# PI <br> MU EPSIL $\mathrm{O}_{\mathrm{N}}$ JOURNAL 

## VOLUME 9 SPRING 1992 NUMBER 6

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PI MU EPSILON JOURNAL
THE OFFICIAL PUBLICATIONOF THE
NATIONAL HONORARY MATHEMATICS SOCIETY

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Problems for solutionand solutionsto problemsshould be mailed directly to the PROBLEMEDITOR Puzzle proposalsand puzzle solutions should be mailed to the EDITOR

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## DEDICATION

This issue of the Pi Mu Epsilon Journal is dedicated to Joe Konhauser. Joseph D. E. Konhauser, Councillor for Pi Mu Epsilon, Editor of the Puzzle Section of the Pi Mu Epsilon Journal, and former Editor of the Journal, died on February 28, 1992, of complications following heart surgery. He was 67 years old. He is survived by his wife, Aileen, and his son, Dan.

Joe earned his bachelor's, master's, and doctorate degrees from Penn State University. From 1949 to 1955 he taught math at Penn State and was a mathematician at HRB-Singer Inc. in State College, PA. He was an associate professor of mathematics at the University of Minnesota for four years before joining the staff at Macalester College in St. Paul, MN, in 1968. He had retired from full time teaching at Macalester in 1991, but had returned to teach his popular geometry course this semester.

Besides his teaching and his wntributions to Pi Mu Epsilon, Joe had been a member of the committees that designed and evaluated tests for the USA Mathematical Olympiad and the William Lowell Putnam Mathematics Competition. He also had served as Fleviews Editor of the American Muthematical Monthly.

Joe had a real talent as a problem poser and solver. He had been Editor of the Puzzle Section since 1983; he had been creating Mathacrostics for the Puzzle Section since 1978. Perhaps even more remarkably, he had been posing a "Problem of the Week ${ }^{\mathrm{h}}$ at Macalester College for over 20 years, without repeating a problem.

Joe will be missed as a mathematician, as a renowned teacher, and as a friend.

## PUZZLE SECTION

The Puzzle Section is for the enjoyment of those readers who are addicted to working doublecrostics or who find an occasional mathematical puzzle or word puzzle attractive. We consider mathematical puzzles to be problems whose solutions consist of answers immediately recognizable as correct by simple observation and requiring little formal proof.

COMMENTS ON PUZZLES 1-7, FALL 1991
For Puzzle \#1, WILLIAM CHAU, RICHARD I. HESS, HENRY LIEBERMAN, and BOB PRIELIPP noted the following: Let $K$ be the area of the triangle, $\boldsymbol{s}$ be its semi-perimeter, and $\boldsymbol{r}$ its inradius. Then $\boldsymbol{r} s=K=2 s$, thus $r=2$. Other relationships satisfied by these triangles were provided by MARK EVANS and CHARLES ASHBACHER.

The answer to Puzzle \#2, the "Bronzebach Conjecture," is yes. Several decompositions were submitted. Perhaps the most concise was by BOB PRIELIPP:

$$
\begin{gathered}
\text { If } \mathrm{n} \text { is odd, } \mathrm{n}=(\mathrm{n}-2)+2 . \quad \text { If } n=4 k, n=\left(\frac{n}{2}-1\right)+\left(\frac{n}{2}+1\right) . \\
\\
\text { If } n=4 k+2, n=\left(\frac{n}{2}-2\right)+\left(\frac{n}{2}+2\right) .
\end{gathered}
$$

Solutions were submitted by CHARLES ASHBACHER, WILLIAM CHAU, VICTOR FESER, RICHARD I. HESS, HENRY LIEBERMAN, and DAVID SHOBE.

The first of the two solutions to Puzzle \#3 was submitted hy DAVID SHOBE; the second solution appeared in The Oxford Guide to Word Games by Tony Augarde, 1984, p. 44.


The solution to Puzzle \# 4 is no. Suppose there were a solution with the set $\{7,8,9\}$. Consider the set containing 15 . To complete the sum of 24 , we need either $1 \& 8$ (but 8 is gone), or 2 $\& 7$ (but 7 is gone), or $3 \& 6$, or $4 \& 5$. In either of these last two cases, there are no pairs of remaining numbers that will go with 14 to reach a sum of 24 . (Solution by VICTOR FESER.) Others submitting solutions were CHARLES ASHBACHER, WILLIAM CHAU, MARK EVANS RICHARD I. HESS, and HENRY LIEBERMAN.

There were several different solutions to Puzzle \# 5. The one that kept the three pieces the most similar in shape was submitted by MARK EVANS.


$$
\begin{aligned}
& \text { where } \mathrm{a}=\frac{1}{3} L-\frac{1}{6} W \\
& \qquad b=\frac{1}{3} W
\end{aligned}
$$

Others submitting solutions were WILLIAM CHAU, RICHARD I. HESS, DAVID SHOBE, and STAN WAGON.

For Puzzle \#6, several solutions were submitted. The most common was


Submitting solutions were WILLIAM CHAU, MARK EVANS, VICTOR FESER, RICHARD I. HESS, and DAVID SHOBE.

For Puzzle \#7, the resistance was found to be the Golden Ratio: $(\sqrt{5}-1) / \mathbf{2}$. Solvers were MARK EVANS, ROBERT GEBHARDT, HENRY LIEBERMAN, and DAVID SHOBE.

SOLUTION TO MATHACROSTIC NO. 33 (FALL 1991)
WORDS:

| A. | openness | K. | entify | U. | Florentine enigma |
| :--- | :--- | :--- | :--- | :--- | :--- |
| B. | Verdict of Twelve | L. | left-handed | V. | theremin |
| C. | extenuate | M. | Ymir | W. | hypergraphics |
| D. roses of grandi | N. | hem and haw | X. | even steven |  |
| E. | Baily's beads | O. | earth | Y. | corkscrew |
| F. yatata yatata | P. | Alan Smithee | Z. | odd |  |
| G. extent | Q. | right | a. | spread |  |
| H. lute | R. | The Great White Way | b. | Modern Times |  |
| I. on the charts | S. | switch | c. | Of Thee I Sing |  |
| J. nines | T. | overtone | d. | spherical cow |  |

AUTHOR AND TITLE: OVERBYE LONELY HEARTS OF THE COSMOS
QUOTATION: The veneer of existence was getting very, very thin, but it was in that last little crack of time - where space foamed into chaos and the spheres rang with harmonies undreamed of and symmetries were enfolded more intricately than a rose, where nothing happened and everything was possible - that the secret of gravity and existence lay.

SOLVERS: THOMAS F. BANCHOFF, Brown University, Providence, RI; JEANETTE BICKLEY St. Louis Community College at Meramec, MO; CHARLES R. DIMINNIE, St. Bonaventure University, NY; ROBERT C. GEBHARDT, Count College of Morris, Randolph, NJ; META HARRSEN, New Hope, PA; HENRY S. LIEBERMAN, Waban, MA; CHARLOTTE MAINES, Rochester, NY STEPHANIE SLOYAN, Georgian Court College, Lakewood, NJ; JOHN L. VANIWAARDEN, Hope College, Holland, MI, and DONNA D. ASHBRIDGE, University of North Carolina - Asheville, NC; ALBERT WILANSKY, Lehigh University, Bethlehem, PA; and BARB ZEEBERG, Denver, CO.

MATHACROSTIC NO. 34

## Proposed by Joseph D. E. Konhauser, shortly before his death.

The 223 letters to be entered in the numbered spaces in the grid will be identical to those in the 23 keyed words at the matching numbers. The key numbers have been entered in the diagram to assist in constructing the solution. When completed, the initial letters on the Words will give the name of an author and the title of a book; the completed grid will be a quotation from that book. Solutions to Mathacrostic No. 34 should be sent to: Richard Poss, Pi Mu Epsilon Journàd, St. Norbert College, De Pere, WI 54115.

## DEFNIIONS

HORDS
A．a migratory Australian cuckoo（2 wds．）
B popular name of Dilworth＇s 1740 ＂A Mea Guide to the English Tongue．＂（comp）

C．study of disease by symptoms

D．probability the first to give theoretical construction for all the five regular solids and to show how to

E．game plan
F．secret asset or ploy（4 wds．）

G．a musical means by which to identify characters，ideas，and objects as and at different times situation
．invented and patented by Kenneth Snel son，it has added a new component to the elegance and airiness
I．must starting point of the space frame for a cycle on a graph（2 wds．）

J．a fanciful product of the mind
K．a field of granular snow
L．used in prescriptions－of each an equal quantity

N．the nickname of the largest simple sporadic group（ 2 wds．）

N．to deprive of vital content or force

0．rural setting for Schubert opera （2 wds．）

P．inciden
Q．leaving no loophole
R．aurked vessel anchored at a chartered point to serve as an aid to navigation

S．famous or unfanous（3 wds．）

T．all out（3）

U．unpublished
$\overline{35} \quad \overline{48} \quad \overline{107} \overline{114} \quad \overline{161} \overline{192} \overline{67} \overline{200} \overline{23}$ $\overline{181} \overline{49} \overline{60} \overline{34} \overline{125} \overline{86} \overline{74} \overline{43} \overline{162}$ $\overline{179} \overline{122} \overline{147} \overline{139} \overline{44} \overline{169} \overline{50} \overline{218} \overline{66}$ $\overline{221} \overline{36} \overline{176} \overline{54} \overline{197} \overline{102} \overline{27} \overline{10} \overline{145}$ $\overline{208}$
$\overline{90} \overline{41} \overline{188} \overline{135} \overline{120} \overline{159} \overline{70} \overline{68}$
$\overline{{ }^{99}} \overline{155} \overline{28} \overline{115} \overline{213} \overline{55} \overline{77} \overline{129} \overline{175}$ $\overline{220} \overline{4} \overline{101} \overline{173} \overline{123} \overline{24} \overline{93} \overline{194}$ $\overline{32} \overline{151} \overline{140} \overline{\overline{178}}$
$\overline{166} \overline{210} \overline{171} \overline{127} \overline{100} \overline{25} \overline{31} \overline{111}$
$\overline{63} \overline{33} \overline{116} \overline{3} \overline{19} \overline{46} \overline{84} \overline{141} \overline{217}$ 185
$\overline{133} \overline{53} \overline{167} \overline{203} \overline{187} \overline{94} \overline{13} \overline{149} \overline{215}$ $\overline{73}$
$\overline{158} \overline{85} \overline{198} \overline{65} \overline{134} \overline{112} \overline{51}$
$\overline{126} \overline{9} \overline{214} \overline{137}$
$\overline{17} \overline{80} \overline{99}$
$\overline{121} \overline{21} \overline{96} \overline{38} \overline{144} \overline{33} \overline{180} \overline{204} \overline{6}$ ${ }^{223}$
$\overline{39} \overline{132} \overline{212} \overline{154} \overline{177} \overline{72} \overline{52} \overline{164} \overline{195}$ $\overline{12}$
$\overline{182} \overline{57} \overline{77} \overline{69} \overline{124} \overline{104} \overline{113} \overline{211} \overline{216}$

$\overline{149} \overline{160} \overline{15} \overline{98} \overline{62} \overline{109} \overline{75}$
$\overline{191} \overline{56} \overline{170} \overline{20} \overline{5} \overline{68} \overline{64} \overline{30}$
$\overline{190} \overline{82} \overline{142} \overline{103} \overline{106} \overline{16} \overline{128} \overline{153} \overline{61}$
$\overline{{ }^{143}} \overline{2} \overline{42} \overline{136} \overline{14} \overline{184} \overline{97} \overline{78} \overline{205}$ $\overline{193} \overline{110} \overline{219} \overline{150} \overline{163}$
$\overline{209} \overline{22} \overline{156} \overline{92} \overline{59} \overline{29} \overline{705}$
$\begin{array}{lll}168 & \overline{186} & \overline{87} \\ \overline{117}\end{array}$
$\overline{118} \overline{91} \overline{199} \overline{71} \overline{193} \overline{157} \overline{207} \overline{152}$

V．biotechnology

Y．originally developed around 1966－70 in an attenpt to understand the strong nuclear attent to understand the strong nucle －leualfication of physics（ 2 ms ．）
$\overline{79} \overline{138} \overline{119} \overline{196} \quad \overline{58} \quad \overline{40} \quad \overline{18} \quad \overline{202} \overline{8}$ $\overline{172}$
$\begin{array}{lllllllll}146 & \frac{81}{2} & 206 & 189 & 45 & 26 & 174 & \overline{11} & \overline{131}\end{array}$ $\overline{222} \overline{95} \overline{201}$

| $1{ }^{1} 6$ | 2 |  | 3 H | 4 | 5 Q | 6 M | 7 | $8^{8} \quad \mathrm{~V}$ | 9 K |  | 10 D | 11 V | 12 N |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 13 I | 14 S |  | 15 P | 16 R |  | 17 L |  | 18 V | 19 H | 20. | 21 H | 22 T | 23 A |
|  | 24 F | 256 |  | 26 甘 | 27 D |  |  | 29 T | 30 a |  | 316 | 32 F | 33 H |
|  | 34 B | 35 A | 36 D | 370 | 38 N | 39 N |  | 40 V | 41 D |  | 42 S | 43 B | 44 C |
| 45 प | 46 H | 470 |  | 48 A | 49 B | 50 C | 51 J | 52 N | 53 I | 54 D | 55 E | 56 Q | 570 |
| 58 |  | 59. | 60 B | 61 R | 62 P | 63 H | 64 a | 65 J | 66 c | 67 A | 68 D |  | 690 |
| 70 D | 71 U |  | 72 N | 73 I | 74 B | 75 P | 760 | 77 E | 78 S | 79 V | 80 | 81 甘 |  |
| 82 R | 83 H |  | 84 H | 85 J | 86 B | 87 T | 88 Q | 89 E | 90 D | 91 U |  | 92 T | 93 F |
| 94 I | 95 \％ | 96 M |  | 97 S | 98 P |  | 99 L |  | 1006 | 101 F | 102 D | 103 R | 1040 |
| 105 T |  | 106 R | 107 A | 1080 |  | 109 P | 1105 | 1116 | 112 J | 1130 | 114 A |  | 115 E |
| 116 H | 117 T | 118 U | 119 V | 120 D | 121 M | 122 C | 123 F | 1240 | 125 B | 126 K | 1276 |  | 128 R |
| 129 E | 1300 | 131 H | 132 N | 133 I | 134 J |  | 135 D | 136 S | 137 K |  | 138 V | 139 C | 140 F |
| 141 H | 142 R | 143 S | 144 A | 145 D | 146 H |  | 147 C | 148 P | 149 I | 150 S | 151 F | 152 U |  |
| 153 R | 154 N |  | 155 E | 156 T | 157 U |  | 158 J | 159 D | 160 P | 161 A | 162 B | 163 S | 164 N |
| 1650 | 166 G | 167 I |  | 168 T | 169 C | 170 Q |  | 1716 | 172 V |  | 173 F | 174 반 |  |
| 175 E | 176 D | 177 N | 178 F | 179 C | 180 M | 181 B | 1820 | 183 U | 184 S | 185 H |  | 186 T | 187 I |
| 188 | 189 H | 190 R | 1910 | 192 A | 193 S | 194 F |  | 195 N | 196 V |  | 197 D | 198 J | 199 U |
| 200 A | 201 V |  | 202 V | 203 I | 204 M | 205 S | 206 y | 207 U | 208 D | 209 T | 2106 | 2110 |  |
| 212 N | 213 E | 214 K | 215 I | 2160 | 217 H | 218 C | 219 S | 220 F | 2210 | $222 甘$ | 223 H |  |  |

## PUZZLES FOR SOLUTION

To give some idea of the types $\boldsymbol{f}$ problem that Joe Konhauser liked to devise, this issue's Puzzle Section will present some of Joe's puzzles that had previously appeared in the Journal. The solution to each puzzle was discussed in the issue immediately following the puzzle's appearance.

## 1. This problem first appeared in the Spring, 1989, issue of the Journal.

In the square array

$$
\begin{array}{lll}
A & B & C \\
C & B & D \\
E & C & F
\end{array}
$$

each letter represents one of the digits 0 through 9. Determine the correspondence, given that:
(1) $\boldsymbol{A B C}$ and $\boldsymbol{C B D}$ are primes,
(2) $\boldsymbol{B} \boldsymbol{B C}$ and $\boldsymbol{C D F}$ are perfect squares, and
(3) $\boldsymbol{A C E}$ and $\boldsymbol{E C F}$ are perfect cubes.

## 2. (Fall, 1983)

Sketch a graph (a finite collection of nodes and arcs) such that exactly three arcs terminate a each node and such that it is not possible to color the arcs with three colors so that no two arcs that are the same color terrninate at the same node.
3. (Fall, 1983)

The eight numbers $\{2,3,4,6,9,14,22, \mathbf{3 1}\}$ have sum $\mathbf{9 1}$ and the property that taken two at a time the 28 sums obtained are all different. Are you able to find 8 positive integers with sum less than 91 with the same property?

## 4. (Spring, 1984)

Using just two colors, in how many distinguishable ways can one color the edges of a regular tetrahedron?

## 5. (Fall, 1984)

The trio of positive integers $\{5, \mathbf{2 0}, \mathbf{4 4}\}$ has the property that the sum of any two of its members is a perfect square. Can you find a set of four distinct positive integers such that the sum of any three is a perfect square?

## 6. (Spring, 1985)

With a pair of compasses draw a circle on a plane. Then, without changing the opening of the compasses, draw a circle on a sufficiently large sphere. Which circle encloses the larger area?

## 7. (Fall, 1987)

Bored in a calculus class, a student started to play with a hand-held calculator. A four-digit number was entered, followed by the "square" key. To the surprise (and delight) of the student, the four terminal digits of the result were the same digits in the same order as those in the number which had been squared. What was that number?

## THE RICHARD V. ANDREE AWARDS

Richard V. Andree, Professor Emeritus of the University of Oklahoma, died on May 8, 1987, at the age of 67. Professor Andree was a Past-President of Pi Mu Epsilon. He had also served the society as Secretary-General and as Editor of the Pi Mu Epsilon Journal. The Society Council has designated the prizes in the National Student Paper Competition as Richard V. Andree Awards

First prize winner for 1991 is Amy Pinegar, for her paper "Inversions and Adjacent Transpositions," which appeared in the fall issue of the Journal. Amy prepared this paper, under the supervision of Dr. David Sutherland, while she was a senior at Middle Tennessee State University. She also presented the paper at the August, 1990, national Pi Mu Epsilon meeting at Columbus, Ohio. Amy will receive $\$ 200$.

