Savage [3] have shown that a purely sequential sampling scheme is not sufficient to provide a universal procedure for interval estimation of a mean with specified width and confidence coefficient. Also it has been shown by Farrell [14] that for estimation of the median, within the class of distributions possessing a unique median, a purely sequential scheme is both necessary and sufficient.

Sequential tests of statistical hypotheses as discussed by Wald [32], are described as "any statisticel test procedure which gives a specific rule, at any stage of the experiment (at the $n t h$ trial for each integral value of $n$ ), for making one of the following three decisions: (1) accept the hypothesis being tested (null hypothesis), (2) to reject the null hypothesis, (3) to continue the experiment by making an additional observation." Thus, after each trial in a succession of trials one of three decisions is made. Either decision
(1) or (2) is made, terminating the procedure, or (3) is made and another trial is performed. Some of the work in this field is contained in the set of references [2], [7], [11], [22], [23], [24], [26], [31], and [32].

Even though a sequential procedure may result in the smallest sample size possible for a given problem, the sequential method does present certain difficulties. It is not always possible or practical to sample from a population at an indefinite number of different times as must be done with this method. The same population may not be available to take more than a limited number of samples or it may be too costly to take samples at different times.

For certain problems it is possible to find a one-step procedure. The advantage of taking only one sample is obvious. Greenwood and Sandomire [20] have solved one such problem which is amenable to this type of sampling procedure. Their procedure gives the sample size required such that a confidence interval can be placed on the standard deviation of a normal population within 100 p percent of the true value. Given $p$ and $\alpha$, their method determines the sample size $n$ such that

$$
P\left[\left|\hat{\sigma}_{n}-\sigma\right|<p \sigma\right] \geq 1-\alpha
$$

where

$$
\hat{\sigma}_{n}=\left[\Sigma\left(x_{i}-\bar{x}\right)^{2} /(n-1)\right]^{1 / 2}
$$

and the $X_{i}$ are $n$ independent observations from a normal density with mean $\mu$ and variance $\sigma^{2}$. Graybill and Connell [17] give a
one-step procedure to obtain sample sizes such that the ratio of variances from two independent normal populations can be estimated within l00p percent of the true ratio with specified confidence coefficient. Similarly Epstein [13], using a one-step procedure, has estimated the mean in the exponential distribution within $100 \delta$ percent of its true value with specified confidence coefficient. In many common densities it is possible to construct a one-step procedure to estimate a parameter within a given percent of its true value.

A two-step sampling procedure for estimation of an unknown parameter can be defined as a procedure for computing an estimate, under certain desired restrictions, based on a sample of size $n$, where $n$ is determined by a first or preliminary sample. Some writers use the terms two-stage sampling, two-sampling, or double sampling to mean the same as two-step sampling.

Blum and Rosenblatt [9] give sufficient conditions for the existence of a two-step procedure for constructing confidence intervals of prescribed widths and confidence coefficients. Let $G$ be a family of distribution functions and let $\theta(0)$ be a real-valued functional defined on $G$. It is desired to make an interval estimate of $\theta(F)$ based on a sample from $F \in G$. For each positive integer $k$ let $F_{k}$ denote the product distribution function on Euclidean $k$-space induced by $F$. The corresponding probability measure is denoted by $P_{F_{n}}$. Let $G(m, \gamma, \delta)=[F \varepsilon G: m(F, \gamma, \delta) \geq m]$ where $m(F, \gamma, \delta)$ is the smallest positive integer such that for all $\Omega \geq m(F, Y, \delta)$ we have

$$
P_{F_{n}}\left[\left|\theta_{n}\left(X_{i}, \ldots, X_{n}\right)-\theta(F)\right| \leq \delta\right] \geq 1-\alpha_{0}
$$

Using this notation, Blum and Rosenblatt give the following theorem: Suppose there exists a decreasing sequence $\left[\mathrm{N}_{\mathrm{g}}\right]$ of Borel subsets of $R^{k}$ such that
(i) $N_{0}=R^{k}$ and $\lim _{j \rightarrow \infty} P_{F_{k}}\left[H_{j}\right]=0$ for every $E \in G$, and
(ii) There exists $\gamma \in(0, \alpha)$ and for each integer $j$ a positive integer $n_{j}$ such that

$$
\inf _{F \in G\left(n_{j}, \gamma, \delta\right)} P_{F_{k}}\left[N_{j}\right]>(1-\alpha) /(1-\gamma) .
$$

Then there exists a two-stage procedure with $k$ observations in the first stage for constructing a confidence interval for $\theta(F)$, of length $2 \delta$ and confidence $1-\alpha$.

A generalization of this theorem is also given for sufficient conditions which require an n-step procedure. Sufficient conditions for the existence of a two-step procedure for constructing confidence intervals of preassigned widths and confidence coefficients using only one observation in the first step are given by Abbott and Rosenblatt [1].

The first two-step procedure was given in 1945 by Stein [30] for estimating the mean $\mu$ from a normal population. In this procedure $d, m$ and $\alpha$ are specified and the sample size $n$ is determined from a preliminary sample of size $m$ such that

$$
\begin{equation*}
P[|\mu-\hat{\mu}| \leq d] \geq 1-\alpha \tag{1.1}
\end{equation*}
$$

where $\hat{\mu}$ is the mean of the combined sample. He also generalizes his method to confidence regions for means of several normal populations with equal but unknown variance.

Four years after the publication of Stein's article Ruben [28], working in ignorance of Stein's [30] results and using a different type of argument, rediscovered them. Ruben's results were achieved more simply and directly and have the advantage that the method generalizes quite naturally to deal with the more difficult problem of sampling from normal populations with unequal and unknown variances.

Procedures to determine the preliminary sample size in Stein's procedure have been given by Seelbinder [29] and Moshman [25]. These methods require some estimate of the range for the variance $\sigma^{2}$ of the population and are concemed with minimum expected sample size.

In Stein's two-step procedure to estimate the mean of a normal population the experimenter specifies $d$ and $1-\alpha$ in (1.1) in advance and the total number of observations is a random variable. In this case the cost of the experiment is not predetermined and may extend beyond the experimenter's resources. To exercise some control over cost Wormleighton [36] generalizes Stein's procedure so that a first sample can be taken to give an estimate of the variance after which the experimenter can decide on the total number of observations and the number of units the estimate of the mean may differ from the true value with a given confidence coefficient. Using Wormleighton's procedure the experimenter is still able to use all his data in making the estimate. Stein's results are extended by Chapman [10] to test hypotheses concerning the ratio of means of two normal populations
with power independent of the unknown variances. To do this, Stein's procedure is used for each population, under certain restrictions dependent upon the hypothesized ratio, and a test involving the difference of two student's t-variables is given. Also Chapman uses Stein's technique to test the hypothesis

$$
H: b=b_{0}
$$

in the regression problem where the $Y_{i}$ are independent random variables with $\sigma_{Y_{i}}$ unknown and

$$
E\left(Y_{i}\right)=a+b x_{i} .
$$

In this case the power is shown to be a function of $\left(b^{\prime}-b_{0}\right)\left(x_{1}-x_{2}\right) / \sqrt{z_{2}}$, where $b^{\prime}$ is the true value of $b, x_{1}$ and $x_{2}$ are the $x$ values at the ends of the range which are used in the first step of Stein's procedure, and $z$ can be chosen to obtain any prescribed power given $b^{\prime}$. This power function is independent of $\sigma_{Y}$ as desired.

Healy [21] also extends Stein's results and gives two-step procedures to construct simultaneous confidence intervals of prescribed widths and confidence coefficients for the following: (1) all normalized linear functions of means, (2) all differences between means, and (3) means of $k$ independent normal populations with common unknown variances. A partial generalization of (2) has been given by Ghurye and Robbins [15]. The method estimates the difference between means from normal populations with different variances. The second sample size, restricted by a cost constraint, is determined on the basis of the size of the preliminary sample and estimate of variance. The
variance of the estimate is given and is shown to be asymptotic to the minimum variance which would be obtained if the variances were known.

Bechhofer, Dunnett, and Sobel [5] give a two-step procedure to rank several normal populations according to their means when these populations have equal but unknown variance. This method is similar to Stein's. They also extend their solution to the more general case where the variances are unequal but the ratios are known. For the case of known variances Bechhofer [4] gives a single-sample multiple decision procedure.

Weiss [33] gives a two-step procedure for obtaining a confidence interval of preassigned width and confidence coefficient for quantiles of a continuous distribution. The only assumption is that the density function is unimodal.

Graybill [16] gives sufficient conditions for two-step estimation in certain parametric cases. The sample size is determined such that the probability is $\beta^{2}$ that the width of a confidence interval, with prescribed confidence coefficient $1-\alpha$, will be less than some preassigned value d. Graybill's theorem is as follows:

Let the chance variable $X$ be the width of a confidence
interval on a parameter $\mu$ based on a sample of size $n$.
Suppose that $X$ depends on $n$ and on an unknown parameter $\theta(\theta$ may be the parameter $\mu)$. Suppose also that there exists a function of $X, \theta$, and $n$, say $g(X ; \theta, n)$, such that if $Y=g(X ; \theta, n)$, then the distribution of $Y$ does not depend on any unknown parameters except $n$. Let $f(n)$ be a function of $n$ such that

$$
P[Y<f(n)]=\beta \quad \text { for any } \quad 0<\beta<1
$$

Let the solution of the equation $g(x ; \theta, n),=f(n)$ for $x$ be $x=h(\theta, n)$, and suppose the following are true for $x>0:$
(a) $g(x ; \theta, n)$ is monotonic increasing in $x$ for every $n$ and $\theta$.
(b) $h(\theta, n)$ is monotonic increasing for every $n$.
(c) $h(\theta, n)$ is monotonic decreasing in $n$ for every $\theta_{0}$
(d) $z$ is random variable which is available from step one of the procedure such that $P[t(z)>\theta]=\beta$ for $0<\beta<1$, where $t(z)$ is a function of $z$ which does not depend on any unknown parameters or on $n$.

Let $d$ and $\beta$ be specified in advance. Then if $n$ is such that the equation

$$
h[t(z), n] \leq d
$$

is satisfied $[t(z)$ is known] then the following inequality is true:

$$
P(X \leq d) \geq B^{2}
$$

This is a two-step procedure which is applicable for many common distributions. In particular, this procedure is applied to the variance of a normal distribution by Graybill and Morrison [19].

Birnbaum and Healy [8] attacked the sample size problem by giving rules for sampling in two steps so as to obtain an unbiased estimator
of a given parameter, having variance equal to, or not exceeding, a prescribed bound. They assume conditions satisfied by many distributions. It is applied to the means of the binomial, Poisson, and hypergeometric distributions, scale parameters in general and of the gamma distribution in particular, the variance of a normal, and a component of variance. This procedure can be applied to problems of interval estimation in two-steps by the use of Tchebycheff's inequality. Graybill and Connell [18] give a two-step procedure to estimate the parameter in the uniform density,

$$
f(u)=1 / \theta \quad ; \quad 0 \leq u \leq \theta,
$$

within $d$ units of the true value with specified confidence coefficient. This procedure is shown to give smaller sample sizes than that possible with Birnbaum and Healy's method using Tchebycheff's inequality.