Second prize winner is Shannon Spittler, for her paper "A Math Problem Within an Antique Clock Label," which also appeared in the fall issue of the Journal. Shannon prepared this paper while she was a junior English major at Miami University in Ohio. She will receive $\$ 100$.

Third prize winner is Judy Kenney, for her paper "Turning Triangles into Circles," which also appeared in the fall issue of the Journal. Judy prepared this paper while she was a senior at the College of St. Benedict. The problem was suggested to her by Dr. Steven Krantz while she was participating in an NSF Summer Research program at Washington University in St. Louis. Judy will receive $\$ 50$.

There were three other student-written papers that appeared in 1991:
"Computerized Segmentation of Liver Structures from CT Images," by Heng Hak Ly, of Illinois Benedictine College. Heng prepared this paper with the help of Dr. Maryellen Giger and Dr. Rose Carney.
"A Note on a Paper $\boldsymbol{f}$ S. H. Friedberg," by Janet Valasek, of Penn State University - New Kensington Campus. Janet prepared this paper with the assistance of Dr. Javier Gomez-Calderon.
"A Pre-Calculus Method for Deriving Simpson's Rule," by John White, of Marshall University.
The current issue of the Journal contains two papers written by students:
"Change Ringing: Mathematical Music" was written by Heather DeSimone while she was a senior at Youngstown State University. She is currently attending graduate school at the College of William and Mary.
"Rings of Small Order" was written by Michael Lin while he was a senior at Moorhead High School, in Moorhead, MN. He is currently a freshman at Stanford University.

Joel Atkins, the winner of third prize in the 1990 Competition, wishes to acknowledge the guidance of Dr. Jack Kinney of Rose-Hulman Institute of Technology.

## RINGS OF SMALL ORDER

Michael H. Lin<br>Stanford University

## Introduction

Since all finite abelian groups have a simple structure, a straightforward way to find all finite rings is to begin with its additive group. If we are given one particular additive group, say $\boldsymbol{G}, \boldsymbol{+}$, to work with, the problem is reduced to finding all binary operations "." on $G$ that are associative and that are left and right diitributive over " + ". This will largely be a matter of trial and error, and thus in general will be computation-intensive. One naive approach would be to try all $\boldsymbol{n}^{\boldsymbol{n}^{2}}$ possible multiplication tables and to check associativity and distributivity for each one of them.

In this paper, a more efficient approach is developed, and a computer program implementing it was written for use on an IBM PC compatible. This program takes as input a standard decomposition of the additive group, and outputs the multiplication tables of all possible rings with that additive group. The program does not determine which outputs are isomorphic. It works for any additive group with order up to $\mathbf{1 2 7}$, although in many cases a complete run would be impractical because of both the amount of generated output and the length of run time.

## Notation

The order of any group $\boldsymbol{G}$ will be denoted by $|G|$; likewise, the order of any element $g$ will be denoted by $|g|$. The cyclic group of order $\mathbf{n}$ will be denoted by $\boldsymbol{C}_{\boldsymbol{n}}$.

Let our given additive group $\boldsymbol{G}$ of order $\mathbf{n}$ be expressed as a direct product of nontrivial cyclic groups $H_{1} \times H_{2} \times \ldots x H_{r}$, where $\left|H_{k}\right|$ divides $\left|H_{k-1}\right|$ for $1<\mathrm{k} \leq r$. (Such a representation is uniquely determined by $G$.)
For $1 \leq \mathrm{k} \leq \boldsymbol{r}$, pick $\boldsymbol{h}_{\boldsymbol{k}}$ in $\boldsymbol{G}$ such that

$$
(h k)=\{0\} \times \ldots \times\{0\} \times H_{k} \times\{0\} \times \ldots \times\{0\}
$$

Let $B=\left\{h_{1}, \boldsymbol{h}_{2}, \ldots, \boldsymbol{h}_{r}\right\}$, so that $(B)=G$.

## Algorithm

1. Input the orders of the $H_{\boldsymbol{k}}$.
2. Compute the addition table and other information about $G$ (such as the multiples and order of each element).
3. Set up a loop so that each passage through the loop assigns a value in $G$ to each of the $r^{2}$ products obtained from $B$. Successive passages through the loop assign every possible value in $G$ to each of the $\boldsymbol{r}^{2}$ products.
4. Check the necessary condition that $\left|h_{j}\right|\left(h_{j} \cdot \boldsymbol{h}_{k}\right)=\left|h_{k}\right|\left(h_{j} \cdot \boldsymbol{h}_{k}\right)=0 \quad \forall h_{j}, \boldsymbol{h}_{k} \in B$.
5. Define the remaining products within $G$ by distributivity: The distributive properties

$$
a(b+c)=a b+a c \quad \text { and } \quad(a+b) c=a c+b d
$$

can be restated as

$$
\left(\sum x_{j}\right)\left(\sum y_{k}\right)=\sum \sum x_{j} y_{k} .
$$

So for any $\boldsymbol{x}=\sum a_{j} h_{j}, y=\sum b_{k} h_{k}$, distributivity gives

$$
\begin{aligned}
x \cdot y & =\left(\sum a_{j} h_{j}\right) \cdot\left(\sum b_{k} h_{k}\right) \\
& =\sum \sum\left(a_{j} h_{j} \cdot b_{k} h_{k}\right) \\
& =\sum \sum a_{j} b_{k}\left(h_{j} \cdot h_{k}\right) .
\end{aligned}
$$

Using the condition of Step 4, it can be shown that the operation "." as given here is welldefined.
6. Check for associativity within $B$; i.e., that

$$
\begin{equation*}
\left(h_{i} \cdot h_{j}\right) \cdot h_{k}=h_{i} \cdot\left(h_{j} \cdot h_{k}\right) \quad \forall h_{i}, h_{j}, h_{k} \in B \tag{*}
\end{equation*}
$$

This is sufficient because, if (*) is satisfied,

$$
\begin{aligned}
\left(\sum a_{i} h_{i} \cdot \sum b_{j} h_{j}\right) \cdot \sum c_{k} h_{k} & =\left[\sum \sum a_{i} b_{j}\left(h_{i} \cdot h_{j}\right)\right] \cdot \sum c_{k} h_{k} \\
& =\sum \sum \sum a_{i} b_{j} c_{k}\left[\left(h_{i} \cdot h_{j}\right) \cdot h_{k}\right] \\
& =\sum \sum \sum a_{i} b_{j} c_{k}\left[h_{i} \cdot\left(h_{j} \cdot h_{k}\right)\right] \\
& =\sum a_{i} h_{i} \cdot\left(\sum b_{j} h_{j} \cdot \sum c_{k} h_{k}\right) .
\end{aligned}
$$

7. If Steps 4 and 6 are both satisfied, we have generated a ring. Output it.
8. Repeat Steps 4 through 7 as indicated in Step 3.

Results
The computer program was written in IBM PC assembly language. In addition to the multiplication tables for each ring and the total number of rings generated, the program also outputs the first $n^{-1}$ powers of each element for each ring. This shows certain properties of the ring at a glance - for example, how many squares are non-zero, and whether all cubes are zero - and thus makes it easier to see which rings might be isomorphic. It also tells at a glance that some pairs of rings are not isomorphic.

The generated rings for some additive groups were hand-classified according to isomorphism. The obtained results agreed with the list published in [1] of all 24 rings, up to isomorphism, of order less that 8 .

The data in the following table were obtained using an 8 Mhz IBM PC/XT clone running a stripped-down version of the computer program. The deleted parts of the program were those that computed the multiples and powers of each element. It should be noted that this trimmed version is significantly faster than the original program.

For cyclic additive groups of order $n$, the program produced a total of $\mathbf{n}$ rings. It was proved in [2] that the number of non-isomorphic rings with additive group $\boldsymbol{C}_{\boldsymbol{n}}$ is the number of divisors of n ; this was verified for n up to 10 .

| Structure of Additive Group | Total Number of Rings Produced | Computation Time |
| :---: | :---: | :---: |
| $C_{2} \times C_{2}$ | 28 (8 non-isomorphic) | $\sim 0.3 \mathrm{sec}$. |
| $C_{4} \times C_{2}$ | 60 (20 non-isomorphic) | 1 sec . |
| $C_{6} \times C_{2}$ | 84 (16 non-isomorphic) | 3 sec . |
| $\mathrm{C}_{8} \times \mathrm{C}_{2}$ | 120 | 7 sec . |
| do $\times C_{2}$ | 140 | 14 sec . |
| $C_{12} \times C_{2}$ | 180 | 25 sec . |
| $C_{14} \times C_{2}$ | 196 | 39 sec . |
| $C_{3} \times C_{3}$ | 121 (8 non-isomorphic) | 17 sec . |
| $C_{6} \times C_{3}$ | 242 | 116 sec . |
| $C_{9} \times C_{3}$ | 405 | 6.2 min . |
| $C_{12} \times C_{3}$ | 484 | 15.0 min . |
| $C_{15} \times C_{3}$ | 605 | 28.9 min. |
| $C_{4} \times C_{4}$ | 616 | 7.8 min. |
| $C_{8} \times C_{4}$ | 1376 | 57 min . |
| $C_{12} \times C_{4}$ | 1848 | 196 min . |
| $C_{16} \times C_{4}$ | 2816 | 7.6 hr . |
| $C_{20} \times C_{4}$ | 3080 | 14.8 hr . |
| $C_{5} \times C_{5}$ | 793 | 106 min . |
| $C_{10} \times C_{5}$ | 1586 | 14.1 hr . |

The above data suggest that if $\boldsymbol{s}$ and t are relatively prime, the number of rings produced for additive group $\boldsymbol{C}_{\boldsymbol{s} \boldsymbol{t}} \times \boldsymbol{C}_{\mathbf{s}}$ is $\boldsymbol{t}$ times the number for $\boldsymbol{C}_{\mathbf{s}} \boldsymbol{x} \boldsymbol{C}_{\mathbf{s}}$.

For large $n$ and small $\Gamma$ (as defined in Notation), Step 5, i.e., the completion of the multiplication table, dominates the other steps in terms of the computation time needed. Also, the time required for one execution of Step 5 is approximately proportional to the size of the multiplication table. Thus, for large $n$ and small $r$, a rough indicator of the total computation time would be

$$
\text { (the number of potential rings that pass Step 4) } \mathbf{x}\left(\boldsymbol{n}^{2}\right) \text {. }
$$

For an additive group of the form $C_{s t} \times C_{s}$, this expression simplifies to $s^{12} t^{3}$.
Additional Observations
While the rings with additive group $\boldsymbol{C}_{\mathbf{4}} \times \boldsymbol{C}_{\mathbf{2}}$ were being hand-classified up to isomorphism, it was noticed that four rings were anti-automorphic; i.e., that there existed a bijection $f$ on each of the rings such that $f(x \cdot y)=f(y) \cdot f(x)$ for all $x$ and y in the ring. The following theorems were then formulated.

## Preliminaries

Let $\boldsymbol{A},,+$ be the abelian group $\boldsymbol{C}_{s^{2}} \times \boldsymbol{C}_{s}=\left(\right.$ a) $\times(6)$. We define $\mathrm{f}: A, \longrightarrow \boldsymbol{A}_{\mathbf{s}}$ by $f(\boldsymbol{p a}+q b)=$ $p(a-b)+q(-b)$ IO all $p, q \in Z$. It can be shown that f is its own inverse; so f is a bijection. Also, it follows immediately from the definition that f preserves the operation " + ",

Theorem 1. Let $R_{s}$ be the ring with additive group $A$, and with multiplication defined by the relations

$$
\begin{aligned}
& a \cdot a=a \cdot b=0 \\
& b \cdot a=b \cdot b=s a .
\end{aligned}
$$

Then $R_{s}$ is a non-commutative ring that is anti-isomorphic to itself. Proof. It can be verified that

$$
|x|(x \cdot y)=|y|(x \cdot y)=0 \quad \forall x, y \in\{a, b\}
$$

and that

$$
(x \cdot y) \cdot z=x \cdot(y \cdot z) \quad \forall x, y, z \in\{a, b\}
$$

It follows from Steps 5 and 6 that $R_{s}$ is a ring. Also, $R_{s}$ is obviously non-commutative.
By the preliminaries, f is a bijection and $\mathrm{f}(x+\mathrm{y})=\mathrm{f}(\mathrm{y})+\mathrm{f}(\boldsymbol{x}) \quad \forall x, y \in \boldsymbol{R}_{\mathbf{s}}$.
We shall show that $\mathrm{f}(\boldsymbol{x} \cdot \mathrm{y})=\mathrm{f}(\mathrm{y}) \cdot \mathrm{f}(\mathrm{x}) \quad V \mathrm{x}, \mathrm{y} \in R_{\mathrm{s}}$.

$$
\begin{aligned}
& f(a \cdot a)=f(0)=0 ; \\
& f(a) \cdot f(a)=(a-b) \cdot(a-b)=a \cdot a-a \cdot b-b \cdot a+b \cdot b=0-0-s a+s a=0 . \\
& f(a \cdot b)=f(0)=0 ; \\
& f(b) \cdot f(a)=-b \cdot(a-b)=-b \cdot a+b \cdot b=-s a+s a=0 . \\
& f(b \cdot a)=f(s a)=s f(a)=s(a-b)=s a-s b=s a-0=s a \\
& f(a) \cdot f(b)=(a-b) \cdot(-b)=-a \cdot b+b \cdot b=0+s a=s a . \\
& f(b \cdot b)=\mathrm{f}(\mathrm{sa})=\ldots=\mathrm{sa} ; \\
& f(b) \cdot f(b)=(-6) \cdot(-6)=6.6=\mathrm{sa} .
\end{aligned}
$$

The general fact that $\mathrm{f}(\mathrm{x} \cdot \mathrm{y})=f(y) \cdot \mathrm{f}(\boldsymbol{x})$ now follows by distributivity.
Therefore, $f$ is an anti-automorphism of $\boldsymbol{R}$,.
Theorem 2. Let $Q$, be the ring with additive group $A$. and with multiplication defined by the relations

$$
\begin{aligned}
& a \cdot b=0 \\
& a \cdot a=b \cdot a=b \cdot b=s a .
\end{aligned}
$$

Then $Q_{3}$ is a non-commutative ring that is anti-isomorphic to itself.
The proof is exactly as for Theorem 1, except
$f(a \cdot a)=f(s a)=\ldots=s a ;$
$f(a) \cdot f(a)=(a-b) \cdot(a-b)=a \cdot a-a \cdot b-b \cdot a+b \cdot b=s a-0-s a+s a=s a$.

These two theorems raise some interesting questions about anti-automorphic non-commutative rings: What conditions upon an additive group are necessary and sufficient for there to exist antiautomorphic non-commutative rings with this additive group? How many anti-automorphic noncommutative rings exist, up to isomorphism, for any given additive group? Can a general description of their multiplication tables be given? What can be said about their structure? What other properties do they have?

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A computer program listing and a sample of the program output can be obtained by writing the author at P.O. Box 4048 Stanford, CA 94909 (e-mail: michelin@leland.stanford.edu).

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The Pi Mu Epsilon Journal was founded in 1949 and is dedicated to undergraduate and beginning graduate students interested in mathematics. Submitted articles, announcements, and contributions to the Puzzle Section and Problem Department of the Journal should be directed toward this group.

Undergraduates and beginning graduate students are urged to submit papers to the Journal for consideration and possible publication. Student papers are given top priority. Expository articles by professionals in all areas of mathematics are especially welcome. Some guidelines are:

1. Papers must be correct and honest.
2. Most readers of the Pi Mu Epsilon Journal are undergraduates; papers should be directed to them.
3. With rare exceptions, papers should be of general interest.
4. Assumed definitions, concepts, theorems, and notations should be part of the average undergraduate curriculum.
5. Papers should not exceed $\mathbf{1 0}$ pages in length.
6. Figures provided by the author should be camera-ready.
7. Papers should be submitted in duplicate to the Editor.

CHANGE RINGING: MATHEMATICAL MUSIC
Heather DeSimone
Youngstown State University

## 1. Introduction

Before the eighth century most church bells were small and rung by hand. These bells were made of iron and did not have good tone quality. Making bells from different alloys began around the eighth century. Using bronze, it became possible to create much bigger and louder bells. It was also discovered that different tones could be made by varying the thickness of the bell wall and the composition of the bell metal. The size of the bell also affected its sound. For example, the bigger bells made deeper sounding notes. At this time, large bells were being installed in church towers all over Europe. At the turn of the thirteenth century a gradual change in the shape of bells took place. The sides became longer and more concave, which improved tone.

As bells became larger and heavier, they became more difficult to ring. Consequently, methods of ringing evolved that did not require shifting the full weight of the bell. One of the methods, which is still employed today, is to swing a bell by a rope attached to the top until it is almost upside-down and then swing it back to complete the other half of the swing. This method was refined by mounting the bell on a half-wheel. A rope was then run around it and down to the floor, which provided a "stay" on the wheel's rim thereby preventing the bell from swinging all the way over. A final improvement was implemented soon after the Reformation when a whole wheel was introduced, allowing complete control over the bell. This improvement not only enabled the bell to stay in an upright position for as long as was needed, but more importantly, it allowed control over the speed of the swing. Pulling harder on the rope as it lowered sped up the swing. Conversely, retarding the rope as the bell swung up slowed down the swing. It was found that if two ringers of two different bells carried out these moves, the bells would change place in their order of ringing 'It was this discovery, when applied to a number of bells, that made 'change-ringing' possible; and this is the foundation on which the whole art of bellringing is based" (Camp, 15).
2. Change Ringing

The basic strategy of change ringing is:
(1) to ring a given set of bells in all possible sequences;
(2) to move in a methodical fashion from one sequence to another; and
(3) to avoid repeating any sequence.

There are n ! possible sequences for n bells. Each number of bells has a specific name as listed in the table.

| Number of Bells | Name | Number of Changes |
| :---: | :---: | :--- |
| 4 | Singles | 24 |
| 5 | Doubles | $\mathbf{1 2 0}$ |
| 6 | Minor | $\mathbf{7 2 0}$ |
| 7 | Triples | 5,040 |
| 8 | Major | 40,320 |
| 9 | Caters | 362,880 |
| 10 | Royal | $3,628,800$ |
| 11 | Cinques | $39,916,800$ |
| 12 | Maximus | $479,001,600$ |

The object of change ringing is to produce all of the permutations on a set of bells according to a set of rules. The highest bell is called the treble bell and the lowest, the tenor. When they ring
in descending order, from treble to tenor, they are said to be in rounds. The rules the bellringers must satisfy are:
(i) the peal must begin and end in rounds;
(ii) no bell may move more than one position from one change to the next; and
(iii) no bell may occupy the same position for more than two successive changes.