Another type of sample size problem has been solved graphically by Birmbaum and Zuckerman [6]. To determine the smallest sample size for which the minimum and the maximum of a sample are the $100 \mathrm{~B} \%$ distribution-free tolerance limits at the probability level $\alpha$, one has to solve the equation

$$
N \beta^{N-1}-(N-1) \beta^{N}=1-\alpha
$$

given by Wilks [34]. The graph presented makes it possible to solve this equation with sufficient accuracy for almost all useful values of $\beta$ and $\alpha$.

The preceding is a resume of the type of work that has been done in the sample size problems. In this dissertation two estimation
problems will be solved with two-step procedures. In Chapter III a two-step procedure will be given to estimate the variance of a normal distribution within $d$ units with a specified confidence coefficient using an inequality derived in Chapter II. With minor modifications, the results of Chapter III can be extended to estimate the mean of the gamma distribution. A two-step procedure derived by a different type of argument will be given in Chapter IV to estimate the mean of a Poisson distribution within $d$ units with a specified confidence coefficient. The results in both Chapters III and IV are compared with Bimbaum and Healy's [8] method ising Tchebycheff's inequality. It is anticipated that the techniques used in this dissertation can be applied to other similar types of groblems. Chapter $V$ presents some of the sample size problems whicn have not yet been solved and discusses some of the problems which are associated with the solutions in Chapters III and IV。

## II. A TCHEBYCHEFF TYPE INEQUALITY FOR GAMMA

### 2.1 Introduction

A Tchebycheff type inequality is useful in many situations in statistics, but for certain densities it may be improved upon. For example, in Chapter III it's desirable to sharpen the inequality somewhat for the gamma density. The purpose of this chapter is to find an inequality that is an improvement of Tchebycheff's inequality for a random variable that is distributed as gamma with parameters $r$ and $\lambda$.

$$
\text { Let } \begin{aligned}
f(x) & =\frac{\lambda(x \lambda)^{r-1}}{\Gamma(x)} e^{-\lambda x}, x>0 . \\
& =0
\end{aligned}
$$

The problem is to prove that

$$
P(|x-r / \lambda|<a r / \lambda)>1-e^{-a \sqrt{2 r-1} / \sqrt{\pi}}
$$

for all $a>0, r \geq 1 / 2$, and $\lambda>0$. Let $v=\lambda x / r, r=n / 2$. Then the problem is equivalent to showing that

$$
\begin{equation*}
\int_{1-a}^{1+a} f_{1}(v) d v>1-e^{-a \sqrt{n-1} / \sqrt{\pi}} \tag{2.1}
\end{equation*}
$$

for all $a>0$ and $n \geq 1$, where $f_{1}(0)$ is the density of a chi-square divided by $n$, its degrees of freedom. Also, the inequality in (2.1) will be compared with

$$
\int_{1-a}^{1+a} f_{1}(v) d v \geq 1-2 / a^{2} n
$$

which is Tchebycheff's inequality for this problem. We shall assume n given. In the proof we shall use $y \sim z$ to mean $y z>0$.

### 2.2 Solution

By definition

$$
\begin{array}{rlr}
f_{1}(v) & =\frac{(n / 2)^{n / 2}}{\Gamma(n / 2)} v^{(n / 2)-1} e^{-(n / 2) v}, 0<v<\infty \\
& =0 & ,-\infty<v \leq 0 .
\end{array}
$$

Let

$$
f_{2}(v)=\frac{\sqrt{n-1}}{2 \sqrt{\pi}} \exp [-|v-1| \sqrt{n-1} / \sqrt{\pi}],-\infty<v<\infty .
$$

Define $h(a)$ by

$$
h(a)=\int_{1-a}^{1+a}\left[f_{1}(v)-f_{2}(v)\right] d v, a \geq 0, n \geq 1
$$

Thus (2.1) is true if

$$
\begin{equation*}
h(a)>0 \text { for all } 0<a<\infty \text { and } n \geq 1 \text {. } \tag{2.2}
\end{equation*}
$$

From Wilton [35] we obtain

$$
\Gamma(m+1)<\sqrt{2 \pi}(m+1 / 2)^{m+1 / 2} e^{-(m+1 / 2)}, m \geq-1 / 2 .
$$

If we let $m=n / 2-1$ we find that

$$
\begin{equation*}
\frac{(n / 2)^{n / 2}}{\Gamma(n / 2)} e^{-n / 2}>\frac{\sqrt{n-1}}{2 \sqrt{\pi}} \quad, \quad n \geq 1 . \tag{2.3}
\end{equation*}
$$

Let

$$
g_{1}(v)=\left[f_{1}(1-v)+f_{1}(1+v)\right] / 2 f_{2}(1+v), v \geq 0 .
$$

From (2.3) we obtain

$$
g_{1}(0)>1 \quad, \quad n \geq 1 \text {. }
$$

Equation (2.2) is true, which implies (2.1) is true, if there exists a $\mathrm{V}_{1}$ such that

$$
\frac{d}{d v} g_{1}(v)\left\{\begin{array}{l}
>0,0 \leq v<v_{1}  \tag{2.4}\\
\leq 0, v_{1} \leq v<\infty
\end{array}\right.
$$

since

$$
\frac{d}{d a} h(a) \sim g_{1}(a)-1 \quad, \quad 0<a<\infty
$$

and

$$
h(\infty)=0 .
$$

We shall show that a $v_{1}$ exists such that (2.4) is true. By definition

$$
g_{1}(v)= \begin{cases}r\left[(1-v)^{p-1} e^{(q+p) v}+(1+v)^{p-1} e^{(q-p) v}\right], & 0 \leq v<1 \\ r(1+v)^{p-1} e^{(q-p) v} & , 1 \leq v<\infty\end{cases}
$$

where $p=n / 2, q=\sqrt{n-1} / \sqrt{\pi}, r=p^{p} e^{-p / q \Gamma(p)}$.
Thus

$$
\begin{aligned}
& \frac{d}{d v} g_{1}(v) \sim \frac{e^{-q v}}{r} \frac{d}{d v} g_{1}(v) \\
& \quad= \begin{cases}{[(q+1)-(p+q) v](1-v)^{p-2} e^{p v}} \\
+[(q-1)-(p-q) v](1+v)^{p-2} e^{-p v}, & 0 \leq v<1 \\
{[(q-1)-(p-q) v](1+v)^{p-2} e^{-p v}} & 1 \leq v<\infty\end{cases}
\end{aligned}
$$

By definition $p>q^{2}$ which implies

$$
(q+1) /(p+q)>(q-1) /(p-q), n \geq 1 .
$$

Therefore

$$
\begin{equation*}
\frac{d}{d v} g_{1}(v) \leq 0, v \geq(q+1) /(p+c) \tag{2,5}
\end{equation*}
$$

To show the existence of such a $v_{1}$ in (2.4) we have proved (2.5) that $v_{1}$ must be less than or equal to $(q+1) /(p+q)$. To find $v_{1}$ we shall discuss three cases.

Case (i): $\quad 4 \leq n<\infty$

Let

$$
\begin{aligned}
& g_{2}(v)=[(1-v) /(1+v)]^{p-2} \\
& g_{3}(v)=e^{-2 p v}[(q-1)-(p-q) v] /[(q+1)-(p+q) v]
\end{aligned}
$$

Hence

$$
\begin{equation*}
\frac{d}{d v} g_{1}(v) \sim g_{2}(v)+g_{3}(v), 0 \leq v<(q+1) /(p+q) \tag{2,6}
\end{equation*}
$$

Differentiating we obtain

$$
\frac{d}{d v} g_{2}(v)=(p-2)[(1-v) /(1+v)] p^{-3}\left[-2 /(1+v)^{2}\right] \leq 0 .
$$

and

$$
\begin{aligned}
\frac{d}{d v} g_{3}(v)= & -2 p g_{3}(v)-e^{-2 p v}(p-q) /[(q+1)-(p+q) v] \\
& +e^{-2 p v}(p+q)[(q-1)-(p-q) v] /[(q+1)-(p+q) v]^{2} \\
\sim & -p[(q-1)-(p-q) v][(q+1)-(p+q) v]-p+q^{2} \\
= & -p\left(p^{2}-q^{2}\right) v^{2}+2 p q(p-1) v-q^{2}(p-1) \\
\leq & -p\left(p^{2}-q^{2}\right)\left[q(p-1) /\left(p^{2}-q^{2}\right)\right]^{2} \\
& +2 p q(p-1)\left[q(p-1) /\left(p^{2}-q^{2}\right)\right]-q^{2}(p-1) \\
\sim & -p+q^{2}<0 .
\end{aligned}
$$

Using (2.5), (2.6), and the knowledge that

$$
g_{2}(0)+g_{3}(0)>0
$$

and

$$
\frac{d}{d v}\left[g_{2}(v)+g_{3}(v)\right]<0,0 \leq v<(q+1) /(p+q),
$$

we see that there exists a $v_{1}<(q+1) /(p+q)$ such that (2.4) is true, implying (2.1) is true for this case.

Case (ii): $1 \leq n \leq 2$

In this case it can be shown that $(2.4)$ is true for $v_{2}=1$, implying (2.1) is true.

Case (iii) : $2<n<4$

By similar, though tedious, manipulations it can be shown that (2.4)
is true for some $v_{1}$ in the interval $(q / p,(q+1) /(p+q))$.
Thus (2.1) is true for all $n \geq 1$.
2.3 Comparison With Tchebycheff's Inequality

Let the density of $v$ be $f_{1}(0)$, a chi-square divided by $n$, its degrees of freedom. By Tchebycheff's inequality

$$
\begin{equation*}
P[|v-1|<a] \geq 1-2 / a^{2} n . \tag{2.7}
\end{equation*}
$$

In this chapter, by (2.1), the corresponding inequality is

$$
\begin{equation*}
P[|v-1|<a] \geq 1-e^{-a \sqrt{n-1} / \sqrt{\pi}} . \tag{2.8}
\end{equation*}
$$

Consider the ratio of $e^{-a} \sqrt{n-1} / \sqrt{\pi}$ and $2 / a^{2} n$, i.e.

$$
k(n, a)=.5 a^{2} n e^{-a \sqrt{n-1} / \sqrt{\pi}}, \quad 0<a<\infty, 1 \leq n<\infty .
$$

If $k(n, a)<1$, then ( 2,8 ) provides a better (larger) lower bound than (2.7) We shall show the values of $a$ and $n$ where $k(n, a)<1$ and hence where the method described in this chapter is better than Tchebycheff's inequality for the gamma density. Figure 2.1 shows that $k(n, a)>1$ only in a limited region. For instance $k(n, a)<1$ for $n \geq 7$ and $0<a<\infty$. Also, $k(n, a)<1$ for $0<a<1.36$ and $1 \leq n<\infty$ 。


Figure 2.1
III. SAMPLE SIZE REQUIRED FOR ESTIMATING THE VARIANCE WITHIN d UNITS OF THE TRUE VALUE

### 3.1 Introduction

The problem of estimating the variance $\left(\sigma^{2}\right)$ of a normal population arises in many experimental situations. J. A. Greenwood and M. M. Sandomire [20] have presented a means of obtaining the sample size required to estimate the variance of a normal population within a given per cent of its true value. An investigator may prefer to estimate the variance within a given number of units. This chapter will provide the sample size required to solve that problem.