The last rule is sometimes relaxed
The six changes on three bells can be rung as follows:
213
231
231
321
312
132
123
or in the reverse order, but only these two ways follow the rules. These bells follow a hunting course. This means each bell works by steps of one to the right or left until the bell is first or last in the change. The first bell moves from the first position, to the second position, and to the third position. The bell then stays in the third position for two consecutive changes before it moves back to the second position. It then moves to the first position and stays there for the last two changes.

Look at the first transition. It can be denoted by the transposition (12) meaning that the bells in position 1 and 2 change places. The two operators applied in the changes on the three bells alternately are $A=(12)$ and $B=(23)$. These generate the entire group of order six. Algebraically, the six changes on three bells can be represented as $(A B)^{3}$ because $A$ and then $B$ are applied three times.

With four bells this is a little more complicated

|  | 1234 | 1342 | 1423 | 1234 |
| :---: | :---: | :---: | :---: | :---: |
| Plain Bob | 2143 | 3124 | 4132 |  |
| Method | 2413 | 3214 | 4312 |  |
|  | 4231 | 2341 | 3421 |  |
|  | 4321 | 2431 | 3241 |  |
|  | 3412 | 4213 | 2314 |  |
|  | 3142 | 4123 | 2134 |  |
|  | 1324 | 1432 | 1243 |  |

These bells also follow a hunting course. In the beginning, bell 1 is moved one position to the right. It then stays in the last position for two changes before moving backwards to its original position. The other three bells follow a similar hunting pattern. As four bells hunt, they create eight changes. In general, if $n$ bells hunt, the hunting generates a group of order $2 n$. The process of hunting on four bells consists of alternately applying the two operators $A=(12)(34)$ and $B=(23)$. As stated before, these generate the first eight changes. Continuing to use these operators, specifically using $\boldsymbol{B}$, would make the next sequence 1234. This is commonly known to bellringers as "replacing rounds!" It is not desirable because all of the possible changes would not have been rung. In order to prevent this, and to continue, we employ the irregular move $\boldsymbol{C}=(34)$. The second eight elements are generated by again applying $\boldsymbol{A}$ and $\boldsymbol{B}$. After the irregular move $\boldsymbol{C}$ is applied again, the third eight elements are generated the same as the first two sets. The Plain Bob method can be algebraically symbolized with the operators as:

## $\left((A B)^{3}(A C)\right)^{3}$

This notation means operators $A$ and $B$ are alternately applied three times
Then $\boldsymbol{A}$ and $\boldsymbol{C}$ are applied. This complete pattern is repeated three times.

## 3. Hamiltonian Circuit

Any particular set of complete changes can be graphically represented as a Hamiltonian cycle Let the nodes of the graph symbolize each change, that is, an ordering in which the bells are rung and the edges connect the possible consecutive changes. Graphing the changes on three bells, the Hamilton circuit is easily found. It is also easy to see that this is the only one since there are no edge left out of the circuit. This shows that there are only two ways of ringing the changes dependingon which direction the identity node is exited.


In general, the number of nodes is equal to $n\}$, where $n$ is the number of bells. The number of edges going to or coming from one node depends on the number of possible changes. In the example of three bells, we can interchange bells 1 and 2,(12) or bells 2 and 3, (23). These are the only two possibilities; therefore, there are two edges per node.

By increasing the number of bells by one, the number of nodes increases, as does the number of permutations. There are four possible ways to change from one sequence to another. The first three, which were discussed earlier, are (12)(34), (23), and (34). The last is switching only the first two, (12). So every one of the 24 nodes has four edges or is connected to four different nodes. This graph is more complicated than the one for three bells. A Hamiltonian circuit is not easy to find in the maze of 48 edges and 24 vertices; however, several can be found. The set of sequences discussed earlier is one example. The figure uses the form given in White [9].


This method is the most commonly rung and the most commonly displayed mathematically; however, there are others.
Names Algebraic Description

## Plain Bob

Reverse Bob
Double Bob
Canterbury
Reverse Canterbury
Double Canterbury
Single Court
Reverse Court
Double Court
St. Nicholas
Reverse St. Nicholas
$\left((A B)^{3} A C\right)^{3}$
$\left(A B A D(A B)^{2}\right)^{3}$
$(A B A D A B A C)^{3}$
$(A B C D C B A B)^{3}$
$\left(D B(A B)^{2} D C\right)^{3}$
$(D B C D C B D C)^{3}$
$(D B C D C B D C)^{3}$
$\left(D B(A B)^{2} D B\right)^{3}$
$\left(D B(A B)^{2} D B\right)^{3}$
$\left(A B(C B)^{2} A B\right)^{3}$
$\left(D B(C B)^{2} D B\right)^{3}$
$(D B A D A B D C)^{3}$
$(A B C D C B A C)^{3}$

It is worth noting that only the first three of these methods satisfy all three rules listed in section 2. The remaining methods fail condition (iii) that states no bell may stay in the same position for more than two consecutive changes.

Using the given operators, two original sequences will be demonstrated. Alternating $D$ and $\boldsymbol{B}$ with every sixth change using the $A$ transition completes the necessary Hamiltonian cycle. Algebraically, this is represented as $\left((D B)^{2}(D A)\right)^{4}$. And using the similar pattern $\left((C B)^{2}(C A)\right)^{4}$ also produces the circuit. The patterns are alike in that the second set replaces the $D$ 's of the first set with $C^{\prime \prime}$ s.

All methods with four bells use exactly 24 of the 48 edges to complete the Hamiltonian circuit. So it seems possible to find two independent cycles on the same graph. None of the above examples are independent of each other. In other words, two completely different Hamiltonian circuits can not be found with the previous instances. Therefore, starting with 48 edges and completing a Hamiltonian circuit in 24 edges does not necessarily mean there are two totally independent Hamiltonian circuits on that graph, even though there are 24 unused edges.

We will now show that such examples exist. By breaking down the earlier diagram into a simpler form where only the $\boldsymbol{B}$ connectors are left and each group of other lines are thought of as separate entities, we have:


From this diagram two independent Hamiltonian circuits can easily be determined. Studying further, we find that there are exactly six different pairs. Also, notice that this graph and the other five are not symmetric. All of the previous examples are symmetric. Now that we know what path to follow going in and out of each vertex, we have to look at the vertices which represent the separate entities. For example, the bottom vertex looks like this:


There are two possibilities to complete the Hamiltonian circuit.


All six vertices are similar, so there are $2^{6}$ possibilities. Multiplying the $2^{6}$ ways times the 6 ways from the $\boldsymbol{B}$ connectors, we have the 384 possibilities for two totally independent hamiltonian circuits on four bells. None of these graphs can be symmetric since the graph of $\boldsymbol{B}$ connectors is not symmetric. Notice this in the following example of two independent hamiltonian cycles.


Without symmetry, a pattern in the letters can not be found, and the changes can not be represented in a short algebraic form like the ones given earlier.

Going on to five bells causes even a bigger problem. The graph of the possibilities, alone, is complicated. There are seven possible changes from sequence to sequence: (12), (34), (45), (23), $(12)(34),(12)(45)$, and $(23)(45)$. So each of the 120 vertices is connected to seven other vertices. That is a total of $\mathbf{4 2 0}$ edges. It takes $\mathbf{1 2 0}$ of these edges to complete a Hamiltonian circuit. Since 420 is not a multiple of $\mathbf{1 2 0}$, independent Hamiltonian circuits can not be found that use all the edges.

For six bells there are $\mathbf{7 2 0}$ changes, $\mathbf{1 2}$ transitions, and $\mathbf{4 3 2 0}$ edges on the graph. Because there is an even number of transitions, the number of edges is a multiple of the number of changes. So it seems likely that independent Hamiltonian circuits that use all of the edges exist. Since there is no easy method for determining which graphs are Hamiltonian and each graph must be considered individually, determining whether there exists more than one Hamiltonian circuit on one graph can not be found using a theorem. Therefore, finding a method to find the Hamiltonian circuits is part of the problem.

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## ELEMENTARY SYMMETRIC POLYNOMIALS, AN INTUITIVE APPROACH

 WITH APPLICATIONS TO COLLEGE ALGEBRA AND BEYONDDaniel Replogle

This article will present an intuitive introduction to the elementary symmetric polynomials and describe some of their uses. It is written with the good college algebra student in mind. Everything it contains should be accessible to the student who has mastered college algebra.

Elementary symmetric polynomials are used frequently in advanced courses in algebra, and are not usually presented until then. However, seeing them earlier might give undergraduates a greater sense of the structure of algebra. Just as journals frequently use methods from advanced calculus to throw light upon topics from standard calculus courses, so topics from advanced algebra may sometimes be used to throw light upon topics from earlier algebra courses. An early introduction to elementary symmetric polynomials is seeing how each term in a polynomial depends upon the roots of that polynomial. This dependence will be shown and stated as a theorem, though no proof of this result will be given.

Consider the following products:

$$
\begin{equation*}
(x-a)(x-b)=x^{2}-(a+b) x+a b . \tag{1}
\end{equation*}
$$

$$
\begin{align*}
(x-a)(x-b)(x-c) & =\left(x^{2}-(a+b) x+a b\right)(x-c)  \tag{2}\\
& =x^{3}-(a+b) x^{2}+a b x-c x^{2}+c(a+b) x-a b c \\
& =x^{3}-(a+b+c) x^{2}+(a b+a c+b c) x-a b c .
\end{align*}
$$

(3) $(x-a)(x-b)(x-c)(x-d)=\left[z^{3}-(a+b+c) x^{2}+(a b+a c+b c) x-a b c\right](x-d)$ $=x^{4}-(a+b+c+d) x^{3}+(a b+a c+a d+b c+b d+c d) x^{2}$
$-(a b c+a b d+a c d+b c d) x+a b c d$.
(4)
$(x-a)(x-b)(x-c)(x-d)(x-e)=x^{5}-(a+b+c+d+e) z^{4}$
$+(a b+a c+a d+a e+b c+b d+b e+c d+c e+d e) z^{3}$
$-(a b c+a b d+a b e+a c d+a c e+a d e+b c d+b c e+b d e+c d e) x^{2}$
$+(a b c d+$ abce + abde + acde $+\boldsymbol{b c d e}) \boldsymbol{x}-\boldsymbol{a b c d e}$.
Noting the above pattern, let $a_{1}, a_{2}, a_{3}, a_{4}, \ldots, a_{k}$ be the zeros of a monic polynomial (a polynomial where the coefficient of the highest degree term is 1 ) of degree $k$. Then define:

$$
\begin{aligned}
s_{1 k}= & a_{1}+a_{2}+a_{3}+\ldots+a_{k} \\
s_{2 k}= & a_{1} a_{2}+a_{1} a_{3}+\ldots+a_{1} a_{k}+a_{2} a_{3}+\ldots+a_{2} a_{k}+a_{3} a_{4}+\ldots+a_{k-1} a_{k} \\
s_{3 k}= & a_{1} a_{2} a_{3}+a_{1} a_{2} a_{4}+\ldots+a_{1} a_{2} a_{k}+a_{1} a_{3} a_{4}+\ldots+a_{1} a_{3} a_{k}+a_{1} a_{k-1} a_{k}+ \\
& \ldots+a_{2} a_{3} a_{4}+\ldots+a_{k-2} a_{k-1} a_{k} \\
s_{4 k}= & a_{1} a_{2} a_{3} a_{4}+\ldots+a_{k-3} a_{k-2} a_{k-1} a_{k} \\
s_{5 k}= & a_{1} a_{2} a_{3} a_{4} a_{5}+\ldots+a_{k-4} a_{k-3} a_{k-2} a_{k-1} a_{k} \\
s_{k k}= & a_{1} a_{2} a_{3} a_{4} \ldots a_{k-2} a_{k-1} a_{k} .
\end{aligned}
$$

The preceding definea the elementary symmetric polynomials on $k$ letters. The following might be helpful to keep in mind:
$\boldsymbol{s}_{\mathbf{1 k}}$ is the sum of all of the zeros,
$\boldsymbol{s}_{\mathbf{2 k}}$ is the sum of all of the disjoint pairwise product8 of zeros,
$\boldsymbol{s}_{\mathbf{3} k}$ is the sum of all of the disjoint $\mathbf{3}$-wise products of zeros, etc.
With all of this in mind and recalling the pattern observed above, we have the following theorem (which can be proved rigorously, for those who desire to do so):

Theorem: The monic polynomial $p(x)$ of degree $k$, with zeros $a_{1}, a_{2}, a_{3}, \ldots, a_{k}$ is given by:

$$
p(x)=x^{k}-s_{1 k} x^{k-1}+x_{2 k} x^{k-2}-s_{3 k} x^{k-3}+\ldots+(-1)^{r} s_{r k} x^{k-r}+\ldots \pm s_{k-1, k} x \mp s_{k k} .
$$

## Applications:

1. Find the monic polynomial $p(x)$ with zeros $1,2, \sqrt{2},-\sqrt{2}$.

$$
\begin{aligned}
s_{14} & =1+2+\sqrt{2}-\sqrt{2}=3 . \\
s_{24} & =1(2)+1(\sqrt{2})+1(-\sqrt{2})+2(\sqrt{2})+2(-\sqrt{2})+(\sqrt{2})(-\sqrt{2}) \\
& =2+\sqrt{2}-\sqrt{2}+2 \sqrt{2}-2 \sqrt{2}-2=0 . \\
s_{34} & =1(2)(\sqrt{2})+1(2)(-\sqrt{2})+1(\sqrt{2})(-\sqrt{2})+2(\sqrt{2})(-\sqrt{2}) \\
& =2 \sqrt{2}-2 \sqrt{2}-2-4=-6 . \\
s_{44} & =1(2)(\sqrt{2})(-\sqrt{2})=-4 .
\end{aligned}
$$

So, $\boldsymbol{p}(\boldsymbol{x})=\boldsymbol{x}^{4}-3 z^{3}+0 x^{2}-(-6) x+(-4)=x^{4}-3 z^{3}+6 x-4$.
2. Find the polynomial $p(x)$ having zeros $i \sqrt{2},-i \sqrt{2}$, and 2 with $p(3)=2$.

$$
\begin{aligned}
s_{13} & =2+i \sqrt{2}-i \sqrt{2}=2 \\
s_{23} & =i \sqrt{2}(-i \sqrt{2})+(i \sqrt{2})(2)-i \sqrt{2}(2) \\
& =2+2 i \sqrt{2}-2 i \sqrt{2}=2 \\
s_{33} & =i \sqrt{2}(-1 \sqrt{2})(2)=4
\end{aligned}
$$

so, $p(x)=r\left(x^{3}-2 x^{2}+2 x-4\right)=r x^{3}-2 r x^{2}+2 r x-4 r$.

So, $\boldsymbol{p}(3)=l l r$. But $\boldsymbol{p}(3)=2$, making $\boldsymbol{r}=2 / 11$. Thus

$$
p(x)=\frac{2}{11} x^{3}-\frac{4}{11} x^{2}+\frac{4}{11} x-\frac{8}{11}
$$

3. Show that if $p$ is a prime number, $p>2$, then the sum of the $p$ th roots of unity is zero, and their product is one.

Each of the pth roots of unity satisfies the equation $\boldsymbol{x}^{p}-1=0$. Further, this polynomial equation has $p$ roots. So, if $\boldsymbol{a}_{1}, \boldsymbol{a}_{2}, \ldots, a_{p}$ are the roots of $\boldsymbol{x}^{p}-1=0$, then

$$
\begin{aligned}
\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{p}\right) & =x^{p}-s_{1 p} x^{p-1}+s_{2 p} x^{p-2}-\ldots+s_{p-1, p} x-s_{p p} \\
& =x^{p}-1
\end{aligned}
$$

(We can be definite about the choice of signs because a prime $>2$ is necessarily odd.) It follows that $\boldsymbol{s}_{1 p}=0$ and $-\boldsymbol{s}_{p p}=-1$. Thus $\boldsymbol{s}_{1 p}=\mathbf{0}$, so the sum of the pth roots of unity is zero (where $p \mathbf{i s}$ a prime $>2$ ).
Also $\boldsymbol{s}_{p p}=\mathbf{1}$, so the product of the pth roots of unity ( $p$ a prime $>2$ ) is one.

## Additional Comments:

The usual method for finding the polynomials in Examples 1 and 2 might be quicker and simpler, but it will not help one to solve problems like that in Example 3. Also, to me, the usual method seems to be just a bit too tedious and it fails to reveal any structure. For comparison, here is how Example $l$ is usually solved:
la. Find the monic polynomial $p(x)$ with zeros $1,2, \sqrt{2},-\sqrt{2}$.

$$
\begin{aligned}
p(x) & =(x-1)(x-2)(x-\sqrt{2})(x-(-\sqrt{2})) \\
& =(x-1)(x-2)(x-\sqrt{2})(x+\sqrt{2}) \\
& =\left(x^{2}-3 x+2\right)\left(x^{2}-2\right) \\
& =x^{4}-3 x^{3}+6 x-4 .
\end{aligned}
$$

A careful look at this and Example 1 above, I think, reveals that the method using symmetric polynomials is somewhat less tedious and reveals more structure.

Daniel Replogle prepared this article shortly after completing his master's at St. Louis University. He is currently a graduate student at the State University of New York at Albany.

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## SOLUTIONS TO ANTIDERIVATIVES

USING A HYPERBOLIC FUNCTIONAL TRANSFORMATION

## Timothy Holland

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Most textbooks for elementary integral calculus include a section titled 'miscellaneous substitutions.' Among the types of problems that these sections generally consider are thcae which involve finding the antiderivatives of rational functions of $\sin (@)$ and $\boldsymbol{\operatorname { c o s }}(\beta)$. The traditional method of solving some of these problems uses the following substitution [1]:

$$
\sin (\beta)=\frac{2 x}{1+x^{2}}, \quad \cos (\beta)=\frac{1-x^{2}}{1+x^{2}}, \quad \text { and } \quad d \beta=\frac{2}{1+x^{2}} d x .
$$

However, the use of the exponential and hyperbolic functions offers an alternative method for solving these integrals. It has the additional benefit of providing a pedagogical tool for expanding the use of the hyperbolic functions in elementary calculus.