Assume a preliminary sample of size $m ; z_{1}, z_{2}, \ldots, z_{m}$, is taken from a normal density with variance $c^{2}$. The unbiased estimator of the variance, $s_{m}^{2}$, is computed by the formula $s_{m}^{2}=(m-1)^{-1} \sum\left(z_{i}-\bar{z}\right)^{2}$. and $d$ and $l-\alpha$ are specified in advance. It is desired to determine $n$, on the basis of the preliminary sample, such that

$$
\begin{equation*}
P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right]>1-\alpha \tag{3,1}
\end{equation*}
$$

where $s_{n+1}^{2}$ is equal to $(1 / n) \sum_{i=1}^{n+1}\left(y_{i}-\bar{y}\right)^{2}$ and where
$y_{1}, y_{2}, \ldots, y_{n+1}$ is a random sample of size $n+1$, from a normal density with variance $\sigma^{2}$ 。

The tables in section 3.3 provide the sample size $n+1$, such that (3.1) is true, for

$$
\begin{aligned}
1-a & =.90, .95, .99 \\
m & =5(5) 20(10) 50(25) 150(50) 300(100) 500(250) 1000 . \\
\frac{s_{m}^{2}}{d} & =.33, .5, .67,1(1) 5(5) 20,30 .
\end{aligned}
$$

The only other known method for solving this problem is given in [8] which requires the use of Tchebycheff's inequality. It can be shown that the method presented in this chapter provides a significantly smaller sample size than does [8]. For some comparisons with [8], see Table 3.4.

### 3.2 Solution

Equation ( 3.1 ) may be written as

$$
\begin{aligned}
P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right] & =E_{n}\{P[(1-a)<v<(1+a) \mid n]\} \\
& =\int_{1}^{\infty} g(n) \int_{(1-a)}^{(1+a)} f_{1}(v) d v d n
\end{aligned}
$$

where $E_{n}$ is expectation with respect to $n ; a=\frac{d}{\sigma^{2}} ; v=\frac{s_{n+1}^{2}}{\sigma^{2}}$; $g(0)$ is the density of $n$, and $f_{1}(0)$ is the density of a chi-square variable divided by $n$, its degrees of freedom. We shall restrict $n$ such that $n \geq 1$. By definition

$$
f_{1}(v)=\frac{\left|\frac{n}{2}\right|^{\frac{n}{2}}}{\Gamma\left|\frac{n}{2}\right|} v^{(n / 2-1)} e^{-(n / 2) v}, \quad 0<v<\infty
$$

$$
=0,-\infty<v \leq 0 .
$$

In Chapter II it was shown that

$$
\int_{1-a}^{1+a} f_{1}(v) d v>\int_{1-a}^{1+a} f_{2}(v) d v \text {, for } a l 1 a>0, n \geq 1
$$

where

$$
f_{2}(v)=\frac{\sqrt{\Omega-1}}{2 \sqrt{\pi}} e^{-\frac{\sqrt{n-1}}{\sqrt{\pi}}|v-1|} \quad,-\infty<v<\infty
$$

and

$$
\int_{1-a}^{1+a} f_{2}(v) d v=1-e^{-\frac{\sqrt{n-1}}{\sqrt{\pi}}} a
$$

If a were known, we might let $n$ be equal to

$$
\frac{\pi \log ^{2} \alpha}{a^{2}}+1
$$

since in that case we would have

$$
\begin{aligned}
P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right] & >E_{n} \int_{1-a}^{1+a} f_{2}(v) d v \\
& =E_{n}(1-\alpha) \\
& =1-\alpha .
\end{aligned}
$$

Because a is assumed unknown let

$$
\begin{equation*}
n=\frac{\pi \log ^{2} \alpha}{d^{2}} k^{2} s_{m}^{4}+1 \tag{3.2}
\end{equation*}
$$

$$
=\frac{R^{2}}{a^{2}} u^{2}+1
$$

where $k$ is some constant, independent of $a$, such that

$$
E_{n}\left\{\int_{1-a}^{1+a} f_{2}(v) d v\right\}=1-\alpha,
$$

and where

$$
\begin{aligned}
R & =\frac{\sqrt{\pi}|\log \alpha|}{m-1} k \\
& =\frac{\sqrt{\pi} \log (1 / \alpha)}{m-1} k
\end{aligned}
$$

and

$$
u=\frac{(m-1) s_{m}^{2}}{\sigma^{2}}
$$

The density of $u$ is chi-square with $m-1$ degrees of freedom; that is

$$
f_{3}(u)=\frac{1}{2^{\frac{m-1}{2}} \Gamma\left(\frac{m-1}{2}\right)} u^{\frac{(m-1)}{2}-1} e^{-(1 / 2) u}, u \geq 0
$$

But

$$
u=\frac{a(n-I)^{I / 2}}{R}
$$

Thus

$$
g(n)=\frac{a^{(m-1) / 2}}{2^{\frac{m+1}{2}} R^{\frac{m-1}{2}} \Gamma\left(\frac{m-1}{2}\right)}, n \geq 1
$$

Therefore

$$
\begin{aligned}
& E_{n}\left\{\int_{1-a}^{1+a} f_{2}(v) d v\right\} \\
& =E_{n}\left\{1-e^{-\frac{\sqrt{n-1}}{\sqrt{\pi}} a}\right\} \\
& =1-\frac{\mathbf{a}^{(m-1) / 2}}{\frac{m+1}{2} \frac{m-1}{2} \Gamma\left(\frac{m-1}{2}\right)} \int_{1}^{\infty}(n-1)^{(m-1) / 4-1} .
\end{aligned}
$$

$$
\begin{aligned}
& =1-\frac{1}{2^{\frac{m-1}{2}} R^{\frac{m-1}{2}}\left(\frac{\sqrt{\pi}+2 R}{2 R \sqrt{\pi}}\right)^{\frac{m-1}{2}}} \\
& =1-\left[1+\frac{2}{(m-1)} \log (1 / \alpha) k\right]-\frac{(m-1)}{2}
\end{aligned}
$$

If we set

$$
\begin{equation*}
k=\frac{(m-1)}{2}\left[\frac{(1 / \alpha)^{\frac{2}{(m-1)}}-1}{\log (1 / \alpha)}\right] \tag{3.3}
\end{equation*}
$$

we have

$$
E_{n}\left\{\int_{1-a}^{1+a} f_{2}(v) d v\right\}=1-\alpha_{0}
$$

Thus if we substitute $k$ in Equation (3.3) into Equation (3.2) we get

$$
\begin{equation*}
n+1=\frac{\pi}{4}\left[(1 / a)^{\frac{2}{(m-1)}}-1\right]^{2}(m-1)^{2} \frac{s_{m}^{4}}{d^{2}}+2 \tag{3,4}
\end{equation*}
$$

We have proved that if the sample size $n+1$, given in Equation ( 3.4 ) is used for the second step sample, the following inequality is satisfied:

$$
P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right]>1-\alpha .
$$

The expected sample size in Equation (3.4) is

$$
\begin{equation*}
E_{n}(n+1)=\frac{\pi}{4}\left[(1 / \alpha)^{\frac{2}{(m-1)}}-1\right]^{2}(m-1)(m+1) \frac{\sigma^{4}}{d^{2}}+2 \tag{3,5}
\end{equation*}
$$

### 3.3 Sample Size Tables

The sample size $n+1$ as given in Equation (3.4) insures that ( 3.1 ) is true. To find the sample size, compute $s_{m}^{2} / d$, where $s_{m}^{2}$ is available from the preliminary sample of the procedure and $d$ is the desired allowable deviation from the true variance, and use Table 3.1, 3.2 , or 3.3 depending on the appropriate $1-\alpha$ level ( $m$ is the sample size on which $s_{m}^{2}$ is computed in the preliminary sample).

To find $n+1$ for values of $s_{m}^{2} / d$ other than those in Tables $3.1,3.2$, and 3.3 use Table 3.4 as follows. Compute $s_{m}^{4} / \mathrm{d}^{2}$, multiply by the entry in Table 3.4 which corresponds to the appropriate l-a level and $m$, and add 2 。

Table 3.5 shows some comparisons between the sample size given in (3.4) and the sample size obtained in [8]. The quantities tabled are

$$
h(m, \alpha)=\frac{n-1}{n^{\prime}-1}=\frac{\pi}{8} \alpha(m-3)(m-5)\left[(1 / \alpha)^{2 /(m-1)}-1\right]^{2} ; m \geq 6
$$

where $n+1$ is given in $(3.4)$ and $n$ is the sample size given in
[8]. It is noted that

$$
h(m, \alpha)=\frac{E(n-1)}{E\left(n^{\prime}-1\right)}
$$

It can be demonstrated that

$$
\begin{aligned}
& h(m, \alpha)<h\left(m, \alpha_{0}\right) \\
& <\lim _{m \rightarrow \infty} h\left(m, \alpha_{0}\right) \\
& =2 \pi e^{-2} \\
& \cong .85 \\
& \text { where } \alpha_{0}=\frac{m-5}{m-1} \quad(m-1) / 2 \text {. This shows that the sample size using }
\end{aligned}
$$

( 3.4 ) is never more than $85 \%$ of the sample size obtained by using the method in [8].

TABLE 3.1
Sample Size $n+1$ such that $P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right]>1-\alpha$


TABLE 3.2
Sample Size $n+1$ such that $P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right]>1-\alpha$

$$
1-\alpha=.95
$$




TABLE 3.3
Sample Size $n+1$ such that $P\left[\left|s_{n+1}^{2}-\sigma^{2}\right|<d\right]>1-\alpha$

$$
1-\alpha=.99
$$




TABLE 3.4
Entries are $(\pi / 4)\left[(1 / \alpha)^{2 /(m-1)}-1\right]^{2}(m-1)^{2}$


TABLE 3.5
Comparison of sample size: $n+1$ given in (2.3), $n$ ' given in [1]

$$
h(m, \alpha)=\frac{n-1}{n^{\prime}-1}=\frac{E(n-1)}{E\left(n^{\prime}-1\right)}
$$

| $m$ | $\alpha$ | .01 | .05 |
| ---: | :---: | :---: | :---: |
| 10 | .437 | .615 | .10 |
| 100 | .344 | .704 | .820 |
| 1000 | .334 | .705 | .832 |

## IV。 SAMPLE SIZE REQUIRED TO ESTIMATE THE PARAMETER IN THE POISSON DISTRIBUTION

### 4.1 Introduction

In this chapter some two-step procedures will be presented to estimate the Poisson Parameter within $d$ units with a specified confidence coefficient. The Poisson density is

$$
\begin{equation*}
P(x ; \lambda)=\frac{e^{-\lambda} \lambda^{x}}{x!}, x=0,1,2, \ldots 0 \tag{4,1}
\end{equation*}
$$

Let $m, d$, and $1-\varepsilon$ be specified in advance and let $x_{1}, x_{2}, \ldots, x_{m}$ be a preliminary sample of size $m$ from $P(0 ; \lambda)$. The problem is to determine $n$, the size of a second sample $y_{1}, y_{2}, \ldots, y_{n}$ from $P(0 ; \lambda)$, based on the values of the first sample, as $m, d$, and $l-\varepsilon$, such that

$$
\begin{equation*}
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right] \geq 1-\varepsilon \tag{4,2}
\end{equation*}
$$

where $\hat{\lambda}_{n}$ is some function of the second sample. If $n$ were fixed the maximum likelihood estimator of $\lambda$ is $\bar{y}_{n}$, the mean of the second sample。 Since $n$ is a random variable the maximum likelihood estimator of $\lambda$ depends on the density of $n$. In this chapter $n$ is not defined in explicit terms and the actual maximum likelihood estimator cannot easily be found. Theorem 4.1 shows that

$$
\begin{equation*}
\hat{\lambda}_{n}=\bar{y}_{n} \tag{4,3}
\end{equation*}
$$

is an unbiased estimator of $\lambda$. Equation (4.3) will be used as the definition of $\hat{\lambda}_{n}$ throughout the remainder of this chapter. Theorem 4.1 assumes $n$ is a proper random variable whose range consists of all the positive integers. Later in this chapter $n$ will be a continuous random variable, but by letting the second sample size be the next largest integer for fractional values of $n$, a new random variable is defined replacing $n$. In this case, equation ( 4.2 ) will be true if it was true for $n$.