> We begin by noting the following:

Theorem 1. $\int \operatorname{sech}(x) d x=2 \tan -^{\prime}\left(\mathrm{e}^{\mathrm{x}}\right)$
Proof:

$$
\int \operatorname{sech}(x) d x=\int \frac{2}{e^{x}+e^{-x}} d x=\int \frac{2 e^{x}}{e^{2 x}+1} d x=2 \tan ^{-1}\left(e^{x}\right)
$$

## Corollary 1. If $\beta=2 \tan ^{-1}\left(e^{x}\right)$, then $d \beta=\operatorname{sech}(x) d x$.

Corollary 1 indicates that if $\beta=2 \boldsymbol{\operatorname { t a n }}^{-\mathbf{1}}\left(\boldsymbol{e}^{\boldsymbol{x}}\right)$, then there is a relationship between the hyperbolic functions of $\boldsymbol{x}$ and the trigonometric functions of $\boldsymbol{\beta}$. We can now establish expressions for the other hyperbolic functions.

Theorem 2. $\sinh (x)=-\cot (@)$.
Proof:

$$
\begin{aligned}
\sinh (x) & =\sinh \left\{\ln \left[\tan \left(\frac{\beta}{2}\right)\right]\right\} \\
& =\frac{1}{2}\left[\tan \left(\frac{\beta}{2}\right)-\cot \left(\frac{\beta}{2}\right)\right] \\
& =\left[\frac{1-\cos (\beta)}{\sin (\beta)}-\frac{\sin (\beta)}{1-\cos (\beta)}\right] \\
& =\frac{[1-\cos (\beta)]^{2}-\sin ^{2}(\beta)}{2 \sin (\beta)[1-\cos (\beta)]} \\
& =\frac{1-2 \cos (\beta)+\cos ^{2}(\beta)-1+\cos ^{2}(\beta)}{2 \sin (\beta)[1-\cos (\beta)]} \\
& =\frac{-2 \cos (\beta)[1-\cos (\beta)]}{2 \sin (\beta)[1-\cos (\beta)]} \\
& =-\cot (\beta)
\end{aligned}
$$

Theorem 3. $\cosh (x)=\csc (\beta)$.
Proof:

$$
\begin{aligned}
\cosh (x) & =\cosh \left(\ln \left[\tan \left(\frac{\beta}{2}\right)\right]\right) \\
& =\frac{1}{2}\left[\tan \left(\frac{\beta}{2}\right)+\cot \left(\frac{\beta}{2}\right)\right] \\
& =\frac{1}{2}\left[\frac{1-\cos (\beta)}{\sin (\beta)}+\frac{\sin (\beta)}{1-\cos (\beta)}\right] \\
& =\frac{[1-\cos (\beta)]^{2}+\sin ^{2}(\beta)}{2 \sin (\beta)[1-\cos (\beta)]} \\
& =\frac{1-2 \cos (\beta)+\cos ^{2}(\beta)+1-\cos ^{2}(\beta)}{2 \sin (\beta)[1-\cos (\beta)]} \\
& =\frac{2[1-\cos (\beta)]}{2 \sin (\beta)[1-\cos (\beta)]} \\
& =\csc (\beta)
\end{aligned}
$$

The following corollaries are the direct results of Theorems 2 and 3:
Corollary 2. $\tanh (x)=-\cos (\beta)$.
Proof: $\quad \tanh (x)=\begin{gathered}\sinh (x) \\ \cosh (x)\end{gathered}=-\cot (\beta)=-\cos (\beta)$.
Corollary 3. $\operatorname{sech}(x)=\sin (@)$.
Proof: $\quad \operatorname{sech}(x)=\begin{gathered}1 \\ 1 \\ \end{gathered}=\sin (@)$.
Corollary 4. $\operatorname{csch}(x)=-\tan (\beta)$.
Proof: $\quad \operatorname{csch}(x)=\frac{1}{\sinh (x)}=\frac{-1}{\cot (\beta)}=-\tan (@)$.
The following examples illustrate how to apply these transformations to some antiderivatives:
Example 1.

$$
\begin{aligned}
\int \frac{d \beta}{1+\sin (@)-\cos (\beta)} & =\int \frac{\operatorname{sech}(x) d x}{1+\operatorname{sech}(x)+\tanh (x)} \\
& =\int \frac{d x}{1+\cosh (x)+\sinh (x)} \\
& =\int \frac{d x}{1+e^{x}} \\
& =x-\ln \left|1+e^{2}\right|+C \quad[\text { Let } x=\ln (y) \text { and integrate by partial fractions.] } \\
& =\ln \left|\tan \left(\frac{\beta}{2}\right)\right|-\ln \left|1+\tan \left(\frac{\beta}{2}\right)\right|+C \\
& =\ln \left|\frac{\tan \left(\frac{\beta}{2}\right)}{1+\tan \left(\frac{\beta}{2}\right)}\right|+C .
\end{aligned}
$$

## Example 2.

$$
\begin{aligned}
\int \frac{d \beta}{3-2 \cos (\beta)} & =\int \frac{\operatorname{sech}(x) d x}{3+2 \tanh (x)} \\
& =\int \frac{d x}{3 \cosh (x)+2 \sinh (x)} \\
& =\int \frac{2 e^{x} d x}{5 e^{2 x}+1} \\
& =\frac{2}{\sqrt{5}} \int \frac{d\left(\sqrt{5} e^{x}\right)}{5 e^{2 x}+1} \\
& =\frac{2}{\sqrt{5}} \tan ^{-1}\left(\sqrt{5} e^{x}\right)+C \\
& =\frac{2}{\sqrt{5}} \tan ^{-1}\left[\sqrt{5} \tan \left(\frac{\beta}{2}\right)\right]+C .
\end{aligned}
$$

## Example 3.

$$
\begin{aligned}
\int \frac{d \beta}{5+4 \sin (\beta)} & =\int \frac{d x}{5 \cosh (x)+4} \\
& =\int \frac{d x}{\frac{5 e^{2 x}+5}{2 e^{x}}+4} \\
& =\int \frac{2 e^{x} d x}{5 e^{2 x}+8 e^{x}+5} \\
& =\left(\frac{2}{5}\right) \int \frac{d\left(e^{x}\right)}{\left(e^{x}+\frac{4}{5}\right)^{2}+\left(\frac{3}{5}\right)^{2}} \\
& =\left(\frac{2}{5}\right) \tan ^{-1}\left(\frac{4+5 e^{x}}{3}\right)+C \\
& =\frac{2}{3} \tan ^{-1}\left[\frac{4+5 \tan \left(\frac{\beta}{2}\right)}{3}\right]+C
\end{aligned}
$$

## Example 4.

$$
\begin{aligned}
\int \sec (\beta) d \beta & =-\int \operatorname{coth}(x) \operatorname{sech}(x) d x \\
& =-\int \frac{\cosh (x) d x}{\sinh (x) \cosh (x)} \\
& =-\int \frac{d x}{\sinh (x)} \\
& =-\int \frac{\sinh (x) d x}{\sinh ^{2}(x)}
\end{aligned}
$$

$$
\begin{aligned}
& =-\int \frac{d[\cosh (x)]}{\cosh ^{2}(x)-1} \\
& =\operatorname{coth}^{-1}[\cosh (x)] \\
& =\operatorname{coth}^{-1}[\csc (\beta)] \\
& =\ln |\sec (\beta)+\tan (\beta)|+C .
\end{aligned}
$$

Example 5.

$$
\int \csc (\beta) d \beta=\int \cosh (x) \operatorname{sech}(x) d x=\int d x=x+C=\ln \left[\tan \left(\frac{\beta}{2}\right)\right]+C
$$

In summary, as an alternative to

$$
\sin (\beta)=\frac{2 x}{1+x^{2}}, \quad \cos (\beta)=\frac{1-x^{2}}{1+x^{2}}, \quad \text { and } \quad d \beta=\frac{2 d x}{1+x^{2}},
$$

the substitutions

$$
\sin (\beta)=\operatorname{sech}(x), \quad \cos (\beta)=-\tanh (x), \quad \text { and } \quad d \beta=\operatorname{sech}(x) d x
$$

can be used to solve many antiderivatives involving rational functions of $\sin (\beta)$ and $\cos (\beta)$.

## Reference

1. S. M. Farrand and N. J. Poxon, Calculus, Harcourt, Brace, Jovanovich, 1984.

Timothy Holland prepared this paper while he was teaching at St. Jude High School and enrolled in a master's program at Alabama State University.

## CHANGES OF ADDRESS

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## THE EASTER DATE PATTERN

## Richard L. Fmnci

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An intriguing date is that of Easter. By a seemingly erratic pattern, it makes its appearance on the calendar each year. Sometimes in March and at other times in April. Sometimes very early and at other times, quite late. But always on Sunday. The date of Easter for a given year was fixed and at other times, quite late. But always on Sunday. The date of Easter for a given year was fixed
in A.D. $\mathbf{3 2 5}$ by the Council of Nicaea. In this ancient decree, Easter became accordingly the first Sunday after the first full moon on or after the vernal (spring) equinox. As both lunar and solar cycles are involved as well as the day of the week pattern, challenging mathematical problems come to light. In particular, how is Easter to be calculated for a given year? What too can be said about the frequency of the various Easter dates and their subtle, hardly noticeable calendar patterns?

Several well known formulas come to mind in pursuit of these and similar mathematical questions. One such approach is a variation on the Easter formula as given by the great mathematician Carl Friedrich Gauss (1777-1855). Before elaborating however on computational techniques, a brief historical note is in order.

Much diversity characterized the Easter observance pattern prior to the year 325. Even in later years, recurring problems arose as a consequence of the far-reaching calendar change of 1582. In that year, the ancient Julian calendar was replaced by the modern calendar of Pope Gregory XIII (the Gregorian calendar). The motivation for the change was essentially one of alignment of dates with seasons. Controversy surrounded the new calendar's introduction; various nations were likewise slow in adoption. Although the Gregorian calendar is the one in present worldwide civil use, some today, for ecclesiastical purposes, celebrate Easter in accordance with the ancient Julian calendar. Coincidentally, the Julian and Gregorian calculations of Easter will occasionally give the same date (as happened in 1865, 1905, and 1954 for example).

A Metonic cycle from ancient times essentially equated $\mathbf{2 3 5}$ full moons with $\mathbf{1 9}$ vernal equinoxes. Hence, a time period of $\mathbf{1 9}$ years denotes the cycle in which the sun and moon patterns eventually prove commensurable. (The Athenian astronomer Meton devised a calendar pattern in 432 B.C. whereby the new moons repeat in 19 year cycles.) More precisely, an integral multiple of one cycle coincided with an integral multiple of the other. However, the calendar reformers of 1582 realized a very slight discrepancy in this equation, namely, the one which blended the lunar cycle with the 19 year solar pattern. The assumption of equality was implicit in the Julian calculation of Easter. The Gregorian correction incorporates the fact that the Dominical Letter of a year (the symbol for the year's first Sunday) and the Golden Number (a given year's place in the overall 19 year cycle) will not, in and of themselves, give the exact Easter date.

Because of the complexity of the relationship between the lunar cycle and the solar cycle, various Easter formulas are restricted to but a single century. Each century thus has its own full moon sequence. Such a complexity of relationships is accounted for concisely by appropriate references to time called EPACTS, (The word "epact" stems from the Greek and denotes the "age" of the moon in days at the start of a new year.) Mathematicians can verify (see the Kluepfel reference) that there are exactly 30 epacts as well as 30 sets of correspondences involving epacts and Golden Numbers. It is therefore possible to construct an Easter formula or set of Easter correspondences which will prove accurate for all time

It can also be shown that the Gregorian calendar's period, namely, its perfect date-day cycle of repetition, is exactly 400 years. Hence, as December 25, 1994 falls on a Sunday, so will Christmas Day 400 years later. By examining any 400 year interval of time, it can be establiihed that, for example, the thirteenth of a month falls more often on Friday than any other day of the week.

Likewise, it can be proved that Presidential Inauguration Day occurs more often on Sunday than any other day of the week (as last happened in 1985 and will next occur in 2013).

A more relevant point is that the Easter period can also be calculated. What then is the smallest interval of time which implies consistently a perfect cycle of Easter date repetition? Suppose a key symbol is associated with each of the thirty Golden Number and Epact associations mentioned above. Let these key symbols for convenience be the integers 0 through 29. It can be shown (see Kluępfel) that any cycle of $\mathbf{1 0 0}$ centuries has a new key symbol (number). This accounts for $\mathbf{3 0}$ (100) or $\mathbf{3 0 0 0}$ centuries. Yet each of these century intervals is associated with one of the nineteen possible Golden Numbers, no two of which are alike. Accordingly, the Easter period becomes 19(3000) centuries or $5,700,000$ years. Acknowledging thus this Easter period of $5,700,000$ years, a tabulation of Easter date frequencies becomes possible. Note among other things that $\mathbf{5 , 7 0 0 , 0 0 0}$ is divisible by $\mathbf{4 0 0}$, in which case the day of the week pattern (Sunday restriction) is maintained.

## VARIATION ON THE GAUSSIAN EASTER FORMULA

The sequence of steps which permits calculating the date of Easter for a particular year is given below

1. $\quad \frac{\text { year }}{19}=\mathrm{A}$ plus remainder $\boldsymbol{B}$
2. $\frac{\text { year }}{100}-C \quad$ plus remainder $D$
3. $\quad \underline{\boldsymbol{G}}=\mathrm{E} \quad$ plus remainder $F$
4. $\frac{\boldsymbol{C}+8}{25}-\boldsymbol{G}$ plus discarded remainder
5. $\frac{C+1-G}{3}=H$ plus discarded remainder
$\frac{19 B+C+15-(E+H)}{30}=$ quotient (discard) plus remainder $Z$
6. $\quad \frac{\boldsymbol{D}}{\boldsymbol{4}}=\boldsymbol{K} \quad$ plus remainder L
$\frac{2 F+2 K+32-(Z+L)}{7}=$ quotient (discard) plus remainder $N$
7. $\quad \frac{B+11 Z+22 N}{451}=P \quad$ plus discarded remainder
8. $\frac{Z+N+114-7 P}{31}=Q \quad$ plus remainder R

Then Q denotes the month and $\mathrm{R}+1$ denotes the day on which Easter falls for a given year,

An illustration reinforces the formula. To calculate Easter for the year 1998, the following letter values are obtained

1. $A=105 \quad B=3$
2. $C=19 \quad D=98$
3. $E=4 \quad F=3$
4. $G=1$
5. $H=6$
6. $Z=21$
7. $K=24 \quad L=2$
8. $N=0$
9. $P=0$
10. $Q=4 \quad R=11$

As Easter is given by month Q and day $R+1$, the actual date of Easter for 1998 is April 12. It is also the most common Easter date of the twentieth century (occurring six times).

## HOW EARLY AND HOW LATE?

Easter may occur as early as March 22. This last occurred in 1818 and before that in 1761 and 1693. Such an early occurrence is actually a rarity. Easter will not come so early again in this century or in the next. Not until the year 2285 will Easter fall on March 22.

At the other extreme, Easter may come as late as April 25. Its last such occurrence was in 1943 and, prior to that, in 1886. Easter will next occur on this latest possible date in the year 2038.

All of the above dates relate to the Gregorian calendar. The calculation of Easter dates according to other schemes frequently deviates from this as mentioned earlier. For example, Easter Sunday in Russia in 1989 occurred on April 30. Such a late date stem from Julian results which are assigned corresponding Gregorian dates.

As noted, there are 35 possible dates for Easter Sunday according to the modern Gregorian calendar. Ten such dates are in March; the remaining twenty-five are in April.

## THE TWENTIETH CENTURY

The Easter Sunday frequency pattern for the twentieth century appears in the graph below. Note that all possible Easter dates are represented except March 22 and April 24.

EASTER SUNDAY FREQUENCY
the twentieth century (1901-2000)


## FIRST 2000 YEARS OF THE GREGORIAN CALENDAR

All possible Easter dates appear in this 2000 year period of time. The least frequent date is March 22, occurring but 13 times. The most frequent are April 4 and April 10, each occurring 83 times.

## ENTIRE EASTER PERIOD OF 5,700,000 YEARS

Not only do all possible Easter dates appear in this vast period of time (that of a perfect Easter date cycle of repetition), it is also easy to tell which of the dates is the least frequent and which is the most Sequent. Note that March 22 (the least frequent date) occurs 27, 550 times. The runner-up is April 25 (occurring 42,000 times). Note likewise that April 19 is the most frequent; it occurs 220, 400 times. This year, 1992, Easter falls on its most frequent date. The midde Easter
date (from March 22 to April 25) is April 8 ; it occurs 192, 850times. Moreover, the average frequency date (from March 22 to April 25) is April 8; it occurs 192,850times. Moreover, the
is obtained by dividing $5,700,000$ by 35 . This average is approximately 162,857 .