Although $\bar{y}_{n}$ is computed on the basis of a random sample of size $n$, the unconditional distribution of the random variable $\bar{y}_{n}$ is independent of $n$. For convenience the subscript $n$ has been added. Theorem 4.1: Let $\hat{\lambda}_{n}=\bar{y}_{n}$, the mean of the second sample $y_{1}, y_{2}, \ldots, y_{n}$ from $P(0 ; \lambda)$. Then

$$
E\left(\hat{\lambda}_{n}\right)=\lambda
$$

## Proof:

By definition

$$
E\left(\hat{\lambda}_{n}\right)=E\left(\bar{y}_{n}\right) .
$$

Thus

$$
\begin{aligned}
E\left(\hat{\lambda}_{n}\right) & =E_{n}\left[E\left(\bar{y}_{n} \mid n\right)\right] \\
& =E_{n}[\lambda] \\
& =\lambda .
\end{aligned}
$$

It is further noted that $\bar{y}_{n}$ given the value of $n$ is the traditional estimator for $\lambda$. It is the minimum variance unbiased estimator and also is a squared error consistent estimator of $\lambda$.

Let

$$
\begin{align*}
& t_{1}\left(n ; \lambda_{2} d\right)=\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d \mid n\right] \\
&=\begin{array}{c}
{[n \lambda+n d]} \\
v=[n \lambda-n d]+1
\end{array}  \tag{4.4}\\
& e^{-n \lambda(n \lambda)^{v}} \\
& v!
\end{align*}
$$

where [k] means the integral value of $k$. The density of $n \bar{y}_{n}$ given $n$ is $P(0 ; n \lambda)$, $i, e_{0}, n \bar{y}_{n}$ given $n$ is Poisson with parameter $n \lambda_{0}$ Thus

$$
\begin{align*}
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right] & =E_{n}\left\{\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d \mid n\right]\right\} \\
& =E_{n}\{\operatorname{Pr}[|n \bar{y}-n \lambda|<n d \mid n]\} \\
& =E_{n}\left\{t_{1}(n ; \lambda, d)\right\} . \tag{4,5}
\end{align*}
$$

In Section 2 we shall prove certain monotonic properties of $t_{1}$ and a generalization of $t_{l}$ which will lead to a determination of $n$, the size of the second sample.

Sections 3, 4, 5 and 6 present various methods for determining $n$. The procedure in Section 3 is the easiest to use and demonstrates the basic technique employed in this chapter to determine $n$. Section 4 gives a second solution which leads to a smaller sample size $n$ but the confidence coefficient is not predetermined. This difficulty is removed by a similar solution in Section 5. This solution allows for a preassigned confidence coefficient and reduces sample size compared with the solution in Section 3 but is more difficult to use. The basic solution is
further generalized in Section 6 but the confidence coefficient is not preassigned.

In Section 8 some comparisons of the results of Sections 3 and 5 are made with the solution to this problem which can be obtained by Birmbaum and Healy's [8] method using Tchebycheff's inequality. These comparisons show a significant reduction in sample size.

To apply the results of this chapter some examples are given in Section 9 along with sample size graphs for procedures developed in Sections 3 and 5. Also Section 7 shows how to use this chapter's methods to estimate the mean in a Poisson stochastic process.

In the remainder of this chapter the following definitions for $z$, $\bar{x}_{m}$, and $H$ will be used:

$$
z=m \bar{x}_{m}
$$

where $\bar{x}_{m}$ is the mean of the preliminary sample, and

$$
\begin{equation*}
H(c ; z)=\sum_{v=0}^{z} \frac{e^{-c z}(c z)^{v}}{v!} \tag{4,6}
\end{equation*}
$$

Also,

$$
a \sim b
$$

will mean that $a$ and $b$ have the same sign。

### 4.2 Monotonic Properties

Equation (4.4) can be rewritten as

$$
\begin{equation*}
t_{1}(n ; \lambda, d)=\sum_{v=[n \lambda-n d]+1}^{[n \lambda+n d]} P(v ; n \lambda) \tag{4,7}
\end{equation*}
$$

where $P(v ; n \lambda)$ is defined by (4.1). Consider

$$
\begin{aligned}
h_{1}(v) & =P(v+1 ; n \lambda) / P(v ; n \lambda) \\
& =n \lambda /(v+1) \\
& \begin{cases}>1, & v<n \lambda-1 \\
<1, & v>n \lambda-1 .\end{cases}
\end{aligned}
$$

Thus, for integral values of $v$, the function $P(v ; n \lambda)$ is monotonic increasing in $v$ for $v>[n \lambda]$ and is monotonic decreasing for $v>[n \lambda]$. Thus the mode of $P(0 ; n \lambda)$ is $[n \lambda]$. If $n \lambda$ is an integer, then $P(0 ; n \lambda)$ is bimodal with modes at $n \lambda-1$ and $n \lambda$ since

$$
\begin{aligned}
P(n \lambda-1 ; n \lambda) & =\frac{e^{-n \lambda}(n \lambda)^{n \lambda-1}}{(n \lambda-1)!} \\
& =\frac{e^{-n \lambda}(n \lambda)^{n \lambda}}{(n \lambda)!} \\
& =P(n \lambda ; n \lambda)
\end{aligned}
$$

Let

$$
\begin{equation*}
t(n ; \lambda, d)=\int_{n \lambda-n d+.5}^{n \lambda+n d-.5} f(v ; n \lambda) d v \tag{4.8}
\end{equation*}
$$

where

$$
f(v ; \lambda)=(1-v+[v]) P([v] ; \lambda)+(v-[v]) P([v]+1 ; \lambda)
$$

and it is assumed that nd is greater than or equal to 2. Otherwise let $t$ be zero. The following two theorems show the relationship of
$t(n ; \lambda, d)$ to $t_{1}(n ; \lambda, d)$ and some monotonic properties of $t$ essential to the determination of $n$.

Theorem 4.2: Let $t_{1}$ and $t$ be defined by (4.7) and (4.8) respectively. Then

$$
E_{n}\{t(n ; \lambda, d)\} \geq 1-\varepsilon
$$

implies equation (4.2) is true, i.e.,

$$
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right] \geq 1-\varepsilon_{0}
$$

Proof: From (4.8) we observe that $f(v ; n \lambda)$ corsists of straight lines joining the values of $P(v ; n \lambda)$ at adjacent integral values of $v$. Therefore

$$
\int_{k}^{k+1} f(v ; n \lambda) d v=.5\{P(k ; n \lambda)+P(k+1 ; n \lambda)\}
$$

for integral values of $k$. Also,

$$
\begin{gathered}
\int_{n \lambda-n d+.5}^{[n \lambda-n d+1.5]} f(v ; n \lambda) d v=0.5\{[n \lambda-n d+1.5]-(n \lambda-n d+.5)\}\{f(n \lambda-n d+.5 ; n \lambda) \\
\\
+P([n \lambda-n d+1.5] ; n \lambda)\}
\end{gathered}
$$

Thus, if $[n \lambda-n d+1,5]=[n \lambda-n d+1]$, then

$$
\begin{aligned}
& \int_{n \lambda-n d+.5}^{[n \lambda-n d+1.5]} f(v ; n \lambda) d v<.5(.5)\{f(n \lambda-n d+05 ; n \lambda)+P([n \lambda-n d+1] ; n \lambda)\} \\
& \leq 25\{f([n \lambda-n d+1] ; n \lambda)+P([n \lambda-n d+1] ; n \lambda)\} \\
& =5 \mathrm{P}([n \lambda-n d+1] ; n \lambda) \text {, }
\end{aligned}
$$

since $f(v ; n \lambda)$ is monotonic increasing for $v$ less than $[n \lambda]$. Otherwise $[n \lambda-n d+1.5]=[n \lambda-n d+2]$, which implies

$$
\begin{aligned}
\int_{n \lambda-n d+05}^{[n \lambda-n d+1,5]} f(v ; n \lambda) d v & \leq \int_{[n \lambda-n d+1]}^{[n \lambda-n d+2]} f(v ; n \lambda) d v \\
& <05\{P([n \lambda-n d+2] ; n \lambda)+P([n \lambda-n d+1] ; n \lambda)\}
\end{aligned}
$$

Similarly, if $[n \lambda+n d=5]=[n \lambda+n d]$, then

$$
\int_{[n \lambda+n d m .5]}^{n \lambda+n d-05} f(v ; n \lambda) d v=.5\{(n \lambda+n d-05)-[n \lambda+n d-.5]\}\{P([n \lambda+n d-.5] ; n \lambda)
$$

$$
\begin{aligned}
& \quad+f(n \lambda+n d-.5 ; n \lambda)\} \\
& <\quad .5(.5)\{P([n \lambda+n d] ; n \lambda)+f(n \lambda+n d-.5 ; n \lambda)\} \\
& \leq \\
& =025\{P([n \lambda+n d] ; n \lambda)+f([n \lambda+n d] ; n \lambda)\} \\
& = \\
& .5 P([n \lambda+n d] ; n \lambda),
\end{aligned}
$$

since $f(v ; n \lambda)$ is monotonic decreasing for $v$ greater than [ $n \lambda$ ]. In the case where $[n \lambda+n d=05]=[n \lambda+n d-1]$, then

$$
\begin{aligned}
\int_{[n \lambda+n d-0.5]}^{n \lambda+n d-.5} f(v ; n \lambda) d v & \leq \int_{[n \lambda+n d-1]}^{[n \lambda-n d]} f(v ; n \lambda) d v \\
& <.5\{P([n \lambda+n d-1] ; n \lambda)+P([n \lambda+n d] ; n \lambda)\} .
\end{aligned}
$$

Therefore

$$
\begin{aligned}
t(n ; \lambda, d) & =\int_{n \lambda-n d+.5}^{n \lambda+n d-.5} f(v ; n \lambda) d v \\
& <\sum_{v=[n \lambda-n d]+1}^{[n \lambda+n d]} P(v ; n \lambda) \\
& =t_{1}(n ; \lambda, d) .
\end{aligned}
$$

Hence

$$
\begin{aligned}
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right] & =E_{n}\left\{t_{1}(n ; \lambda, d)\right\} \\
& >E_{n}\{t(n ; \lambda, d)\}
\end{aligned}
$$

If nd is less than 2, then the inequality is still true which completes the proof.