EASTER SUNDAY FREQUENCY
the first two-thousand years of the Gregorian Calendar (1583-3582)


EASTER SUNDAY FREQUENCY
the entire Easter period of $5,700,000$ years


TABLE III

Various questions and conjectures arise in examining a long list of consecutive Easter dates． Consider for example the thousand year Easter listing given below．

Easten surian

| far Easter | Vtar Easter | Vtar Easter | Year Easter | uer Easter | Year Easter | Year Easter | Year Easter | Year Easter | Year fuster |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  | 1603 | 101 | 103s apr 10 | 1866 har 26 |  |  | 19 |
| 1619 dp | 1611 apr | 1112 Apr | 1613 Apr | 1614 tar 10 | 1115 Apr 19 | 1616 Air | 1117 Mar 21 | 1618 Apr | 31 |
| 1629 Apr 19 | 1621 Apr 11 | 1622 Mar 27 | 1633 Apr 16 | 1624 Apr | 1623 KT 10 | 1626 Apr 12 | 1627 App | 1628 Agr 23 | 1629 Apr 15 |
| 1630 Mar 31 | 1631 Apr 20 | 1632 Apr 11 | 1633 Mar 27 | 1634 Apr 16 | 1633 Apr | 1636 Mar 23 | 1137 Apr I2 | 1638 mpr | 1639 Apr 24 |
| 1640 Apr | 1141 har 31 | 1622 Apr 20 | 1143 Apr | 1644 har 27 | 1649 ${ }^{\text {ded }} 16$ | 1646 | 1647 Apr 21 | 164 APr | 14 |
| 1150 Apr | 1651 Apr 9 | 1152 mar 31 | 1633 Apr 13 | 1654 dy | 1635 mar 2 | 1654 Apr 16 | 1157 納 | 1658 Apr 21 | 1659 App 13 |
| 1660 Mar 21 | 1661 Apr 17 | 1622 apr | IU3 Mar 25 | 1664 app 13 | 1659 Apr | 1664 ApP 25 | 1667 Apr 10 | 1688 mpr | 1669 Apr 21 |
| 1670 Apr | 1171 Mar | 1672 App 17 | 1173 Apr | 1674 kar 25 | 1673 Apr | 1476 Apr | 1177 Aper 18 | 1478 App | 1679 |
| 1680 Apr 21 | 1681 Apr 6 | 1882 har 29 | 1183 Apr 18 | 1694 Mpr | 1685 Apr 2 | 16as Apr 14 | 1687 har | 1689 Apr 1 | 1639 Apr 10 |
| 1699 mar 2 | 1691 Apr 15 | 1692 Apr | 1693 Mar 22 | IH4 Apr II | 1695 Apr | 1696 Apr 22 | 1697 Mpr | 1698 Mar 10 | 1699 Apr 19 |
| 1700 Apr 11 | 1701 Mar 27 | 1702 Apr | 1103 Agr | 1704 Mar | 1705 Apr 12 | 1701 Apr | 1707 Apr 24 | 1708 Apr | 1709 Mar 31 |
| 1710 Apr 2 | 1714 Apr | 1712 h | 1713 hpr | 1714 Apr | 1715 Apr 21 | 1716 Apr 12 | 1717 Bar 28 | 1718 Apr 1 | 179 Apr |
| 1720 Mar | 1721 Apr 13 | 172 Apr | 1723 mar 28 | 1724 Apr | 1721 apr 1 | 1726 tor | 1727 Apr | 1720 HiT 2 | 1729 |
| 1770 Agr | 1731 Mar 25 | 1732 apr 13 | 1733 Apr 5 | 1734 Apr 29 | 1735 Apr 10 | 1736 Apr | 1737 Apr 21 | 1738 Apr | 1739 mar 21 |
| 1740 Apr 17 | 1741 apr 2 | 1742 Mar 25 | 1743 Apr 14 | 174 Apr 5 | 1745 Apr 18 | 1746 Apr 10 | 1747 Abr | 1748 Apr | 174 |
| 1750 Mzr 29 | 1751 Apr It | 1752 Apr | 1753 Apr 22 | 1754 Apr 14 | 1755 har 30 | 176 | 1757 |  |  |
| 1760 Apr | 1761 Ma | 1762 apr 11 | 1763 Apr | 1764 Apr 22 | 1715 Apr | 176 shar | 1767 Apt 19 | 1768 Apr | 176 |
| 1770 Apr 1 | 1771 Mar 31 | 172 Abr 19 | 1773 Apr 11 | 1774 Apr | 1775 apr 16 | 1776 | 1777 Mar | 1778 Agr 19 | 1779 Aor 4 |
| 1780 Mar 26 | 1781 Apr 15 | 1102 kar 31 | 1783 Apr 20 | 1784 Apr 11 | 1785 mar 27 | 1786 Apr 11 | 1787 mpr | 1788 fit 23 | 1789 Apr 12 |
| 1790 Apr 4 | 1791 apr 24 | 1792 Apr | ${ }^{1793}$ mar 31 | 1799 Abr 20 | 1793 App | 1798 har 27 | 1797 far 16 | 1799 Apr | 1749 mar 24 |
| 1800 Apr 13 | 1801 Apr | 1802 Apr 18 | 1803 Agr 10 | 1804 Apr | 1805 Apr 14 | 1801 apr | 1807 IT 79 | 1808 Apr | 1809 |
| 1810 Apr | 1811 Apr 14 | 1812 Mar 29 | 1813 Apr 18 | 1814 App 10 | 1815 mar 21 | 1811 Apr 14 | 1817 Apr | 1818 Mar | 1819 Apr 11 |
| 1820 Apr | 1821 Apr 22 | 1822 Apr | 1823 har 10 | 1824 Apr | 1823 Apr | 1824 Mar 21 | 1827 Apr 15 | 1878 apr | 1821 Apr 19 |
| 1830 Apr 11 | 1811 Apr | 1832 Apr | 1833 Apr | 1834 Mar | 1835 Apr 19 | 18\％${ }^{\text {d }}$ | 1837 har 21 | 1838 har 15 | 1839 har 31 |
| 1840 Apr 19 | 1841 apr 11 | 1842 | 1843 Apr 16 | 1844 Apr | 1845 Har 23 | 1846 Ap | 1847 App | 184 A Apr ${ }^{\text {2J }}$ | 1849 Apr |
| 1850 Mar | 1851 Apr | 1852 Agr 11 | 1853 har 27 | 1854 Agr 16 | 1855 Apr | 1836 Mar 23 | 1837 apr 12 | 1850 Apr |  |
| 1860 ap | 1861 Mar 31 | 1862 Apr 20 | 1863 Apr | 1814 mar 27 | 1815 Apr 16 | 1864 fpr | 1867 Air 21 | 1658 Apr 12 | 1869 fir 28 |
| 1870 Apr | 1871 Apr | 1812 nar 31 | 1873 Agr ${ }^{13}$ | 1874 dor | 1875 Mar | 1871 Apr 11 | 1877 Apr | 1878 Apr 21 | 1879 der is |
| 1880 Mar | 1881 Apr | 1882 apr | 1883 far 25 | 1884 Apr 13 | 1885 Apr | 1881 Apr 25 | IN 7 Apr 10 | 188 Apr 1 | ce9 Ar |
| 1890 Apr | 1811 Mar | 1812 Apr | 1893 Apr | 1894 har 25 | 1895 Apr 14 | 1896 Apr | mi App 18 | 1898 Apr 10 | 1899 aper 2 |
| 1900 Apr | HOI Apr | 1902 Mar | 1803 Apr 12 | 1704 Apr | 1905 Apr 23 | 1906 Apr 15 | 1907 Mar 31 | 1900 Agr 19 |  |
|  | 1911 1921 Aer 18 27 | 1912 Apr | H13 har 23 | 1914 apr 12 | 1915 Apr | 1916 Apr 23 | 1917 Apr | 1918 har 31 | 20 |
| 1930 for 20 | 1991 har 27 | ${ }_{1932 \mathrm{Apr}} 193 \mathrm{Har}$ | 1933 Agr | ${ }_{1}^{1924}$ Apr 20 | 1935 Apr | $1926 \mathrm{Agr}{ }^{4}$ | 1927 for 17 | 1978 apr | 1929 har 31 |
| 1940 fir 24 | 1914 Apr 13 | 1942 Apr | 1943 | 1941 Apr | （193 | 1936 apr 12 | 19377 Mat 28 | 1939 Apr 17 | 1939 apr 9 |
| 1950 Apr | 1931 Mar 25 | 1992 Apr 13 | 1953 Apr | 1951 Apr 18 | 1935 apr 10 | 1956 | 1997 | 1989 | 17 |
| Hid Apr | 1961 Apr | 1962 Apr 22 | 1963 apr 14 | 1984 Mar 29 | 1965 A | 1965 apr 10 | 1967 Mar 21 | 1968 Apr 14 |  |
| 1970 Mar | 1971 Apr 11 | 1972 Apr 2 | H13 App 11 | 1974 Apr 14 | 1975 har 30 | 1976 apr 18 | 1971 Apr 10 | 1978 Har 21 |  |
| 1980 Apr | 1981 apr 19 | 1982 Apr 11 | 1993 Apr | 1994 Apr 22 | 1985 Apr 7 | 1984 har 1 | 1887 Apr 19 | 1988 Abr | 1989 for 21 |
| 1990 Apr 15 | 1991 mar 31 | 1992 apr 19 | 1993 apr 11 | 1994 Apr 3 | 1995 Apr 11 | 1996 Apr | 1997 har 30 | 1 178 Apr 12 | 1999 far 4 |
| 2000 Apr 23 2010 Apr | 2001 Apr 15 | 2002 nar | 2003 Apr | 2004 Apr 11 | 2005 har 27 | 2008 Apr 16 | 2007 Apr | 2008 mar | 2008 Apr 12 |
| 2020 Apr 12 | 2021 Apr | 2022 Apr 17 | 2013 Mar 31 | 2014 apr 20 | 2015 Apr |  |  | 2018 | 2019 Agr 21 |
| 2050 Apr 21 | 2031 Apr 13 | 2032 har 28 | 2033 apr 17 | 2034 Apr | 2035 nar 25 | 2023 Apr | 2027 har | 2028 Apr 11 | 2079 spr 1 |
| 2010 hpr 1 | 2041 Apr 21 | 2042 Aor | 2043 Sir 29 | 2044 Apr | 2045 Apr 9 | 2036 Apr 13 | 2037 Apr | ${ }_{2}^{2038} \mathrm{Apr} 25$ | 2039 Apr 10 |
| 2050 Apr 10 | 2051 Apr 2 | 2052 Apr 21 | 2053 Apr 6 | 2054 Mre It | 2058 Apr 1 | ${ }^{20 \%}$ Apr | 2057 Apr |  | 2049 Apr 18 205 V Mar 30 |
| 2060 Apr 18 | 2061 Apr 10 | 2012 \＃rar 26 | 2063 Apr 15 | 2064 Apr | 2065 Mar 29 | 2066 Apr It | 2067 Apr | 2068 Apr | 2069 |
| 70 har 30 | 2071 Apr 19 | 2072 Apr 10 | 2073 Mar 21 | 2074 Apr 15 | 2075 Apr 7 | 2076 Apr 19 | 2071 Apr II | 2078 Apr | ${ }_{2079}^{2069}$ apr 23 |
| go Apr | 2081 far 30 | 2082 Apr 19 | 2083 Apr | 2004 kar 21 | 2085 apr 15 | 2086 Bar | 2087 Apr 20 |  |  |
| 2090 Agr 11 | 2091 Apr | 30 | 2093 Apr 12 | 2094 Apr 1 | 2095 Apr 24 | 2096 Apr 15 | 2097 mar | 2098 | 2099 Air I2 |

EASIER Stuoar

|  |  |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |  |  |  |  |
|  | 2111 mar 2 |  |  |  |  | 211 |  |  |  |
| 1270 Ap | 2121 Apr | 2122 Mar | 2123 |  |  |  | 2127 | 2121 Apr 18 |  |
| ， | 2131 Apr is | 2132 Apr | 2133 Apr 19 |  |  |  |  |  |  |
|  |  |  |  |  |  |  |  |  |  |
| 150 Apr |  | 215 |  | 2134 Mar |  |  | 215 | 2158 |  |
| 1160 Mar 23 | 216 |  | 211 |  |  |  |  |  |  |
|  | 2171 Apr | 2172 Apr | 2173 Apr |  |  | 2171 |  |  |  |
| did | 2181 Ap | 2182 Apr 2 |  |  |  | 2185 | 218 | 2188 |  |
|  | 219 |  |  |  |  |  |  |  |  |
| neo hor |  | 2 Agr | 2293 | 2204 Ap | 2705 |  | mi Aor 19 |  |  |
| 2210 Ap |  |  | 2213 apr |  |  |  |  |  |  |
| 2220 AD | 2221 | Mar | 2223 | 2224 | 272 | 22 | 227 |  |  |
| ： | ：311 | 2332 for 目 | 273 | 223 |  |  |  | 22：8 |  |
|  |  |  | 2243 | 2244 Mar 3 |  | 22 | 234 | 2248 Apr 16 |  |
| 230 | 229 | 52 | 2235 | 254 | 2ns | 25 | 2257 | 2258 | 2258 Apr |
|  | 2261 | 2262 Ap | $2 \mathrm{2bJ} \mathrm{Mar}$ | 226 | 2265 |  |  |  |  |
|  | 2771 | 2222 | 2273 | 22 |  |  |  |  |  |
|  | 2881 | 2322 M | 2283 A | 22 |  |  |  |  |  |
| Har 30 |  | 2792 A | ， |  |  |  |  |  |  |
|  |  |  | 103 |  |  | 2306 | 307 |  |  |
|  | 131 | 12 | 2313 | 2314 | 2315 | 2316 | 2317 |  |  |
| \％ | 2341 | 222 | 23 | 2321 |  |  |  |  |  |
|  | 2331 | 2 | 2133 | 23 |  | 2336 A |  |  |  |
| （14）Apr 14 | 2341 |  | 2313 |  |  |  |  |  |  |
| So Bar 26 |  | 2332 | 2351 |  | 2355 | 1 |  |  |  |
|  | 231 | 2162 | 2333 n | 2354 cpr |  |  |  |  |  |
| 3 ab 30， 19 |  |  |  | 2374 Mar |  |  |  |  |  |
| SO Mar 23 |  | 2：82 |  |  |  |  |  |  |  |
|  |  |  |  |  |  |  |  |  |  |
| 110 hor 15 | 2411 | 2412 | ${ }^{4} 113$ |  |  |  |  |  |  |
| Apr 5 |  |  | 2423 |  | 2425 | 242 |  |  |  |
| 0 O fipr 14 | 2431 | 2432 | 2433 | 2414 |  | 436 |  |  |  |
| 10 A | 2413 |  | 析 |  |  |  |  |  |  |
| （5i）Apr 3 |  |  | 2451 |  |  |  |  |  |  |
| 60 do |  | 2412 Apr | 2433 | 2464 |  | 216 |  |  |  |
|  |  |  |  |  |  | 2176 Ad |  |  |  |
| ${ }^{4} \mathrm{M}$ | 2481 nd | 22 | 2133 n |  |  |  |  |  |  |
| 190 mp | 2191 Ma |  |  |  |  | 2598 |  |  |  |
|  | 2501 A | 2502 Mar 26 <br> 2512 Apr | ${ }_{2513}^{2503 ~ A p r ~} 15$ | 2514 | 2315 | 2516 |  |  |  |
|  | 2521 Ma | 2522 | 2523 | 2524 | 2525 | 2526 | 3 |  |  |
| So A | 2331 Apr | 352 Ma | 3333 Apr | 253 | 2535 har | 23564 | 253 |  |  |
|  |  | 2542 Ab | 8313 | 254 |  | 246 | 2541 Apr |  |  |
| A |  |  | 2553 |  | 2555 Apr 13 | 2556 |  |  |  |
|  |  |  |  |  |  | 2566 Apr 6 |  |  |  |
|  |  |  |  |  |  |  |  |  |  |
|  |  |  |  |  |  |  |  |  |  |

A typical question concerns the month pattern. In particular, "Can Easter occur in March in both of any two consecutive years?" The table above strongly suggests the answer is "no." In order to be certain, one must examine the entire Easter period, or, by consideration of the formula, establish that Q cannot be 3 (the March number) for consecutive year numbers. By a computer analysis of the Q question for the entire Easter period of $5,700,000$ years, the answer of "no" is verified. Accordingly, Easter cannot come in March two years in a row. Another question suggested by the long Easter listing is "Can Easter occur on corresponding Sundays in any two consecutive years?" Corresponding Sundays in consecutive years are those Sundays which differ by exactly 52 weeks. For example, April 9, 1950 and April 8, 1951 are corresponding Sundays. The first is an Easter date; the second is not. Knowing for instance that April 12, 1998 is Easter (as calculated earlier), can one now assert for afact that neither April 13, Sunday, 1997 nor April 11, Sunday, 1999 is an Easter date? Once again, by a computer analysis of the Easter period of $5,700,000$ years, the answer is "consecutive Easters cannot occur on corresponding Sundays." Other questions likewise stem from the Easter listing. A few are included here for the purpose of additional exploration.

1. It appears that in consecutive years, Easter dates can be no closer datewise than 8 days. Consider for example the Easter dates April 18, 1965, and April 10, 1966. Note too that the earlier year always seems to contain the later date number. Are these suggested patterns valid?
2. The shortest interval of time separating a given Easter and its next like date occurrence is five years. This happened, for example, on (Easter) March 29, 1959, and (Easter) March 29, 1964. What is the greatest interval of time separating a given Easter and its next like date occurrence? Note that long intervals of time are suggested by examining Easter lists. Among them are the Easter dates March 22, 1818, and March 22, 2285, which span an interval of 467 years.
3. Is it possible for a decade to consist entirely of April Easters? What is the greatest number of consecutive April Easters possible?
4. Ponce de Leon, the European discoverer of Florida, gave the area its name ("Florida" or "flowery Easter") on Easter Sunday in the year 1513. What was the exact date of this Julian calendar Easter and could it be the same as its projected Gregorian date counterpart, namely, April 6?
5. Can two like date Easters be exactly $\mathbf{4 0 0}$ years apart? Recall that the period of the Gregorian calendar is 400 years
The above are but a few of the many questions contained in the mathematical subtlety of the Easter date pattern.

The Gregorian calendar will prove many times out of line with the seasons in the course of $5,700,000 y$ ears. Actually, the calendar proves a day in error every 3323 years. In less than 100,000 years, the present calendar's marginal inconsistency with the seasons will magnify and measure roughly a month. Accordingly, it must be stressed that the computations above rest on the assumptions implicit in the Gregorian calendar's construction (a calendar likely to be modified or abandoned in the years ahead). Still, in its present, highly familiar form, it affords an opportunity for the mathematically curious to explore an intriguing pattern of numbers and number relationships.

Appreciation is expressed to Victor Gummersheimer and Johnny Lai for their computer assistance in the preparation of this manuscript.

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## ERRATA

The following errors appeared in the Fall, 1991, issue of the Pi Mu Epsilon Journal:
Page 280 On line 4 and again on line 5, " $a c-\mathrm{bd} \neq 0$ " should have read " $a d-\mathrm{bc} \# 0$."
Page 292 "Theorem 2" and "Example 10" should have been "Theorem 7" and "Example 11."
'Page 296 Line 7 should have read " $365.25 / 365.25+11365.25=24$ hours $/ X$."
Page 297 On line 5, "then $\mathrm{O}_{1}, \mathrm{O}_{2}$, and $\mathrm{O}_{3}$ " should have read "then $\mathrm{O}_{1} \mathrm{O}_{2}$ and $\mathrm{CO}_{3}$."
Page 299 On line 4 of Theorem 2, " $A_{1}^{\prime}=C_{2}$ " should have read " $A_{n}^{\prime}=C_{2}$."

The Editor apologizes for any problems that these errors might have caused.

## EXTENDING A FAMILIAR INEQUALITY

## Norman Schaumberger

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The following problem appeared on the 1973 USA Mathematical Olympiad:
Prove that if $\boldsymbol{a}, \boldsymbol{b}$, and c are positive real numbers, then

$$
\begin{equation*}
a^{a} b^{b} c^{c} \geq(a b c)^{(a+b+c) / 3} \tag{1}
\end{equation*}
$$

A simpler version,

$$
\mathrm{a}^{\mathrm{a}} \mathrm{~b}^{\mathrm{b}} \geq(a b)^{(a+b) / 2}
$$

is a familiar exercise in a number of texts.
The usual proofs of (1) use a not particularly simple elementary argument or Jensen's inequality. [See M. S. Klamkin, USA Mathematical Olympiads 1972-1986, MAA, 1988, p. 81.]