Theorem 4.3: Let $t(n ; \lambda, d)$ be defined by (4.8). Then assuming $n d \geq 2$, we have
(a) $\frac{\partial t}{\partial d}>0$
(b) $\frac{\partial t}{\partial \lambda}<0$
(c) $\frac{\partial t}{\partial n}>0$ 。

Proof:
(a) Differentiating $t$ with respect to $d$ we obtain

$$
\begin{aligned}
\frac{\partial t}{\partial d} & =n f(n \lambda+n d-.5 ; n \lambda)+n f(n \lambda-n d+.5 ; n \lambda) \\
& >0 .
\end{aligned}
$$

(b) Also

$$
\frac{\partial t}{\partial \lambda}=n f(n \lambda+n d-.5 ; n \lambda)-n f(n \lambda-n d+.5 ; n \lambda)+\int_{n \lambda-n d+.5}^{n \lambda+n d-.5} \frac{\partial t}{\partial \lambda}(v ; n \lambda) d v,
$$

where

$$
\frac{\partial f}{\partial \lambda}(v ; n \lambda)=(1-v+[v]) \frac{\partial P}{\partial \lambda}([v] ; n \lambda)+(v-[v]) \frac{\partial}{\partial \lambda} p([v]+1 ; n \lambda) .
$$

But

$$
\begin{aligned}
\frac{\partial P}{\partial \lambda}(v ; n \lambda) & =\left(\frac{v}{\lambda}-n\right) P(v ; n \lambda) \\
& =n P(v-1 ; n \lambda)-n P(v ; n \lambda) .
\end{aligned}
$$

Thus

$$
\begin{aligned}
\frac{\partial f}{\partial \lambda}(v ; n \lambda)= & n(1-v+[v]) P([v]-1 ; n \lambda)+n\{-(1-v+[v])+ \\
& (v-[v])\} P([v] ; n \lambda) \\
& -n(v-[v]) P([v]+1 ; n \lambda) \\
= & n\{(1-v+[v]) P([v]-1 ; n \lambda)+(v-[v]) P([v] ; n \lambda)\}-n\{(1-v+[v]) P([v] ; n \lambda) \\
& +(v-[v]) P([v]+1 ; n \lambda)\} \\
= & n f(v-1 ; n \lambda)-n f(v ; n \lambda) .
\end{aligned}
$$

## Therefore

$$
\begin{aligned}
& \frac{\partial t}{\partial \lambda}=n f(n \lambda+n d-.5 ; n \lambda)-n f(n \lambda-n d+.5 ; n \lambda)+n \int_{n \lambda-n d+.5}^{n \lambda+n d-.5}\{f(v-1 ; n \lambda)-f(v ; n \lambda)\} d v
\end{aligned}
$$

$$
\begin{aligned}
& <0 \text {, }
\end{aligned}
$$

since $f(v ; n \lambda)$ is monotonic increasing for $v$ in the interval
( $n \lambda-n d=0, n \lambda-n d+5$ ) and is monotonic decreasing in the interval ( $n \lambda+n d-1,5, n \lambda+n d-5$ ) 。
(c) Finally, differentiating with respect to $n$ we have

$$
\frac{\partial t}{\partial n}=(\lambda+d) f(n \lambda+n d-.5 ; n \lambda)-(\lambda-d) f(n \lambda-n d+.5 ; n \lambda)+\int_{n \lambda-n d+.5}^{n \lambda+n d-.5} \frac{\partial}{\partial n} f(v ; n \lambda) d v,
$$

where

$$
\frac{\partial f}{\partial n}(v ; n \lambda)=\lambda f(v-1 ; n \lambda)-\lambda f(v ; n \lambda)
$$

since $n$ appears everywhere $\lambda$ does in $f(v ; n \lambda)$ and from part (b)

$$
\frac{\partial f}{\partial \lambda}(v ; n \lambda)=n f(v-1 ; n \lambda)-n f(v ; n \lambda) .
$$

Thus

$$
\frac{\partial t}{\partial n}=(\lambda+d) f\left(n \lambda+n d-{ }_{0} 5 ; n \lambda\right)-(\lambda-d) f(n \lambda-n d+.5 ; n \lambda)+\lambda \int_{n \lambda-n d+.5}^{f(v ; n \lambda) d v-\lambda} \begin{gathered}
n \lambda-n d-.5
\end{gathered} \begin{gathered}
n \lambda+n d-.5 \\
f(v ; n \lambda) d v . \\
n \lambda+n d-1.5
\end{gathered}
$$

Let

$$
\begin{aligned}
& L=n \lambda, \\
& D=n d,
\end{aligned}
$$

and

$$
S_{ \pm}= \pm(L \pm D) f\left(L \pm D \bar{\Psi}_{n} 5 ; L\right) \overline{+} \int_{L \pm D-1 \bar{\Psi}_{0} 5}^{L+D \bar{\Psi}_{0} 5} f(v ; L) d v
$$

Therefore

$$
\begin{aligned}
\frac{\partial t}{\partial n} & \sim n \frac{\partial t}{\partial n} \\
& =s_{+}+s_{-} .
\end{aligned}
$$

Let $\quad s_{ \pm}=L \pm D \mp .5-[L \pm D \mp .5]$.

For convenience the subscript on $s$ will be deleted below. Hence

$$
\begin{aligned}
& S_{ \pm}= \pm(L \pm D)\{(1-s) P([L \pm D \mp \cdot 5] ; L)+s P([L \pm D \overline{+} .5]+1 ; L)\} \\
& \overline{+} L\{5(1-s) f(L \pm D-1 \mp .5 ; L)+.5(1-s) P([L \pm D+5] ; L) \\
& \left.+.5 \mathrm{sP}\left(\left[\mathrm{~L}+\mathrm{D} \overline{\mathrm{~F}}_{0} 5\right] ; \mathrm{L}\right)+.5 \mathrm{sf}(\mathrm{~L}+\mathrm{DF} .5 ; \mathrm{L})\right\} \\
& =\mp_{0} 5 L(1-s)^{2} P\left(\left[L \pm D \bar{H}_{0} 5\right]-1 ; L\right) \\
& \pm\{(L \pm D)(1-s)-5 L(1-s) s-5 L-5 L s(1-s)\} P([L \pm D+.5] ; L) \\
& \pm\left\{(L \pm D) s-.5 L s^{2}\right\} P([L \pm D+.5]+1 ; L) \text { 。 }
\end{aligned}
$$

But

$$
\begin{aligned}
P([L \pm D \bar{\dagger}, 5]-1 ; L) & =\frac{[L \pm D \overline{+}, 5]}{L} P([L \pm D \bar{\mp}, 5] ; L) \\
& =\frac{(L \pm D \overline{+}, 5-s)}{L} P([L \pm D \overline{+}, 5] ; L) .
\end{aligned}
$$

Hence

$$
\begin{aligned}
S_{ \pm} & \sim 2 S \\
= & \pm\left\{-\left(L \pm D F_{0} 5-s\right)(1-s)^{2}+2(L \pm D)(1-s)-2 L s(1-s)-L\right\} P([I \pm D \overline{+}, 5] ; L) \\
& \pm\left\{2(L \pm D) s-L s^{2}\right\} P\left(\left[L+D F_{0} 5\right]+1 ; L\right) \\
\sim & \left(2 S_{ \pm}\right) L / P\left(\left[L \pm D \Psi_{0} 5\right]+1 ; L\right) \\
= & \pm\left\{\left(-2 s+s^{2}\right) L \pm\left(1-s^{2}\right) D+(s \pm .5)(1-s)^{2}\right\}\left(L \pm D \bar{t}_{0} 5+1-s\right) \\
& \pm\left\{\left(2 s-s^{2}\right) L \pm 2 s D\right\} L
\end{aligned}
$$

$$
\begin{aligned}
= & L D+\left(.5 \overline{+s} \pm s^{2}\right) L \pm\left(1-s^{2}\right) D^{2} \\
& +\left\{1 \overline{+s}-(3 \overline{+1}) s^{2}+2 s^{3}\right\} D \pm(s \pm .5)(1-s)^{2}(1-s \overline{+} .5) \\
= & L D+\{.5 \overline{+s}(1-s)\} L \pm(1-s)(1+s) D^{2}+s(1-s) D \\
& +(1-s)\left(1 \overline{+s}-2 s^{2}\right) D \pm(1-s)^{2}(s \pm .5)\{1-(s \pm .5)\} \\
> & L D+.25 L m D^{2}-.25 \\
= & (L-D) D+.25(L-1) \\
> & 0
\end{aligned}
$$

Therefore

$$
\begin{aligned}
\frac{\partial t}{\partial n} & \sim s_{+}+s_{-} \\
& >0
\end{aligned}
$$

which completes the proof.
4.3 One Point Solution

Let $m, d, \alpha$ and $\beta$ be specified in advance and observe a preliminary sample $x_{1}, x_{2}, \ldots x_{m}$. Define $n_{1}$ such that

$$
\begin{equation*}
t\left(n_{1} ; \lambda, d\right)=1-\alpha \tag{4.9}
\end{equation*}
$$

(note Figure 4.1).


Figure 4.1

Determine the value $n$ of the random variable, the size of the second sample, such that

$$
\begin{equation*}
t\left(n ; c \bar{x}_{m} ; d\right)=1-\alpha \tag{4,10}
\end{equation*}
$$

where $t$ is defined by $(4,8), H$ by $(4,6)$, and $C$ is defined by

$$
\begin{equation*}
H(c ; z)=\beta . \tag{4.11}
\end{equation*}
$$

From (4.9) and (4.10) we have

$$
\begin{equation*}
t\left(n_{1} ; \lambda, d\right)=t\left(n ; c \bar{x}_{m}, d\right) \tag{4,12}
\end{equation*}
$$

Thus, using Theorem 4.3 (c) and (b), we obtain

$$
\begin{align*}
\operatorname{Pr}\left(n>n_{1}\right) & =\operatorname{Pr}\left[t(n ; \lambda, d)>t\left(n_{1} ; \lambda, d\right)\right] \\
& =\operatorname{Pr}\left[t(n ; \lambda, d)>t\left(n ; c \bar{x}_{m}, d\right)\right] \\
& =\operatorname{Pr}\left(\lambda<c \bar{x}_{m}\right) \tag{4,13}
\end{align*}
$$

Theorem 4.4: Let $H(c ; z)$ be defined as in (4.6). Then

$$
\operatorname{Pr}\left(\lambda<c \bar{x}_{m}\right)=1-H(c ; z) .
$$

Proof: The random variable $m \bar{x}_{m}$ is distributed as Poisson with mean m. Therefore

$$
\operatorname{Pr}\left(\bar{x}_{m}>\lambda_{1-\beta}\right)=I-\beta,
$$

where $\lambda_{1-\beta}$ is defined such that

$$
\sum_{v=0}^{m \lambda} \frac{e^{-m \lambda}(m \lambda)^{v}}{v!}=\beta .
$$

This implies that

$$
\operatorname{Pr}\left(\lambda<\lambda^{\prime}\right)=1-\beta
$$

where $\lambda^{\prime}$ is defined by

$$
\underset{v=0}{\sum_{\mathrm{x}}^{\mathrm{m}}} \frac{e^{-m \lambda^{\prime}}\left(m \lambda^{\prime}\right)^{v}}{v!}=\beta .
$$

Hence

$$
\begin{aligned}
\operatorname{Pr}\left(\lambda<c \bar{x}_{m}\right) & =1-\sum_{v=0}^{\bar{x}_{m}} \frac{e^{-m c \bar{x}_{m}\left(m c \bar{x}_{m}\right)^{v}}}{v!} \\
& =1-H(c ; z),
\end{aligned}
$$

since $\lambda^{\prime}$ and $\beta$ are in a one to one correspondence with each other, thus completing the proof.