We start our proof by noting that if $x>0$, then

$$
\begin{equation*}
x \ln x \geq x-1 \geq \ln x \tag{2}
\end{equation*}
$$

with equality iff $\mathrm{x}=1$.
We use (2) to extend (1) and then to obtain an important limiting relation for the power mean.
The right side of (2) follows immediately from the observation that $f(x)=x=1-\operatorname{In} x$ has an absolute minimum at $\mathrm{x}=1$ because $\mathrm{f}^{\mathrm{t}}(\mathrm{x})=1-1 / \boldsymbol{x}=0$ iff $\mathrm{x}=1$, and $f^{\prime \prime}(x)=1 / \mathrm{x}^{2}$ is positive for $\mathrm{x}>0$. If we now replace x by $1 / x$ in $\mathrm{x}-1 \geq \ln \mathrm{x}$ we get $\mathrm{x} \ln \mathrm{x} \geq \mathrm{x}-1$ which completes (2).

Let $\mathrm{A}=(\mathrm{a}+\mathrm{b}+\boldsymbol{c}) / \boldsymbol{3}$ and substitute $\mathrm{x}=\boldsymbol{a} / \boldsymbol{A}, \boldsymbol{x}=\boldsymbol{b} / \boldsymbol{A}$, and $\mathrm{x}=\boldsymbol{c} / \boldsymbol{A}$ successively into (2). Adding gives

$$
\frac{a}{A} \ln \frac{a}{A}+\frac{b}{B^{=}} \ln \frac{b}{A}+\frac{c}{A} \ln \frac{c}{A} \geq \frac{a+b+c}{A} \cdot 3 \geq \ln \frac{a}{A}+\ln \frac{b}{A^{A}}+\ln \frac{c}{A^{A}}
$$

Hence

$$
\ln \left[(a / A)^{a / A}(b / A)^{b / A}(c / A)^{c / A}\right] \geq 0 \geq \operatorname{In}\left(\frac{a b c}{A^{3}}\right)
$$

or

$$
\begin{equation*}
\left(\frac{a^{a} b^{b} c^{c}}{A^{a+b+c}}\right)^{1 / A} \geq 1 \geq \frac{a b c}{\frac{2}{2}} \tag{3}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
a^{a} b^{b} c^{c} \geq A^{a+b+c} \geq(a b c)^{(a+b+c) / 3} \tag{4}
\end{equation*}
$$

This double inequality gives (1) and somewhat more. Also, there is equality iff $\boldsymbol{a} / \boldsymbol{A}=1, \boldsymbol{b} / \boldsymbol{A}=1$, and $c / A=1$. That is, iff $\mathrm{a}=\mathrm{b}=\mathrm{c}$.

The power mean, $M_{r}$, of order $r$ is defined by

$$
M_{r}=\left(\frac{1}{n} \sum_{i=1}^{n} a_{i}^{r}\right)^{1 / r},
$$

where $a_{i}>0(i=1,2, \ldots, n)$ and $\boldsymbol{r} \# 0$ are real numbers. Thus $M_{1}$ and $M_{2}$ are the arithmetic mean and root mean square. If $n=3$, then $\boldsymbol{M}_{\boldsymbol{r}}^{\boldsymbol{r}}=\left(\mathrm{a}^{\mathrm{r}}+\mathrm{b}^{r}+\mathrm{c}^{\mathrm{r}}\right) / 3$.

Putting $\mathrm{x}=\boldsymbol{a}^{\boldsymbol{r}} / \boldsymbol{M}_{\boldsymbol{r}}^{\boldsymbol{r}}, \mathrm{x}=\boldsymbol{b}_{\boldsymbol{r}}^{\boldsymbol{r}} / \boldsymbol{M}_{\boldsymbol{r}}^{\boldsymbol{r}}$, and $\mathrm{x}=\boldsymbol{c}_{\boldsymbol{r}}^{\boldsymbol{r}} / \boldsymbol{M}_{\boldsymbol{r}}^{\boldsymbol{r}}$ in (2) and adding gives

$$
\begin{aligned}
\frac{a^{r}}{M_{r}^{r}} \ln \frac{a^{r}}{M_{r}^{r}}+\frac{b^{r}}{M_{r}^{r}} \ln \frac{b^{r}}{M_{r}^{r}}+\frac{c^{r}}{M_{r}^{r}} \ln \frac{c^{r}}{M_{r}^{r}} & \geq \frac{a^{r}+b^{r}+c^{r}}{M_{r}^{r}}-3 \\
& \geq \ln \frac{a^{r}}{M_{r}^{r}}+\ln \frac{b^{r}}{M_{r}^{r}}+\ln \frac{c^{r}}{M_{r}^{r}}
\end{aligned}
$$

In a similar way to that used to get (3), it follows that

$$
\left[\frac{\left(a^{r}\right)^{a^{r}}\left(b^{r}\right)^{b^{r}}\left(c^{r}\right)^{c^{r}}}{\left(M_{r}^{r}\right)^{a^{r}+b^{r}+c^{r}}}\right]^{1 / M_{r}^{r}} \geq 1 \geq \frac{a^{r} b^{r} c^{r}}{M_{r}^{3 r}} .
$$

If $\boldsymbol{r}>0$, raising to the $\boldsymbol{M}_{\boldsymbol{r}}^{\boldsymbol{r}} / \boldsymbol{r}$ power gives

$$
\frac{a^{a^{r}} b^{b^{r}} c^{c^{r}}}{M_{r}^{a^{r}+b^{r}+c^{r}}} \geq 1 \geq\left(\frac{a b c}{M_{r}^{3}}\right)^{M_{r}^{r}}
$$

or

$$
\begin{equation*}
a^{a^{r}} b^{b^{r}} c^{c^{r}} \geq M_{r}^{a^{r}+b^{r}+c^{r}} \geq(a b c)^{\left(a^{r}+b^{r}+c^{r}\right) / 3} \tag{5}
\end{equation*}
$$

More generally, the same kind of argument can be used to get

$$
\begin{equation*}
a_{1}^{a_{1}^{r}} a_{2}^{a_{2}^{r}} \ldots a_{n}^{a_{n}^{r}} \geq M_{r}^{a_{1}^{\tau}+a_{2}^{r}+\ldots+a_{n}^{r}} \geq\left(a_{1} a_{2} \ldots a_{n}\right)^{\left(a_{1}^{\tau}+a_{2}^{r}+\ldots+a_{n}^{r}\right) / n} \tag{6}
\end{equation*}
$$

Inequality (4) is a special case of (5). If $\boldsymbol{r}<0$, the inequalities in (5) are reversed. Thus, for example, if $\boldsymbol{r}=-1$, then $\boldsymbol{M}_{-1}$ is the harmonic mean and (5) becomes

$$
a^{\frac{1}{a}} b^{\frac{1}{b}} c^{\frac{1}{c}} \leq M_{-1}^{\left(\frac{1}{a}+\frac{1}{b}+\frac{1}{c}\right)} \leq(a b c)^{\left.\left(\frac{1}{a}+\frac{1}{b}+\frac{1}{a}\right) / 3\right)}
$$

If $\boldsymbol{r}=0, M_{\boldsymbol{r}}$ is not defined. However, the geometric mean, $\sqrt[3]{a b c}$, is usually denoted by $M_{0}$, The standard proof that

$$
\begin{equation*}
\lim _{r \rightarrow 0} M_{r}=\sqrt[n]{a_{1} a_{2} \ldots a_{n}} \tag{7}
\end{equation*}
$$

uses L'Hospital's Rule and the theory of exponential functions. [See Hardy, Littlewood, and Polya, Inequalities, Cambridge University Press, Cambridge, 1952, p. 15.]

Equation (7) follows at once from the observation that (5) can be written as

$$
\left(a^{a^{r}} b^{b^{r}} c^{c^{r}}\right)^{1 /\left(a^{r}+b^{r}+c^{r}\right)} \geq M_{r} \geq(a b c)^{1 / 3}
$$

If $\boldsymbol{r} \rightarrow 0,\left(a^{a^{r}} b^{b^{r}} c^{c^{r}}\right)^{1 /\left(a^{r}+b^{r}+c^{r}\right)}$ tends to $\sqrt[3]{a b c}$ and we get (7) for $n=3$. The general case can be proved in an analogous manner using (6).

## PROOF OF THE CONVERGENCE OF A SEQUENCE OF RADICALS

## Andrew Cusumano <br> Gret Neck, NY

The purpose of this paper is to investigate the expression

$$
S=\sqrt{1+2 \sqrt{1+3 \sqrt{1+4 \sqrt{1+\ldots}}}} .
$$

We first note that $S$ represents a sequence $\left\{a_{k}\right\}_{k=1}^{\infty}$ of real numbers. Checking a few terms, either by hand or with an easily written computer program, leads us to conjecture that the sequence converges to 3 . In order to prove this, we first show that it is a monotone increasing sequence that is bounded above by 3 . Thus, $\lim _{k \rightarrow \infty} \boldsymbol{a}_{k}=\mathrm{a}$ exists and $\mathbf{a} \leq 3$. Finally, we show that $\mathbf{a}=3$.
$\mathrm{T}_{0}$ write $S$ as an increasing sequence and see that it is bounded above by the number 3 , we note that

$$
\begin{aligned}
& a_{1}=\sqrt{1+2}<\sqrt{1+2(4)}=3 \\
& a_{2}=\sqrt{1+2 \sqrt{1+3}}<\sqrt{1+2 \sqrt{1+3(5)}}=3 \\
& a_{3}-\sqrt{1+2 \sqrt{1+3 \sqrt{1+4}}}<\sqrt{1+2 \sqrt{1+3 \sqrt{1+4(6)}}}=3 \\
& \vdots \\
& a_{k}=\sqrt{1+2 \sqrt{1+3 \sqrt{1+\ldots+(k-1) \sqrt{1+k \sqrt{1+(k+1)}}}}} \\
& \quad<\sqrt{1+2 \sqrt{1+3 \sqrt{1+\ldots+(k-1) \sqrt{1+k \sqrt{1+(k+1)(k+3)}}}}=3} \\
& \vdots
\end{aligned}
$$

Since $\left\{a_{k}\right\}_{k=1}^{\infty}$ is an increasing bounded sequence, bounded above by 3 , we know that $\lim _{k \rightarrow \infty} \boldsymbol{a}_{\boldsymbol{k}}=\mathbf{a}$ exists and $\mathbf{a} \leq 3$.

Thus we can write

$$
\begin{equation*}
a=\sqrt{1+2 \sqrt{1+3 \sqrt{1+4 \sqrt{1+\ldots}}}} \tag{1}
\end{equation*}
$$

To prove that $\mathbf{a}=3$, we construct a sequence $\left\{b_{n}\right\}_{n=1}^{\infty}$ such that $\boldsymbol{b}_{\boldsymbol{n}}<\mathbf{a} \leq \mathbf{3}$ for every $\boldsymbol{n}$, and then show that $\lim _{n \rightarrow \infty} \boldsymbol{b}_{\boldsymbol{n}}=3$. In order to construct our $\boldsymbol{b}_{\boldsymbol{n}}$, we must get some idea of how far each $\boldsymbol{a}_{\boldsymbol{k}}$ is from 3. If we consider, for example, just the part

$$
\sqrt{1+6 \sqrt{1+7 \sqrt{1+8 \sqrt{1+\ldots}}}},
$$

this is clearly at least as big as $\boldsymbol{x}$, where

$$
x=\sqrt{1+6 \sqrt{1+6 \sqrt{1+6 \sqrt{1+\ldots}}}}
$$

We then note that $\sqrt{1+6 x}=$ a;. By using the quadratic formula on the equation $\boldsymbol{x}^{2}-6 \boldsymbol{x}-\mathbf{1}=0$, we can see that $\mathbf{x}>6$. We can then compare

$$
\sqrt{1+2 \sqrt{1+3 \sqrt{1+4 \sqrt{1+5(6)}}}}
$$

to

$$
\sqrt{1+2 \sqrt{1+3 \sqrt{1+4 \sqrt{1+5(6+1)}}}}
$$

which is exactly equal to 3 . We can now generalize this approach in order to construct the $\boldsymbol{b}_{\boldsymbol{n}}$. Notice that

$$
d_{n}=\sqrt{1+(n+1) \sqrt{1+(n+2) \sqrt{1+(n+3) \sqrt{1+\ldots}}}}>x_{n}
$$

where

$$
x_{n}=\sqrt{1+(n+1) \sqrt{1+(n+1) \sqrt{1+(n+1) \sqrt{1+\ldots}}}}=\sqrt{1+(n+1) x_{n}} .
$$

Since $\boldsymbol{x}_{\boldsymbol{n}}^{2}-(\mathbf{n}+1) \boldsymbol{x}_{\boldsymbol{n}}-\mathbf{1}=0$, the quadratic formula shows that $\boldsymbol{x}_{\boldsymbol{n}}>\boldsymbol{n}+1$.
By replacing $\boldsymbol{d}_{\boldsymbol{n}}$ by $\mathbf{n}+\mathbf{1}$ in (1), we get $\boldsymbol{b}_{\boldsymbol{n}}$, where

$$
b_{n}=\sqrt{1+2 \sqrt{1+3 \sqrt{1+\ldots+(n-1) \sqrt{1+n(n+1)}}}}
$$

Since $\mathrm{n}+\mathbf{1}<\boldsymbol{x}_{\boldsymbol{n}}<d_{n}$, we have $\boldsymbol{b}_{\mathrm{n}}<\mathbf{a}$. Thus, for every $\boldsymbol{n}, \boldsymbol{b}_{n}<\mathbf{a} \leq \mathbf{3}$.
To complete our proof, we need only show that $\lim _{n \rightarrow \infty} b_{n}=3$. To do this, we first note that if $\mathbf{0}<\boldsymbol{x}<\boldsymbol{y}, 0<w$, and $0<\boldsymbol{u}<\mathbf{1}$, then $\mathbf{u}<\sqrt{\boldsymbol{u}}$, and

$$
\frac{x}{y}<\frac{1+w x}{1+w y}<\frac{\sqrt{1+w x}}{\sqrt{1+w y}} .
$$

It follows that

$$
\frac{n+1}{n+2}<\frac{\sqrt{1+n(n+1)}}{\sqrt{1+n(n+2)}}<\frac{\sqrt{1+(n-1) \sqrt{1+n(n+1)}}}{\sqrt{1+(n-1) \sqrt{1+n(n+2)}}}<\ldots<\frac{b_{n}}{3}<\frac{a}{3} \leq \frac{3}{3}=1
$$

where

$$
3=\sqrt{1+2 \sqrt{1+3 \sqrt{1+\ldots+(n-1) \sqrt{1+n(n+2)}}}}
$$

Since $\lim _{n \rightarrow \infty} \frac{n+1}{n+2}=1$, it follows that $\lim _{n \rightarrow \infty} b_{n}=3$. Thus a $=3$, and our proof is complete.

## A CLOSED FORM FOR A FAMILY OF SUMMATIONS

## Russell Euler <br> Northwest Missouri State University

Let p be an integer such that $\mathrm{p} \geq 2$. It can be shown, by using the asymptotic relationship

$$
\binom{n}{p} \sim \frac{n^{p}}{p!} \quad \text { as } n \rightarrow \infty
$$

from page 33 of [2], that

$$
\begin{equation*}
I(p)=\sum_{n=p}^{\infty}\binom{n}{p}^{-1} \tag{1}
\end{equation*}
$$

converges. The series (1) was evaluated in [1] by using partial sums. In this paper, a closed form for (1) will be obtained by using special functions. The special functions that will be used are reviewed first.

The gamma function is denoted by $\Gamma(x)$ and defined by

$$
\Gamma(x)=\int_{0}^{\infty} e^{-t} t^{x-1} d t
$$

for $\boldsymbol{x}>0$. The gamma function has the property that $\Gamma(x+1)=x \Gamma(x)$ provided $\boldsymbol{x}$ is neither zero nor a negative integer. In particular, for $n=0,1,2, \ldots, \quad \Gamma(n+1)=n!$.

The factorial function is defined by

$$
(a)_{n}=a(a+1) \cdots(a+n-1) \text { for } n \geq 1 \text { and }(a)_{0}=1 \text { for } a \neq 0 .
$$

In particular, $n!=(1)_{n}$ and, from page 9 of $[2]_{1}(a)_{n+k}=(a+n)_{k}(a)_{n}$.
The (Gaussian) hypergeometric function is denoted by ${ }_{2} F_{1}(a, b ; c ; \boldsymbol{x})$ and defined by

$$
{ }_{2} F_{1}(a, b ; c ; x)=1+\sum_{n=1}^{\infty} \frac{(a)_{n}(b)_{n} x^{n}}{(c)_{n} n!}
$$

provided c is neither zero nor a negative integer. If none of the parameters a, $\boldsymbol{b}$, or c are zero or a negative integer, it is known that this series is absolutely convergent for $|x|<1$, divergent for $|z|>1$, and is absolutely convergent for $|x|=1$ provided $\mathrm{a}+\mathrm{b}-\mathrm{c}<0$.

To evaluate (1), first notice that $I(p)$ can be written as

$$
\begin{align*}
I(p) & =p!\sum_{n=p}^{\infty} \frac{(n-p)!}{n!} \\
& =p!\sum_{n=0}^{\infty} \frac{n!}{(n+p)!} \tag{2}
\end{align*}
$$

However,

$$
\begin{align*}
\frac{n!}{(n+p)!} & =\frac{(1)_{n}}{(1)_{p+n}} \\
& =\frac{(1)_{n}}{(1+p)_{n}(1)_{p}} \tag{3}
\end{align*}
$$

Substituting (3) into (2) and simplifying yields

$$
I(p)=\sum_{n=0}^{\infty} \frac{(1)_{n}}{(1+p)_{n}}
$$

Hence,

$$
\begin{align*}
I(p) & =\sum_{n=0}^{\infty} \frac{(1)_{n}(1)_{n}}{(1+p)_{n} n!} \\
& ={ }_{2} F_{1}(1,1 ; 1+p ; 1) . \tag{4}
\end{align*}
$$

It has been shown in [3], page 49, that if $\operatorname{Re}(c-a-b)>0$ and if $c$ is neither zero nor a negative integer,

$$
{ }_{2} F_{1}(a, b ; c ; 1)=\frac{\Gamma(c) \Gamma(c-a-b)}{\Gamma(c-a) \Gamma(c-b)} .
$$

Therefore, identity (4) becomes

$$
\begin{equation*}
I(p)=\frac{\Gamma(p+1) \Gamma(p-1)}{\Gamma(p) \Gamma(p)} . \tag{5}
\end{equation*}
$$

Since $\Gamma(x)=(x-1) \Gamma(x-1),(5)$ simplifies to give

$$
\sum_{n=p}^{\infty}\binom{n}{p}^{-1}=\frac{p}{p-1}
$$

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## INQUIRIES

Inquiries about certificates, pins, posters, matching prize funds, support for regional meetings, and travel support for national meetings should be directed to the Secretary-Treasurer, Robert M. Woodside, Department of Mathematics, East Carolina University, Greenville, NC 27858, 919-7576414.