Therefore, from $(4.11),(4.13)$ and Theorem 4.4. we obtain

$$
\begin{align*}
\operatorname{Pr}\left(n>n_{1}\right) & =1-H(c ; z) \\
& =1-\beta \tag{4.14}
\end{align*}
$$

From Theorem 4.3 (c) we obtain

$$
t(n ; \lambda, d) \geq \begin{cases}0 & 0<n<n_{1} \\ t\left(n_{1} ; \lambda, d\right), & n \geq n_{1}\end{cases}
$$

which is depicted by the shaded area in Figure 4.1. Thus, by choosing n as defined in $(4,10)$, we obtain

$$
\begin{aligned}
E_{n}[t(n ; \lambda, d)] & >0 \cdot \operatorname{Pr}\left(0<n<n_{1}\right)+t\left(n_{1} ; \lambda, d\right) \operatorname{Pr}\left(n \geq n_{1}\right) \\
& =(1-\alpha)(1-\beta)
\end{aligned}
$$

by (4.9) and (4.14). From Theorem 4.2 we conclude that

$$
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right]>1-\varepsilon,
$$

where

$$
1-\varepsilon=(1-\alpha)(1-\beta)
$$

Graphs are given in Section 9 to find $n$ for various values of $m, d, \bar{x}_{m}$ and $l-\varepsilon$ 。
4.4 Two Point Solution:

Let $m, d, \alpha, \beta$, and $\delta$ be preassigned and let $x_{1}, x_{2, \ldots,}, m$ be a preliminary sample. Define $n_{1}, n_{2}$ such that
(a) $t\left(n_{1} ; \lambda, d\right)=\gamma(1-\alpha)$
(b) $t\left(n_{2} ; \lambda, d\right)=1-\alpha$
where $\gamma$ is calculated by (4.19) (note Figure 4.2).


Figure 4.2

Define $c_{1} ; c_{2}$ such that
(a) $H\left(c_{1} ; z\right)=(1-\delta) \beta$
(b) $H\left(c_{2} ; z\right)=B$.

Determine the value $n$ of the random variable, the size of the second sample, such that

$$
\begin{equation*}
t\left(n ; c_{2} \bar{x}_{m}, d\right)=1-\alpha \tag{4,18}
\end{equation*}
$$

and calculate $\gamma$ using

$$
\begin{equation*}
t\left(n ; c_{1} \bar{x}_{m}, d\right)=\gamma(1-\alpha) \tag{4.19}
\end{equation*}
$$

Proceeding as in Section 3 and using (4.16), (4.18) and (4.19) we have

$$
t\left(n_{i} ; \lambda, d\right)=t\left(n ; c_{i} \bar{x}_{m}, d\right), \quad i=1,2 .
$$

Thus, by Theorem 4.3 (c) and (b), we obtain

$$
\begin{aligned}
\operatorname{Pr}\left(n>n_{i}\right) & =\operatorname{Pr}\left[t(n ; \lambda, d)>t\left(n_{i} ; \lambda, d\right)\right] \\
& =\operatorname{Pr}\left[t(n ; \lambda, d)>t\left(n ; c_{i} \bar{x}_{m} d\right)\right] \\
& =\operatorname{Pr}\left(\lambda<c_{i} \bar{x}_{m}\right), \quad i=1,2 .
\end{aligned}
$$

Hence, by (4.17) and Theorem 4.4, we have

$$
\begin{equation*}
\operatorname{Pr}\left(n>n_{1}\right)=1-(1-\delta) \beta \tag{4.20}
\end{equation*}
$$

and

$$
\operatorname{Pr}\left(n>n_{2}\right)=1-\beta
$$

which implies that

$$
\begin{align*}
P_{r}\left(n_{1}<n<n_{2}\right) & =1-(1-\delta) \beta-(1-\beta) \\
& =\delta \beta \tag{4.21}
\end{align*}
$$

From Theorem 4.3 (c) we obtain

$$
t(n ; \lambda, d) \geq\left\{\begin{array}{cl}
0 & 0<n \leq n_{1} \\
t\left(n_{1} ; \lambda, d\right), & n_{1}<n \leq n_{2} \\
t\left(n_{2} ; \lambda, d\right), & n>n_{2}
\end{array}\right.
$$

This is represented by the shaded portion in Figure 4.2. Thus by (4.16), (4.20) and (4.21) we conclude that

$$
\begin{aligned}
E_{n}[t(n ; \lambda, d)] \geq 0 \cdot \operatorname{Pr}\left(0<n \leq n_{1}\right) & +t\left(n_{1} ; \lambda, d\right) \operatorname{Pr}\left(n_{1}<n \leq n_{2}\right) \\
& +t\left(n_{2} ; \lambda, d\right) \operatorname{Pr}\left(n>n_{2}\right) \\
= & \gamma(1-\alpha) \delta \beta+(1-\alpha)(1-\beta) \\
= & (1-\alpha)(1-\beta+\gamma \delta \beta) .
\end{aligned}
$$

By Theorem 4.2 this implies

$$
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right]>1-\varepsilon
$$

where

$$
1-\varepsilon=(1-\alpha)(1-\beta+\gamma \delta \beta) .
$$

The value of $n$ is determined by (4.18) and it is observed that this would result in the same value as determined in the one point solution assuming $\alpha$ and $\beta$ are the same. The value of $1-\varepsilon$, the confidence coefficient, is increased however. The amount of increase depends upon $\gamma$ which is calculated after the value of $n$ is detemmined. Thus the confidence coefficient is not preassigned as would normally be desired. Section 5 shows one way to avoid this difficulty.

### 4.5 Two Point Solution with Preassigned Confidence Coefficient

Specify $m, d$, and $\alpha, \beta, \gamma, \delta$ such that $(1-\alpha)(1-\beta+\gamma \beta)$
equals some prescribed value, say $1-\varepsilon_{0}$ Let $x_{1}, x_{2}, 0,0, x_{m}$ be a preliminary sample, Define $n_{1}, n_{2}$ as in (4.16) and $c_{1}, c_{2}$ as in (4.17). Determine the values $n^{\prime} n^{\prime \prime}$ of two random variables such that
(a) $t\left(n^{\prime} ; c_{1} \bar{x}_{m}, d\right)=\gamma(1-\alpha)$
(b) $t\left(n^{\prime \prime} ; c_{2} \bar{x}_{m}, d\right)=(1-\alpha)$.

From (4.16) and (4.22) we have
(a) $t\left(n_{1} ; \lambda, d\right)=t\left(n^{\prime} ; c_{1} \bar{x}_{m}, d\right)$
(b) $t\left(n_{2} ; \lambda, d\right)=t\left(n^{\prime \prime} ; c_{2} \bar{x}_{m}, d\right)$

Thus, by Theorem $4.3(c)$ and (b) we obtain
(a) $\operatorname{Pr}\left(n^{\prime}>n_{1}\right)=\operatorname{Pr}\left[t\left(n^{\prime} ; \lambda, d\right)>t\left(n_{1} ; \lambda, d\right)\right]$

$$
=\operatorname{Pr}\left[t\left(n^{\prime} ; \lambda_{1} d\right)>t\left(n^{\prime} ; c_{1} \bar{x}_{m}, d\right)\right]
$$

$$
=\operatorname{Pr}\left(\lambda<c_{1} \bar{x}_{m}\right)
$$

(b) $\operatorname{Pr}\left(n^{\prime \prime}>n_{2}\right)=\operatorname{Pr}\left[t\left(n^{\prime \prime} ; \lambda, d\right)>t\left(n_{2} ; \lambda, d\right)\right]$

$$
\begin{aligned}
& =\operatorname{Pr}\left[t\left(n^{\prime \prime} ; \lambda, d\right)>t\left(n^{\prime \prime} ; c_{2} \bar{x}_{m}, d\right)\right] \\
& =\operatorname{Pr}\left(\lambda<c_{2} \bar{x}_{m}\right) .
\end{aligned}
$$

Hence, by (4.17) and Theorem 4.4, we have
(a) $\operatorname{Pr}\left(n^{\prime}>n_{1}\right)=1-(1-\delta) \beta$
(b) $P_{r}\left(n^{\prime \prime}>n_{2}\right)=1-\beta$.

Now choose the value $n$ of the random variable, the size of the second sample, such that

$$
\begin{equation*}
n=\max \left(n^{\prime}, n^{\prime \prime}\right) \tag{4.25}
\end{equation*}
$$

Therefore, by (4.24), we have
(a) $\operatorname{Pr}\left(n>n_{1}\right) \geq \operatorname{Pr}\left(n^{\prime}>n_{1}\right)$

$$
\begin{equation*}
=1-(1-\delta) \beta \tag{4,26}
\end{equation*}
$$

(b) $\operatorname{Pr}\left(\mathrm{n}>\mathrm{n}_{2}\right) \geq \operatorname{Pr}\left(\mathrm{n}^{\prime \prime}>\mathrm{n}_{2}\right)$

$$
=1-\beta .
$$

Define $n_{1}^{\prime}$ and $n_{2}^{\prime}$ such that
(a) $\operatorname{Pr}\left(n>n_{1}^{\prime}\right)=1-(1-\delta) \beta$
(b) $\operatorname{Pr}\left(n>n_{2}^{\prime}\right)=1-\beta$,
which implies that

$$
\begin{equation*}
\operatorname{Pr}\left(n_{1}^{\prime}<n \leq n_{2}^{\prime}\right)=\delta \beta \tag{4,28}
\end{equation*}
$$

comparing (4.26) and (4.27) it is seen that
(a) $n_{1}^{\prime} \geq n_{1}$
(b) $n_{2}^{\prime} \geq n_{2}$,
which, by (4.16) and Theorem 4.3(c), implies
(a) $t\left(n_{1}^{\prime} ; \lambda, d\right) \geq t\left(n_{1} ; \lambda, d\right)$

$$
\begin{equation*}
=\gamma(1-\alpha) \tag{4.29}
\end{equation*}
$$

(b) $t\left(n_{2}^{\prime} ; \lambda, d\right) \geq t\left(n_{2} ; \lambda, d\right)$

$$
=1-\alpha_{0}
$$

From Theorem 4.3(c) we have

$$
t(n ; \lambda, d) \geq \begin{cases}0 & , \quad 0<n \leq n_{1}^{\prime} \\ t\left(n_{1}^{\prime} ; \lambda, d\right), & n_{1}^{\prime}<n \leq n_{2}^{\prime} \\ t\left(n_{2}^{\prime} ; \lambda, d\right), & n>n_{2}^{\prime}\end{cases}
$$

Thus, by $(4.27 b),(4.28)$, and $(4.29)$, we obtain

$$
\begin{aligned}
E_{n}[t(n ; \lambda, d)] \geq 0 \cdot \operatorname{Pr}\left(0<n \leq n_{1}^{\prime}\right) & +t\left(n_{1}^{\prime} ; \lambda, d\right) \operatorname{Pr}\left(n_{1}^{\prime}<n \leq n_{2}^{\prime}\right) \\
& +t\left(n_{2}^{\prime} ; \lambda, d\right) \operatorname{Pr}\left(n>n_{2}^{\prime}\right) \\
& \geq r(1-\alpha) \delta \beta+(1-\alpha)(1-\beta)
\end{aligned}
$$

```
=(1-\alpha)(1-\beta+\gamma\delta\beta)
=1-\varepsilon.
```

Thus, by Theorem 4.2 , if $n$ is determined by (4.25) we conclude that

$$
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right]>1-\varepsilon,
$$

a predetermined confidence coefficient.
This solution is more difficult to work with but yields smaller second sample sizes. In Section 9 graphs are given for various values of $m, d, \bar{x}_{m}$ and $1-\varepsilon$. Thus, in using the graphs an experimenter does not have to specify $\alpha, \beta, \gamma$ and $\delta$.

### 4.6 The $k$ Point Solution

Let $m, d, k, \beta_{i}(i=1,2, \ldots \ldots, k)$ and $\gamma_{1}$ be specified in advance and define $c_{i}(i=1,2, \ldots, k)$ such that

$$
\begin{equation*}
H\left(c_{j} ; z\right)=\beta_{i}, \quad i=1,2, \ldots, k \tag{4.30}
\end{equation*}
$$

Determine the value $n$ of the random variable, the size of the second sample, by ( $a$ ), and $\gamma_{i}(i=2,3, \ldots, k)$ such that
(a) $t\left(n ; c_{1} \bar{x}_{m}, d\right)=\gamma_{1}$
(b) $t\left(n ; c_{i} \bar{x}_{m}, d\right)=\gamma_{i}, \quad i=2,3, \ldots, k$ 。

Further define $n_{i}(i=1,2, \ldots, k)$ such that

$$
\begin{equation*}
t\left(n_{i} ; \lambda, d\right)=\gamma_{i}, i=1,2, \ldots, k \tag{4.32}
\end{equation*}
$$

(note Figure 4.3).


Figure 4.3

From (4.31) and (4.32) it is seen that

$$
t\left(n_{i} ; \lambda, d\right)=t\left(n ; c_{i} \bar{x}_{m}, d\right), \quad i=1,2, \ldots, k
$$

Thus, by Theorem 4.3(c) and (b), we obtain

$$
\begin{aligned}
\operatorname{Pr}\left(n>n_{i}\right) & =\operatorname{Pr}\left[t(n ; \lambda, d)>t\left(n_{i} ; \lambda, d\right)\right] \\
& =\operatorname{Pr}\left[t(n ; \lambda, d)>t\left(n ; c_{i} \bar{x}_{m}, d\right)\right] \\
& =\operatorname{Pr}\left(\lambda<c_{i} \bar{x}_{m}\right), \quad i=1,2, \ldots, k
\end{aligned}
$$

Hence, by (4.30) and Theorem 4.4; we have

$$
\operatorname{Pr}\left(n>n_{i}\right)=1-\beta_{i}, \quad i=1,2, \ldots, k
$$

which implies that

$$
\begin{equation*}
\operatorname{Pr}_{r}\left(n_{i}<n \leq n_{i+1}\right)=\beta_{i+1}-\beta_{i}, \quad i=1,2, \ldots, k \tag{4.33}
\end{equation*}
$$

where we define

$$
n_{k+1}=\infty
$$

and

$$
\beta_{k+1}=1 .
$$

From Theorem 4.3(c) it is seen that

$$
t(n ; \lambda, d) \geq\left\{\begin{array}{cl}
0, & 0<n \leq n_{2} \\
t\left(n_{1} ; \lambda, d\right), & n_{1}<n \leq n_{2} \\
t\left(n_{2} ; \lambda, d\right), & n_{2}<n \leq n_{3} \\
\cdots & \\
t\left(n_{k} ; \lambda, d\right), & n>n_{k}
\end{array}\right.
$$

(note the shaded portion in Figure 4.3).
Thus, by (4.32) and (4.33) we have

$$
\begin{aligned}
E_{n}[t(n ; \lambda, d)] & \geq \sum_{i=1}^{k} t\left(n_{i} ; \lambda, d\right) \operatorname{Pr}\left(n_{i}<n \leq n_{i+1}\right) \\
& =\sum_{i=1}^{k} \gamma_{i}\left(\beta_{i+1}-\beta_{i}\right) .
\end{aligned}
$$

By Theorem 4.2 this implies that

$$
\operatorname{Pr}\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right]>1-\varepsilon
$$

where

$$
1-\varepsilon=\sum_{i=1}^{k} \gamma_{i}\left(\beta_{i+1}-\beta_{i}\right) .
$$

Since the $\gamma_{i}(i=2,3, \ldots, k)$ are not determined in advance, the confidence coefficient is not predetermined in this extension of the two point solution in Section 4 .
4.7 The Poisson Stochastic Process

In a Poisson stochastic process we have a counting process $\{N(t), t \geq 0\}$ such that $\{N(t), t \geq 0\}$ has stationary independent increments and

$$
\operatorname{Pr}[N(t)-N(s)=k]=\frac{e^{-\lambda(t-s)}[\lambda(t-s)]^{k}}{k!}, \quad k=0,1,2, \ldots \ldots
$$

To apply the method of this chapter to estimate $\lambda$, pick some constant $\tau>0$. Then let the random variables $X_{i}, Y_{j}, i=1,2, \ldots, j=1,2, \ldots$ be defined such that

$$
X_{i}=N\left(t_{i}+\tau\right)-N\left(t_{i}\right), \quad Y_{j}=N\left(s_{j}+\tau\right)-N\left(s_{j}\right)
$$

where the intervals $\left(t_{i}, t_{i}+\tau\right)$ and $\left(s_{j}, s_{j}+\tau\right)$ for $i=1,2, \ldots$, $j=1,2 \ldots$ are disjoint. Therefore the $X_{i}$ and $Y_{i}$ are independent and identically distributed as Poisson with parameter $\tau \lambda$. Thus, to estimate $\lambda$ within $d_{l}$ units use the method presented here with $d=\tau d_{1}$, with $x_{i}$ equal to the value of the random variable $x_{i}$, and $y_{i}$ equal to the value of the random variable $Y_{i}$. Therefore we obtain

$$
\operatorname{Pr}\left[\left|\tau \lambda-\bar{y}_{n}\right|<\tau d_{1}\right]>1-\varepsilon
$$

which implies that

$$
\operatorname{Pr}\left[\left|\lambda-\bar{y}_{n} / \tau\right|<d_{1}\right]>1-\varepsilon
$$

the desired result.

```
4.8 Birnbaum and Healy's Solution
```

Birnbaum and Healy [8] give a two-step procedure to estimate $\lambda$ with the unbiased estimator $\bar{y}_{n}$ having variance not exceeding a prescribed bound. By using Tchebycheff's inequainty their method gives

$$
P\left[\left|\hat{\lambda}_{n}-\lambda\right|<d\right]>1-\varepsilon
$$

if

$$
\begin{equation*}
n=\left(m \bar{x}_{m}+z\right) / m \varepsilon d^{2} \tag{4.34}
\end{equation*}
$$

where $\bar{x}_{m}$ is the mean in a preliminary sample of size m.
Table 4.1 gives some comparisons of the second sample size when determined according to Bimbaum and Healy's solution, the one point solution from Section 3, and the two point solution from Section 5 for various values of $1-\varepsilon, m, d$, and $\bar{x}_{m}$. In these examples sample sizes were reduced by varying amounts from 50 to $90 \%$.

TABLE 4.1
Comparisons of Second Sample Sizes
Column I gives $n$ for Birnbaum and Healy's solution (equation 4.34).
Column II gives $n$ for the one point solution (equation 4.10), Column III gives $n$ for the two point solution (equation 4.25).

| 1-ع | m | $\bar{x}_{\mathrm{m}}$ | d | I | II | $n^{\prime}$ | n" | III |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| . 90 | 5 | 100 | 2 | 205.5 | 78 | 72 | 73 | 73 |
| . 95 | " | " | " | 501 | 112 | 98 | 104 | 104 |
| . 99 | " | " | " | 2505 | 196 | 180 | 178 | 180 |
| . 95 | 1 | 100 | 2 | 505 | 130 | 121 | 112 | 121 |
| " | 2 | " | " | 502.5 | 120 | 109 | 108 | 109 |
| " | 5 | " | " | 501 | 112 | 98 | 104 | 104 |
| " | 10 | " | " | 500.5 | 108 | 94 | 102 | 102 |
| " | 25 | " | " | 500.2 | 104 | 91 | 100 | 100 |
| " | $\infty$ | " | " | 500 | 97 | 84 | 97 | 97 |
| . 95 | 5 | 20 | 2 | 101 | 26 | 23.5 | 22 | 23.5 |
| " | " | 50 | " | 251 | 58 | 53 | 53 | 53 |
| " | " | 100 | " | 501 | 130 | 98 | 104 | 104 |
| " | " | 250 | " | 1251 | 265 | 233 | 252 | 252 |
| . 95 | 5 | 100 | 1 | 2004 | 444 | 400 | 416 | 416 |
| " | " | " | 2 | 501 | 112 | 98 | 104 | 104 |
| " | " | " | 10 | 20.04 | 4.4 | 4.0 | 4.2 | 4.2 |
| . 95 | 1 | 10 | 1 | 220 | 91 | 95 | 53 | 95 |
| " | 5 | " | " | 204 | 60 | 55 | 47 | 55 |
| " | 10 | " | " | 202 | 52 | 47 | 46 | 47 |
| " | 25 | " | " | 200. 8 | 47 | 42 | 44 | 44 |
| " | $\infty$ | " | " | 200 | 40 | 34 | 40 | 40 |
| . 95 | 5 | 2 | . 1 | 4400 | 1800 | 1900 | 1230 | 1900 |
| " | 10 | " | " | 4200 | 1450 | 1440 | 1060 | 1440 |
| " | 25 | " | " | 4080 | 1160 | 1130 | 940 | 1130 |
| " | $\infty$ | " | " | 4000 | 800 | 680 | 800 | 800 |
| . 95 | 56 | 2.23 | . 2 | 1125 | 285 | 250 | 245 | 250 |

### 4.9 Second Sample Size Graphs

To determine the second sample size, first specify $m, d$, and 1- $\varepsilon$ in advance. Then take a preliminary sample of size $m$ and compute $m \bar{x}_{m}$, the sum of the $m$ observations. If $m \bar{x}_{m}$ is less than 500 use Figure 4.4 to find $c$ for a one point solution or to find $c_{1}$ and $c_{2}$ if a two point solution is desired. For a one point solution the lower curve labeled $c$ is to be used for $1-\varepsilon=.90$ or .95. Similarly the upper curve labeled $c$ is to be used for $1-\varepsilon=$.99. If a two point solution is desired the value of $1-\varepsilon$ is immaterial. For $\overline{m x}_{m}$ greater than 500 the same procedure applies to Figure 4.5. Next compute $c \bar{x}_{m}$ or $c_{1} \bar{x}_{m}$ and $c_{2} \bar{x}_{m}$. Consider the one point solution. If the ratio of $c \bar{x}_{m}$ to $d$ is less than $30: 1$ for $1-\varepsilon=.99$, less than $50: 1$ for $1-\varepsilon=.95$, or less than $60: 1$ for $1-\varepsilon=.90$ use Figure 4.6. For larger ratios it is necessary to use Figure 4.7. Plot a point of the form $\left(k d, k c \bar{x}_{m}\right)$, for some value of $k$, on the graph. The value of $k$ chosen is immaterial but the larger the value used the more accurate will be the result. With a straight edge connect the plotted point with the origin. At the point of intersection between this line and the appropriate curve, depending upon the size of $1-\varepsilon$, record the value on the horizontal axis. This value is nd and thus the second sample size, $n$, is found by dividing this by $d$. The method to compute $n^{\prime \prime}$, needed for the two point solution, is identical with the above except $c_{2}$ is used instead of $c$. To determine $n^{\prime}$, Figure 4.8 is used for ratios of $c_{1} \bar{x}_{m}$ to $d$ less than $33: 1$ for $1-\varepsilon=099$, less than $54: 1$ for $1-\varepsilon=.95$, and less than $63: 1$ for $1-\varepsilon=, 90$. Larger ratios require Figure 4.9, Again the procedure is the same as that for finding $n^{\prime \prime}$ except $c_{2}$ is
replaced by $c_{1}$. Finally, to determine the second sample size in the two point solution take the larger of $n '$ and $n "$.