# WHAT IS 'LOCALLY COMPACT ${ }^{\text {n }}$ ? 

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Each textbook in Topology has its own way of defining what it means for a space to be locally compact. Some authors make an effort to give an equivalent characterization under some additional assumptions about the topological space (see [2]). Essentially there are four wncepts with the name of local compactness and the relations among these have been only partially studied. Even though local compactness, by its very title, is a local property (recall that a property is said to be local if it can be specified for any single point in the space), there has been only global study (that is, a study of the spaces where the property is assumed for every point in the space) of it in the literature. In this paper, we study it locally at a point. Implications among these concepts will be discussed at a particular point. Moreover, we present examples to help understand the impossibility of reverse implications.

Throughout this paper X represents an arbitrary topological space and $\boldsymbol{x}$ denotes a fixed point of X. Schnare [3] discussed two definitions of local compactness, which are rephrased here to define them as properties of space X at a point $\boldsymbol{x}$ as follows: A topological space X is called weakly locally compact, or simply w-compact, at $\boldsymbol{x}$ iff there is a compact neighborhood of $\boldsymbol{x}$ in the space X . X is called mildly locally compact, or m -compact, at $\boldsymbol{x}$ iff there is a neighborhood of $\boldsymbol{x}$ whose closure is compact. A topological space is said to be 1 -compact iff it is 1 -compactat each of its points where $\mathbf{1}$ " is " $w$ ", " $m$ ", or any other letter that makes sense in the following discussion. Schnare [3] showed that a w-compact space is m-compact iff the closure of any compact set is compact. Later, Gross [4] introduced a third definition of local compactness, which is modified here as a property at a particular point $\boldsymbol{x}$. A space X is called bit locally compact, or b -compact, at $\boldsymbol{x}$ iff each neighborhood of $\boldsymbol{x}$ contains a compact neighborhood of $\boldsymbol{x}$.

It is well known that all these wncepts are equivalent in Hausdorff spaces and regular spaces. In fact, in such spaces, these are equivalent to one more concept called strongly locally compact. X is said to be strongly locally compact, or s-compact, at $\boldsymbol{x}$ iff each neighborhood of $\boldsymbol{x}$ contains a compact closed neighborhood of $\mathbf{z}$. The particular choice of terminology becomes apparent after observing that s -compact is strongest, w-compact is the weakest, and b -compact, m-compact lie in between for any general spaces. That is, we have the following implications in any general topological space X at the point $\boldsymbol{x}$ :


These implications are strict. Moreover, b-compact and m-compact are incomparable in a general topological space.

Even though compact spaces are obviously m-compact (and thus w-compact), compactness does not imply either s-compactness or b-compactness. Consider the one-point compactification of the space $Q$ of rational numbers. This is a $T_{1 \frac{1}{2}}$-space (a space in which each compact set is closed). It
can be shown that it is neither s-compact nor b-compact. This example also tells us that even in compact $\boldsymbol{T}_{1 \frac{1}{2}}$-spaces

$$
\text { m-compact } \xrightarrow{\longrightarrow} \text { b-compact }
$$

at $\boldsymbol{x}$. However, b-compact certainly implies m-compact in $T_{1 \frac{1}{2}}$-spaces. In fact, this implication holds even under a weaker assumption on the topological space. ${ }^{2}$ To explain this assumption, we need the following definition. A space is called an $\mathbf{R}$-space iff the closure of a compact set is compact. Clearly any regular or $\boldsymbol{T}_{\mathbf{1 \frac { 1 } { 2 }}^{-}}$(hence $\boldsymbol{T}_{\mathbf{2}^{-}}$) space is an R -space. since m -compactness at $\boldsymbol{\tau}$ is equivalent to the statement that there is a compact closed neighborhood of $\boldsymbol{x}$ in X , it is immediate that in any R-space

```
b-compact \(\longrightarrow\) m-compact
```

at $\boldsymbol{x}$ and

$$
\text { m-compact } \longleftrightarrow \text { w-compact }
$$

at $\boldsymbol{x}$. Of course, b-compact does not imply m-compact in general spaces. Gross [4] has an example of a b-compact normal space which is not m-compact. An easy example is the following: Consider an infinite set $X$ with a distinguished point $\boldsymbol{x}$ in which a set is declared to be open if it is either empty or it contains $\mathbf{z}$. This is a $\mathrm{T}_{0}$-space (a space in which distinct points have distinct closures) which is b-compact at $\boldsymbol{x}$ but not m -compact.

An infinite set with cofinite topology reveals that even in compact, $\boldsymbol{T}_{\mathbf{1}}$ - and R-spaces
b- and m-compact $\nrightarrow \mathrm{s}$-compact at $\boldsymbol{z}$.
But in-compact and s-compact at $\mathfrak{x}$ are equivalent in a topological space which is $\boldsymbol{T}_{\mathbf{2}}$ at $\boldsymbol{x}$. A space X is called $\boldsymbol{T}_{\mathbf{2}}$ at $\boldsymbol{x}$ iff for any point y of X different from $\boldsymbol{x}$, there exist two disjoint open sets $\boldsymbol{G}$ and H in X containing y and $\boldsymbol{x}$, respectively. It is easy to verify that a topological space X is $\boldsymbol{T}_{\mathbf{2}}$ at a point $\boldsymbol{x}$ iff to each compact set A not containing $\boldsymbol{x}$ there correspond two disjoint open sets $L$ and M such that $\mathrm{A} \subseteq L$ and $\boldsymbol{x} \boldsymbol{6} \mathrm{M}$.

Let us show that if X is $\boldsymbol{T}_{\boldsymbol{2}}$ at $\boldsymbol{x}$ and m-compact at $\boldsymbol{x}$ then it is s-compact at $\boldsymbol{x}$. Let $\boldsymbol{G}$ by any open set containing $z$. Let N be a compact closed neighborhood of $\boldsymbol{x}$ (by m-compactness at $\boldsymbol{x}, \mathrm{N}$ exists). Write $\mathrm{A}=\mathrm{N} \cap \boldsymbol{G}^{\prime}$, where $\boldsymbol{G}^{\prime}$ represents the complement of $\boldsymbol{G}$. Clearly A is a closed subset of the compact space N and hence compact. Since $\boldsymbol{x} \notin \mathrm{A}$ and X is $\boldsymbol{T}_{\mathbf{2}}$ at $\boldsymbol{x}$, there are disjoint open sets $\boldsymbol{L}$ and $\boldsymbol{M}$ such that $\mathrm{A} \subseteq L$ and $\boldsymbol{x} \in \mathrm{M}$. Now $M \subseteq L^{\prime}$ and

$$
\bar{M} \subseteq L^{\prime} \subseteq A^{\prime}=N^{\prime} \cup G
$$

(the bar indicates the closure of the set), which means $\bar{M} \cap \mathrm{~N} \subseteq \boldsymbol{G}$. Thus $\boldsymbol{H}=\mathrm{M} \cap \mathrm{N}^{\circ}\left(\mathrm{N}^{0}\right.$ is the interior of N ) is an open set containing $\mathscr{z}$ and

## $\bar{H} \subseteq \bar{M} \cap \bar{N}=\bar{M} \cap N \subseteq G$.

Moreover, $\bar{H}$, being a closed subset of compact set $N$, is compact. This shows that X is s-compact at z .

At this point note that, for any R -space that is $\boldsymbol{T}_{\mathbf{2}}$ at $\boldsymbol{x}$ the implications

$$
\text { w-compact } \longrightarrow \text { m-compact } \longrightarrow \text { s-compact }
$$

hold at $\boldsymbol{x}$, hence all compactness concepts are equivalent.
Notice that a space which is s-compact at $\boldsymbol{x}$ is regular at $\boldsymbol{x}$; that is, to each open set $\boldsymbol{G}$ containing $\boldsymbol{x}$ there corresponds an open set H such that $\boldsymbol{x} \in \mathrm{H} \subseteq \bar{H} \subseteq G$. In fact, this property of the space assures the equivalence of all these concepts. To prove this, let us assume that $\mathrm{X} \mathbf{1 8}$. w-compact at $\boldsymbol{x}$ and regular at $\boldsymbol{x}$. We show that X is s-compact at $\boldsymbol{x}$. Let $\boldsymbol{G}$ be any open set
containing $\boldsymbol{z}$. Since X is w-compact at $\boldsymbol{z}$, there is a compact neighborhood N of $\boldsymbol{z}$. Put $\mathrm{A}=\boldsymbol{G} \cap \mathrm{N}^{0}$. Then by regularity at $\boldsymbol{z}$, there exists an open set $H$ such that

## $x \in H \subseteq \bar{H} \subseteq A$

Clearly $\overline{\boldsymbol{H}}$ is compact (because a closed subset of a compact space is compact) and $\overline{\boldsymbol{H}} \subseteq$ G. Thus X is s -compact at $\boldsymbol{\varepsilon}$. Thus all these concepts are equivalent in spaces which are regular $\overline{\mathrm{a}} \mathrm{t} \boldsymbol{x}, T_{2}$ at $z$ with R-property, or Hausdorff spaces.

We close our discussion with an analysis of some of the standard properties of local compact spaces. Clearly any local compactness is $\boldsymbol{x}$ is closed hereditary (i.e., preserved under closed subspaces). However, only $\mathbf{s}$ - and b-compactness are open hereditary (i.e., preserved under open subspaces). The one point compactification of the space Q is m -compact (hence w -compact) in which the open set Q is neither m -compact nor w -compact. A w-compact dense subset B of a $T_{1+}$-space $X$ is open. Indeed, suppose $\mathbf{b} \in B$. Since $\mathbf{B}$ is w-compact, there exist a compact subset $\mathbf{C}$ of B and an open subset G of X such that $\mathrm{b} \in B \cap G \subseteq \mathrm{C} . C$ is closed in X , because X is a $T_{1 \frac{1}{2}}$-space. Since B is dense in X and $G$ is open in $X, \quad \bar{G}=\overline{B \cap G}$. Thus

$$
b \in G \subseteq \bar{G}=\overline{B \cap G} \subseteq \bar{C}=C \subseteq B
$$

This shows that $B$ is a neighborhood of $\mathbf{b}$. This being true for any $\boldsymbol{b} \in \mathrm{B}, \mathrm{B}$ is open.
References

1. J. L. Kelly, General Topology, D. Van Nostrand Company, Inc., Princeton, New Jersey, (1966).
2. A. Wilansky, Topology for Analysis, Ginn and Company, Massachusetts, (1970).
3. P. S. Schnare, "Two Definitions of Local Compactness," American Mathematical Monthly 72 (1965), pp 764-765
4. J. L. Gross, "A Third Definition of Local Compactness," American Mathematical Monthly 74 (1976), pp 1120-1122

## ATTENTION FACULTY ADVISORS

To have your chapter's report published, send copies to Robert M. Woodside, SecretaryTreasurer, Department of Mathematics, East Carolina University, Greenville, NC 27858 and to Richard L. Poss, Editor, St. Norbert College, De Pere, WI 54115.

## ON THE CALCULUS OF RESIDUES

## Prem N. Bajaj <br> The Wichita State University

In this note we give a paradox in the calculation of residues at a pole; a paradox in the sense that an incorrect procedure gives a correct answer. Some of the well-known examples of this type are: the incorrect cancellation of 6 in 16/64, of 9 in 19/95, or of 2 in $(1+\boldsymbol{x})^{\mathbf{2}} /\left(1-\boldsymbol{x}^{2}\right)$ gives the correct answer. See also [1]

Let $f(z)=g(z) / z^{n}$ where $g$ is analytic and has a zero or order $m$ at the origin; $m, n$ being positive integers and $\mathrm{m}<\boldsymbol{n}$. At $\boldsymbol{z}=0, \boldsymbol{f}$ has a pole of order $\mathrm{n}-\mathrm{m}$ and we discuss its residue R .

Considering, INCORRECTLY, $f$ to he_a_pole_forder $n$, at the origin, we have,

$$
\begin{equation*}
R=\left.\frac{1}{(n-1)!} \frac{d^{n-1}}{d z^{n-1}}\left(z^{n} \frac{g(z)}{z^{n}}\right)\right|_{z=0}=\frac{1}{(n-1)!} g^{(n-1)}(0) \tag{1}
\end{equation*}
$$

However, the pole off at $\boldsymbol{z}=0$ is actually of order $\mathrm{n}-\boldsymbol{m}$, and, so

$$
\begin{equation*}
R=\left.\frac{1}{(n-m-1)!} \frac{d^{n-m-1}}{d z^{n-m-1}}\left(z^{n-m} \frac{g(z)}{z^{n}}\right)\right|_{z=0}=\left.\frac{1}{(n-m-1)!} \frac{d^{n-m-1}}{d z^{n-m-1}}\left(\frac{g(z)}{z^{m}}\right)\right|_{z=0} \tag{2}
\end{equation*}
$$

Now let $\boldsymbol{g}(\boldsymbol{z})=\boldsymbol{z}^{\boldsymbol{m}} \boldsymbol{G}(\boldsymbol{z})$ so that $G$ is analytic at $\boldsymbol{z}=0$ and $G(0) \# 0$. Using Leibnitz's theorem, we have

$$
g^{(n-1)}(z)=z^{m} G^{(n-1)}(z)+\binom{n-1}{1} m z^{m-1} G^{(n-2)}(z)+\ldots+\binom{n-1}{m} m!G^{(n-1-m)}(z)
$$

so that

$$
g^{(n-1)}(0)=\frac{(n-1)!}{(n-m-1)!} G^{(n-1-m)}(0)
$$

reducing (2) to (1).
Incorrectly obtained result (1) can also be seen to be true by using the power series

$$
\sum_{k=0}^{\infty} \frac{z^{k}}{k!} g^{k}(0)
$$

of $g$. However, the above approach illustrates an application of Leibnitz's theorem - generally forgotten or ignored by students - for finding the nth derivative of the product of two functions.

## Reference

1. R . Katz and S. Venit, "Partial Differentiation of Functions of a Single Variable," Pi Mu Epsilon: Journal 7 (1982), 405 -406.

## PROBLEM DEPARTMEN

## Edited by Clayton W. Dodge

University of Maine
This department welcomes problems believed to be new and at a level appropriate for the readers of this journal. Old problems displaying novel and elegant methods of solution are also invited. Proposals should be accompanied by solutions if available and by any information that will assist the editor. An asterisk $\left(^{*}\right)$ preceding a problem number indicates that the proposer did not submit a solution.

All communications should be addressed to C. W. Dodge, Math. Dept., University of Maine, Orono, ME 04469. Please submit each proposal and solution preferably typed or clearly written on a separate sheet (one side only) property identified with name and address. Solutions to problems in this issue should be mailed by December 15, 1992.

## Correction

761. [Fall 1991] Proposed by Murray S. KJamkin, University of Alberta, Edmonton, Alberta, Canada.

Determine all functions $f(\boldsymbol{x})$ such that

$$
f(x)=\sum_{n=0}^{\infty} a_{n} x^{n} \text { and } \frac{1}{f(x)}-\sum_{n=0}^{\infty}(-1)^{n} a_{n} x^{n}
$$

The error was that the exponent on the (-1) was incorrectly given as $\boldsymbol{n}+\mathbf{I}$.

## Problems for Solution

771. Proposed by Alan Wayne, Holiday, Florida.

In the base six addition

$$
E V E+E V E+E V E+A N D=1310
$$

the digits of the addends have been unambiguously replaced by letters. Restore the digits. Where was EVE?
772. Proposed by Robert C. Gebhwdt, Hopatcong, New Jersey.

Let $\boldsymbol{x} \mathbf{x} 44$ be a four-digit number and $\boldsymbol{y} \boldsymbol{y}$ be a two-digit number in base $b>4$. Find $\boldsymbol{x}$ and $y$ in terms of $b$ so that $(y y)^{2}=x \not x 4$ in every such base $b>4$ (such as $88^{2}=7744$ in base ten).
773. Proposed by Leon Bankoff, Los Angeles, California

In a given circle $(\mathrm{O})$ a chord $C D$ is drawn to intersect diameter $A O B$ at point $E$. Three circles are inscribed, the first two in the sectors $B E C$ and $B E D$, and the third in the opposite segment $C E D$. Let the circle in sector $B E C$ touch $C E$ at $\boldsymbol{J}$ and let the circle in sector $B E D$ touch $D E$ at $N$. See the figure. If the three inscribed circles have equal radii,
a) show that $C D$ is perpendicular to $A S$,
b) find the ratio $A E / E B$,
c) find the ratio $A D / A B$,
d) find the ratio $C D / A B$,
e) show that the rectangle $J K M N$ on $\boldsymbol{J N}$ as base and with opposite side KM passing through $\boldsymbol{A}$ circumscribes the third inscribed circle, and
f) show that the rectangles $J K L D$ and $N M W$ are golden rectangles.


Problem 773
774. Proposed by Robert C. Gebhwdt, Hopatcong, New Jersey.

The first player in a game who acquires 250 points is the winner. Because player A is a better player than player B, he gives player B a 50-point handicap. Similarly player B gives player C a $\mathbf{5 0}$-point handicap and player C gives player D a 50-point handicap. What handicap should player A give player D?
775. Proposed by Norman Schaumberger, Bronx Community College, Brons, New York.