As an example, consider an actual experiment in which flowers were exposed to low level irradiation and the number of discolored sectors per petal were counted for 56 petals.

TABLE 4.2

| Frequency Count | 0 1 2 3 4 5 <br> 7 11 16 11 8 2 | 7 | 1 |
| :--- | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |

In this experiment $m=56, \bar{x}_{m}=2.23$ and $s^{2}=2.22$. Assuming this follows the Poisson distribution the true value of $\lambda$ can be estimated within .2 of a unit with confidence coefficient $1-\varepsilon=.95$ as follows. We compute $m \bar{x}_{m}=125$. From Figure 4.4 we see that $c=1.287$ which implies that $c \bar{x}_{m}=2.87$. The ratio of $c \bar{x}_{m}$ to $d$ is $14.35: 1$ enabling us to use Figure 4.6. The largest value of $k$ possible to use is 1000 , which corresponds to the point $(200,2870)$. The intersection of the straight line joining this point with the origin and the $1-\varepsilon=095$ curve has a value of nd $=57$ on the horizontal axis. Thus the second sample size for the one point solution is $n=57 / d=285$. For the two point solution $c_{1}$ and $c_{2}$ are read from Figure 4.4 as $c_{1}=1.377$ and $c_{2}=1.122$ implying $c_{1} \bar{x}_{m}=3.07$ and $c_{2} \bar{x}_{m}=2.50$. Again setting $k=1000$ we plot the point $(200,3070)$ in Figure 4.8 , connect it to the origin, and the intersection with the $1-\varepsilon=095$ curve is at $n ' d=49$ on the horizontal axis. Thus $n^{\prime}=49 / d=245$. Similarly, plotting the point $(200,2500)$ in Figure 4.6 yields $n^{\prime \prime} d=50$ implying $n^{\prime \prime}=50 / d=250$. Therefore the two point second sample size is $n=250$.





Figure 4.7



To derive the graphs for the one point solution various values of $B$ and $a$ were tried for the particular case where $m=5, \bar{x}_{m}=100$, and $d=2$ such that $(1-\beta)(1-\alpha)=1-\varepsilon$. The optional values creating the smallest sample size for $1-\varepsilon=.99$ were $\beta=.0005$, $1-\alpha=.9905$. For $1-\varepsilon=.90$ and $1-\varepsilon=.95$ the optimum values for $\beta$ were so close that $\beta=.002$ worked for both with negligible loss in sample size but a gain in simplicity. The $1-\alpha$ values were therefore. 902 and . 952 respectively. These values were used in Figures 4.4, 4.5, 4.6, and 4.7. In the two point solution it was also necessary to optimize $\gamma$ and $\delta$ (see Section 5). Again the case where $m=5, \bar{x}_{m}=100$, and $\mathrm{d}=2$ was chosen to optimize. In particular, for $1-\varepsilon=95$ the optimum was at $1-\varepsilon=.9520, \beta=.1379, \gamma=.9865$, and $\delta=.999$. This allowed the use of the same graph as in the one point solution with $1-\varepsilon=.95$ to compute $n^{\prime \prime}$. Therefore to optimize in the $I-\varepsilon=90$ case I set $1-\alpha=.902$ and the other values were optimized at $\beta=.1003, \gamma=.9800$, and $\delta=.998$. Finally, to optimize for $1-\varepsilon=.99$ I set $1-\alpha=.9905$ and $(1-\delta) \beta=.0001$ (in the $1-\varepsilon=.90$ case (1- $\delta$ ) $\beta$ optimized at. 0002 and for $1-\varepsilon=.95$ at .00014). This led to the optimum values $\beta=.0961$, $\gamma=.9958$, and $\delta=.999$. Because the sample size is fairly insensitive to small changes in $\beta$, a common value of $\beta=.10$ was chosen. Similarly, the same value of $\delta=.999$ was used with negligible loss. Thus $(1-\delta) \beta=.0001$ for all three cases. This consolidation then forced new values of $\gamma$ such that $1-\varepsilon=(1-\alpha)(1-\beta+\beta \gamma \delta)$ as desired. The new values of $\gamma$ were $\gamma=.981, .981$, and . 996 for $1-\varepsilon=.90,95$ and . 99 respectively. These changes greatly reduced the number of graphs involved but did not materially change the sample size.

Figures 4.4 and 4.5 graph $c$ versus $z$ for given values of $H(c ; z)$ where $H(c ; z)$ is defined by (4.6). In particular, the curves from top to bottom represent $\mathrm{H}=.0001, .0005, .002$, and . 10 . Figures 4.6 and 4.7 graph $n c \bar{x}_{m}$ versus ad for various values of $t\left(n ; c \bar{x}_{m}, d\right)$ which is defined by ( 4.8 ). This type of graph is practical because $t\left(n ; c \bar{x}_{m}, d\right)=t\left(1 ; n c \bar{x}_{m}, n d\right)$. From left to right the cuxves correspond to $t=.902, .952$, and .9905. Figures 4.8 and 4.9 are similar to Figures 4.6 and 4.7 with the curves corresponding to $\tau=\gamma(1-\alpha)=.8849$, .9339, and . 9865 from left to right.

The computations in Figures 4.4 and 4.5 were made from existing tables for values of $z \leq 50$. Table 4.3 displays these. For larger values of $z$ the normal approximation to the Poisson was used. No loss of accuracy was evident in several examples which were checked. The values for Figures $4.6-4.9$ were obtained by use of the IBM 1620 computer, using a program which computed the actual values by summation of the appropriate Poisson distribution.

Table 4.3


## V. UNSOLVED PROBLEMS

The two procedures given in Chapters III and IV clearly are an improvement over existing methods. These solutions, however, raise additional problems to be investigated.

Both procedures supply a lower bound but the upper bound should also be given consideration. There is some information on an upper bound in the Poisson problem. For instance, in the one point solution

$$
\begin{aligned}
\operatorname{Pr}\left[\left|\lambda-\hat{\lambda}_{n}\right|<d\right] & <1 \cdot \operatorname{Pr}\left[n \geq n_{1}\right]+(1-\alpha) \operatorname{Pr}\left[n<n_{1}\right] \\
& =1 \cdot(1-\beta)+(1-\alpha) \beta=1-\alpha \beta
\end{aligned}
$$

For $1-\varepsilon=.90,95, .99$ this gives upper bounds of .999804, .999904, and .99999525 respectively. These are not very useful. Since they were obtained using the optimal values for $\alpha$ and $\beta$ there exist other values of $\alpha$ and $\beta$ with $(1-\alpha)(1-\beta)=l-\varepsilon$, that will yield larger sample sizes and at the same time decrease the upper bound; $1-\alpha=1-\beta=\sqrt{1-\varepsilon}$ for example. An upper bound for a large sample size would also be an upper bound for a small sampie, thus the upper bound could be lowered. Other techniques should be investigated to put better limits on the comfidence coefficient, since a confidence coefficient much larger than the level desired increases sample size and wastes resources.

Another problem left unanswered is whether or not nore information can be used. For instance, both procedures described here neglect
the first sample once it has been used to determine the size of the second sample. Could the preliminary sample be used somehow with the second sample in the estimator? This seems to be a complex problem. As a starting difficulty, what estimator should be used?

An important problem to consider is that of the size of the first sample. To determine this some knowledge is neec̃ed of the approximate size of the parameter. A three-step procedure may bequired, where the first step would be needed to find a first, rough approximation to the size of the parameter.

In Chapter IV, Section 7, it was shown how to astimate the parameter in a Poisson process with a two-step procedure. It was necessary to specify a constant $\tau$ to do this. Some investigation should be made to devise a scheme for picking the best value for $\tau$. Perhaps a three-step solution would be necessary. The first step would be used to find $\tau$.

The graphs in Section 9 of Chapter IV were derived by finding the combination of $\alpha, \beta, \gamma, \delta$ which minimized the second sample size for the particular values $m=5, d=2$, and $\bar{x}_{\mathrm{m}}=100$. This is fully described in Chapter IV, Section 9. If the second sample size was minimized for each individual set of values of $m, d$, and $\bar{x}_{m}$, reductions would be made in sample size. The feasibility of doing this should be investigated. For one indication of how good the particular optimization procedure used is, for the two point solution, examine the difference between $n^{\prime}$ and $n^{\prime \prime}$. Table 40.1 shows that $n^{\prime}=95$ and $n^{\prime \prime}=53$ for $1-\varepsilon=.95, \mathrm{~m}=1, \bar{x}_{\mathrm{m}}=10$, and $\mathrm{d}=1$, indicating the optimization was not good for these values. In fact $n$ is greater than the one point
solution sample size which is 91 . Similar results occur for $1-\varepsilon=.35$, $m=5, \bar{x}_{m}=2$, and $d=.2$. It is noted though that langer values of $m$ reduce this difference in both cases and also decrease the two point solution second sample size below that of the one point solution. An investigation could be made into procedures for solving the problems presented in Chapters III and IV by minimizing expected sample size. This would require a different approach from those in Chapters III and IV。

Also of interest would be a Bayes' type of solution in which various loss functions could be considered.

The solution for the variance of a normal problem in Chapter III can be generalized to estimate the mean in the gamna distribution with little modification. To solve other problems with this technique it would be necessary to first derive improvements on Tchebycheff's inequality for the distribution involved. The solution for the Poisson mean problem in Chapter $V$ seems to be very useful and the same technique could be used in many problems requiring a two-step procedure. In fact, it may result in lower sample sizes in estimating the variance of the normal. This problem and many others should be investigated using the technique in Chapter IV

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Two-step sampling procedures are presented to estimate the variance of a normal distribution and the mean of a Poisson distribution within d units with a specified confidence coefficient.

The procedure to estimate the variance of a normal is based on a Tchebycheff type inequality derived especially for the gama distribution. A different type of argument, which could be applied to many other distributions, was used to solve the problem for the Poisson distribution.

Sampling sizes are presented in tables and graphs to implement the two solutions. Also, favorable comparisons are made with existing methods.

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