If $H$ is the harmonic mean of the positive numbers $a_{1}, a_{2}, \ldots, a_{\mathrm{n}}$, prove that

$$
H^{\frac{1}{a_{1}}+\frac{1}{a_{2}}+\cdots+\frac{1}{a_{n}}} \geq a_{1}^{\frac{1}{a_{1}}} a_{2}^{\frac{1}{a_{2}} \cdots a_{n}^{\frac{1}{a_{n}}}}
$$

776. Proposed by Russell Euler, Northwest Missouri State University, Maryville, Missouri. Let $n$ be a fixed positive integer and let

$$
P_{k}=1^{\mathbf{k}}+2^{\mathbf{k}}+\cdots+n^{\mathbf{k}}
$$

Write as a polynomial in $P_{1}$ the expression

$$
15^{4}\left(P_{1}^{4}+P_{2}^{4}+P_{3}^{4}+P_{4}^{4}\right)
$$

777. Proposed by Seung-Jin Bang, Seoul, Korea

It is well known that $\ln (n+1)<S_{\mathrm{a}}<\operatorname{In} n$, where

$$
S_{n}=1+\frac{1}{2}+\frac{1}{3}+\cdots+\frac{1}{n}
$$

It is also known (OutMathematiconum 11 (1985) p. 109) that

$$
n(n+1)^{1 / n}-n<S_{n}<n-(n-1) n^{-1 /(n-1)}
$$

$\ln (n+1)<n(n+1)^{1 / n}-n$ and $n-(n-1) n^{-1 /(n-1)}<1+I n n$ for all $n$ a 2
778. Proposed by Laura L. Kelleher and Frank P. Battles, Massachusetts Maritime Academy, Buzzards Bay, Massachusetts.

It is readily established that the arc length along the curve $\boldsymbol{y}=\boldsymbol{\operatorname { c o s h }} \boldsymbol{x}$ on any interval $[\boldsymbol{a}, \boldsymbol{b}]$ and the area under the graph of this same function on this same interval are numerically equal. For what other functions, if any, is this curious fact true?

## 779. Proposed by W. Moser, McGill University, Montreal, Canada,

If $0<a \leq x \leq y \leq 1 / a$, then prove that

$$
\begin{gathered}
x+\frac{1}{x} \leq a+\frac{1}{a}, \quad \frac{x}{y}+\frac{y}{x} \leq \frac{y}{a}+\frac{a}{y} \\
\frac{x}{y}+\frac{y}{x} \leq a x+\frac{1}{a x}, \quad \text { and } \\
(x+y)\left(\frac{1}{x}+\frac{1}{y}\right) \leq\left(a+\frac{1}{a}\right)^{2} .
\end{gathered}
$$

780. Proposed by R. S. Luthar, University of Wisconsin Center, Janeswille, Wisconsin.

Let $A B C D$ be a parallelogram with $L A=60^{\circ}$. Let the circle through $\boldsymbol{A}, \boldsymbol{B}$, and $\boldsymbol{D}$ intersect $\boldsymbol{A C}$ at $\boldsymbol{E}$. See the figure. Prove that $\boldsymbol{B D ^ { 2 }}+\boldsymbol{A B} \cdot A D=A E \cdot E C$.


Problem 780

Problem 781
781. Proposed by the late Jack Garfunkel, Flushing, New York.

Erect squares $A D E F, B D K L$, and $C D G H$ as shown in the figure, on the segments $A D, D C$, and $\boldsymbol{B D}$, where $\boldsymbol{D}$ is any point on side $C A$ of given triangle $\boldsymbol{A B C}$. Let $A ; Y$, and $Z$ be the centers of the erected squares. Prove that triangles $A B C$ and $X Y Z$ are similar and the ratio of similarity is $\sqrt{ }$.
782. Proposed by Murray S. Klamkin, University of Alberta, Edmonton, Alberta, Canada.

In O. Bottema et al, Geometric Inequalities, Wolters-Noordhoff, Gronigen, 1969, item 12.55, p. 118 , it is stated that for a triangle $A B C$ with no angle $\geq 2 \pi / 3$ that

$$
2\left(R_{1}+\mathrm{R},+R_{3}\right)^{2} \geq\left(\mathrm{a}^{2}+\mathrm{b}^{2}+c^{2}\right)+4 F \sqrt{3}
$$

where $\boldsymbol{R}_{1}, \boldsymbol{R}_{\mathbf{2}}$ and $\boldsymbol{R}_{\mathbf{3}}$ are the respective distances from an arbitrary point $\boldsymbol{P}$ inside the triangle to its vertices, $\boldsymbol{a}$, b, and $\boldsymbol{c}$ are the triangle's side lengths, and $\boldsymbol{F}$ is its area. Item 12.55 further states that for a triangle in which $L A \geq 2 \pi / 3$,

$$
\left(R_{1}+R_{2}+R_{3}\right)^{2} \geq(\mathrm{b}+\mathrm{c})^{2}
$$

Show that the first inequality is true for all triangles.
783. Proposed by the late Jack Garfunkel, Flushing, New York.

If, $\boldsymbol{A}, \boldsymbol{B}$, and $\boldsymbol{C}$ are the angles of a triangle $A B C$, then prove that

$$
\frac{\sum \sin ^{2} A}{\sum \cos ^{2}\left(\frac{A}{2}\right)} \geq \frac{\Pi \sin A}{\Pi \cos \left(\frac{A}{2}\right)} .
$$

## Solutions

403. [Fall 1977, Fall 1983, Fall 1984]. Proposed by David L. Silverman, Lest Los Angeles, California

Two players play a game of 'Take It or Leave It" on the unit interval $\mathbf{( 0 , 1 )}$. Each player privately generates a random number from the uniform distribution and either keeps it as his score or rejects it and generates a second number which becomes his score. Neither player knows, prior to his own play, what his opponent's score is or whether it is the result of an acceptance or a rejection. (However, variants based on modifying this condition, either unilaterally or bilaterally, are interesting.)

The scores are compared and the player with the higher score wins $\$ 1.00$ from the other.
a. What strategy will give a player the highest expected score?
b. What strategy will give a player the best chance of winning?
c. If one player knows that his opponent is playing so as to maximize his score, how much of an advantage will he have if he employs the best counter-strategy?
II. Comment and solution by Peter Griffin, California State University Sacramento, California

Unfortunately the published solution [Fall 1984] is incorrect in two major particulars. Part (b) is coincidentally correct because the expectation for $\boldsymbol{m}=(\sqrt{\mathbf{5}-1)} / 2$ turns out to be itself. This is not a valid argument for establishing that this strategy will beat any other strategy, however. The answer to part (c) is wrong because it is based on the same fallacious reasoning as (b).

Both parts (b) and (c) have nothing to do with expectation, but involve how often one strategy will do better than another. Here is a true historical analogy: In the 1960 World Series the New York Yankees averaged 8 runs to only 4 runs for the Pittsburgh Pirates, yet they lost the World Series. The Pirates won more often (4 games to 3), but when the Yankees won, they tended to win by many, many more runs. I will sketch briefly how to solve parts (b) and (c).
"Using an a-rule" shall mean keeping your first random number if it exceeds the criterion: number $a$, otherwise discarding it and being left with the second random number regardless of what it
is. The distribution rule of one's final score using an a-rule is $\boldsymbol{F} \boldsymbol{F}(\boldsymbol{x})=\boldsymbol{\operatorname { P r }}[\boldsymbol{A} \mathrm{Sx}]$. This fundion has two formulas:

$$
\begin{gathered}
\text { for } x \leq \mathrm{a}, F_{\mathrm{a}}(x)=\operatorname{Pr}\left[R_{1} \leq \mathrm{a}\right] \cdot \operatorname{Pr}\left[R_{2} \leq \mathrm{a}\right]=o r \text {, while } \\
\text { for a }<x, F_{2}(x)=\operatorname{Pr}\left[a<R_{1} \leq x\right]+\operatorname{Pr}\left[R_{1} \leq a\right] \cdot \operatorname{Pr}\left[R_{2} \leq \mathrm{a}\right] \\
=(x-a)+a x=(a+1) x-a .
\end{gathered}
$$

From these formulas we derive the density function for an a-rule to be $f_{\mathbf{a}}(\boldsymbol{x})=\boldsymbol{F}_{\mathbf{a}}{ }^{\prime}(\boldsymbol{x})$, so

$$
f_{n}(x)=\mathrm{a} \text { if } x \leq \mathrm{a}, \quad \text { and } \quad f_{\mathrm{a}}(x)=\mathrm{a}+\mathbf{1} \text { if } \boldsymbol{x}>\mathrm{a} .
$$

The density is integrated over various intervals to End the probability of being in the interval. To find the probability that the random score $\boldsymbol{A}$ is approximately equal to $\boldsymbol{x}$, we use $f_{\mathbf{A}}(\boldsymbol{x}) d x$.

The probability that an a-rule beats a $\boldsymbol{b}$-rule, assuming a $\leq \boldsymbol{b}$, is found by approximating
$\Sigma \operatorname{Pr}[A=x] \cdot \operatorname{Pr}[B<x]$ with $\Sigma F_{b}(x) \cdot f_{t}(x) d x$, which gives $\operatorname{Pr}[A>B]=\Sigma F_{b}(x) f_{t}(x) \Delta x$, and hence

$$
\int_{0}^{1} F_{b}(x) f_{a}(x) d x
$$

$$
=\int_{0}^{a} b x \cdot a d x+\int_{a}^{b} b x \cdot(a+1) d x+\int_{b}^{1}[(b+1) x-b](a+1) d x
$$

$$
=\frac{1}{2}+\frac{(a-b)(1-a b-b)}{2}=P(a, b) .
$$

Note that $\boldsymbol{P}(\boldsymbol{a}, \boldsymbol{b})$ is a quadraticin each variable separately. First, fix b and maximize $P(\boldsymbol{a}, \boldsymbol{b})$ as a function of a, yielding

$$
\frac{\partial p}{\partial a}=\frac{1-b-2 a b+b^{2}}{2}=0 \text { implies } \hat{a}=\frac{b^{2}-b+1}{2 b} .
$$

For $b=1$, this gives $6=1 / 2$. But for $b=1 / 2, a=3 / 4$, which is not admissible, being greater than $1 / 2$ Admiisibiility of 6 requires that $\left(b^{2}+b+1\right) / 2 b \leq b$, which implies that $\boldsymbol{b}^{2}+b-1 a 0$, and hence that $\mathrm{b} \geq m=\left(\mathrm{V} 5^{-1}\right) / \mathbf{2}$.

Next fix a and maximize $\boldsymbol{P}(\boldsymbol{a}, \boldsymbol{b})$, the probability that B loses to $\boldsymbol{A}$. This requires that

$$
\frac{\partial p}{\partial b}=\frac{-a^{2}-a-1+2 a b+2 b}{2}=0 \text { implies } \hat{b}=\frac{a^{2}+a+1}{2(a+1)}
$$

Admissibility of $\hat{b}$ follows from $\hat{b} 2 \mathrm{a}$, implying $\left(\mathrm{a}^{2}+\mathrm{a}+1\right) /(2(a+1)) 2 \mathrm{a}$, so $\mathrm{a}^{2}+\mathrm{a}-1 \mathrm{~s} 0$, and finally a $\leq \mathrm{m}=(\mathrm{V} 5 \cdot \mathbf{1}) / \mathbf{2}$. For $\mathrm{a}=\mathbf{1 / 2}$, we get $\hat{b}=7 / 12$ and $P(1 / 2,7 / 12)=95 / 192$, so the 7112 -rule beats the $1 / 2$-rule $1-95 / 192=97 / 19$ of the time. This is better than

$$
1-P\left(\frac{1}{2}, \frac{5}{8}\right)=1-\left(\frac{1}{2}+\frac{-\frac{1}{8}\left(1-\frac{5}{16}-\frac{5}{8}\right)}{2}\right)=\frac{1}{2}+\frac{1}{256} .
$$

It is no surprise that the published solution mentioned simulations that gave $50.4 \%$ wins, since $\mathbf{1 / 2}+$ $\mathbf{1} / 256=0.504$. The 7112 -rule gives $1 / 2+\mathbf{1 / 1 9}=0.505$. So, $7 / 12$ is the answer to Part (c), not $5 / 8$.

Because $\boldsymbol{m}$ satisfies the relation $\boldsymbol{m}=\mathbf{1} /(\mathbf{1}+\mathrm{m})$ and $\mathbf{1}=\boldsymbol{m}+\boldsymbol{m}^{2}$, it is not hard to show that
$P(m, b)>1 / 2$ for all $\mathrm{b}>m$ and $P(b, m)<1 / 2$ for all $\mathrm{b}<\boldsymbol{m}$, which demonstratesthat $\boldsymbol{m}$ will beat any rule other than itself, which it ties. It provides a saddle point of the function $P(a, b)$, namely $(m, m)$.

The following BASIC program simulates the $7 / 12$ strategy versus the $1 / 2$ strategy $\mathbf{1 0 , 0 0 0 , 0 0 0}$ times, printing the results every 10,000 games. To simulate the $5 / 8$ strategy, replace each 7 by 5 and each 12 by 8 in line 12 . My results of running this program were that $7 / 12$ beats $\mathbf{1 / 2}$ with frequency 0.505215 and $\mathbf{5} / 8$ beats $\mathbf{1 / 2}$ with frequency 0.503897 , which compare with the ideals $97 / 19=0.505218$ and $\mathbf{1 2 9} / \mathbf{2 5 6}=0.503906$.

## 2 FOR J = 1 TO 1000 <br> 3 RANDOMIZE

4 FOR I = 1 TO 10000
$6 \mathrm{~A}=\mathrm{RND}$
8 IF A $<.5$ THEN A $=2 *$ A
$10 \mathrm{~B}=$ RND
12 IF B < 7/12 THEN B $=12^{*} \mathrm{~B} / 7$
14 IF A > B THEN 18
$16 \mathrm{~W}=\mathrm{W}+1$
18 NEXT I
20 PRINT J, 10000*W/J
22 NEXT J
731. [Spring 1990, Spring 1991] Proposed by Roger Pinkham, Stevens Institute of Technology, Hoboken, New Jersey.
a) Show that on the lattice points in the plane having integer coordinatesone cannot have the vertices of an equilateral triangle.
*b) What about a tetrahedron in 3 -space?
VI. Comment by Seung-Jin Bang Seoul, Korea.

In The Newsletter of the Korean MathematicalSociety, No. 27 (July 1991) p. 17, there appears a solution by a colleague and me to the generalization of part (a) that states that no regular $(2 n+1)$ gon can have all rational coordinatesin the Euclidean plane. We furthermore point out that the result is true for any regular $\boldsymbol{n}$-gon provided that $\boldsymbol{n}$ is not a power of $\boldsymbol{2}$.
733. [Fall 1990,Fall 1991] Proposed by Roger Pinkham, Stevens Institute of Technology, Hoboken, New Jersey,

If $p(x)$ is a polynomial and $p(x) \geq 0$ for all $x$, then

$$
p+p^{\prime}+p^{\prime \prime}+\ldots \geq 0
$$

for allx.
IV. Solution by David Yavenditti, Alma, Michigan.

Let $S(x)=p(x)+p^{\prime}(x)+p^{\prime \prime}(x)+\cdots$, so that $S^{\prime}(x)=p^{\prime}(x)+p^{n}(x)+p^{m}(x)+\cdots$. Then $\boldsymbol{S}(\boldsymbol{x})=\boldsymbol{p}(\boldsymbol{x})+\boldsymbol{S}^{\prime}(\boldsymbol{x})$. S i c e p is a polynomial that is always nonnegative, thenp attains a minimum. Now $\boldsymbol{S}$ is a polynomial with the same leading term as p, so $\boldsymbol{S}$ also must have a minimum. Since $\boldsymbol{S}$ is a polynomial, then $\boldsymbol{S}$ attains its minimum value at $\boldsymbol{x}=c$ only if $\boldsymbol{S}^{\prime}(\boldsymbol{c})=0$. Then, for all real $\boldsymbol{x}$,

$$
S(x) \geq \min \{S(x)\}=S(c)=p(c)+S^{\prime}(c)=p(c) \geq 0
$$

745. [Spring 1991] Proposed by Alan Wayne, Holiday, Florida.

Find all solutions to

$$
\begin{array}{r}
E N I D \\
+\quad D I D \\
\hline D I N E
\end{array}
$$

I. Solution by Victor G. Feser, University of Mary, Bismarck, North Dakota.

Since there are four symbols, we solve the problem in base $B$, where $B \geq 4$. From the $\boldsymbol{B}^{3}$ column, $\mathrm{D}=E+1$. From the units column $\mathrm{W} \rightarrow E$, where the arrow is read "yields" and is equivalent to "congruent mod $B$." It signifies that 1 may be carried into the next column. Then $2 \mathrm{E}+2 \rightarrow E$, so $E+2 \rightarrow 0$, whence $E=B-2$. Now $D=B=1$ and 1 is carried to the $B$ column (the tens column in base ten).

If 1 is carried to the $B^{2}$ column, then it becomes $\mathbf{1}+N+(B \cdot I) \rightarrow I$, and we have $N=I$, which is not allowed. So the $B$ and $\boldsymbol{B}^{2}$ columns must read $21+1=N$ and $N+(B \cdot 1)=I+B$. Thus $N=I+I$ and, from the $B$ column, $I=0$ and $N=1$. Hence, for each base $B \geq 4$, the unique solution is


To complete the problem we show that there is no solution for any negative base $B \leq-4$. Since successive powers of a negative number alternate signs, if we carry $l$ from a column in an addition, it carries into the next column as -1 . From the $B^{3}$ column we get $\mathrm{E}^{-1} 1=D$, and from the units column, $\mathrm{W} \rightarrow E$. Substituting, as before, we get $2 \mathrm{E}=2 \rightarrow E$, so $E=2$ and $D=1$. Now $21 \rightarrow N$, that is, $21=$ $N+\boldsymbol{c}|\boldsymbol{B}|$, where $c=0$ or 1 . Also $-\boldsymbol{c}+N+\mathbf{1}=I+|\boldsymbol{B}|$, which demands that $N \geq|\boldsymbol{B}|$, an impossibility. So there is no solution for any negative base.
II. Comment by the Proposer.

Easy, wasn't it?
Also solved for any positive base by CHARLES ASHBACHER, Hiawatha, LA, SEUNG-JIN BANG, Seoul, Korea, WILLIAM CHAU, New York, $N Y$, HENRY S. LIEBERMAN, Waban, MA, BOB PRIELIPP,Universityof Wisconsin-Oshkosh, KENNETHM. WILKE, Topeka, $\boldsymbol{K S}$, and the PROPOSER. Base ten solutions were submitted by JOHN T. ANNULIS, University of Arkansas-Monticello, FRANK P. BATTLES, Massachusetts Maritime Academy, Buzzards Bay, JAMES E. CAMPBELL, Indiana University at Bloomington, CAVELAND MATH GROUP, Western Kentucky University, Bowling Green, MARK EVANS, Louisville, KY, HOWARD FORMAN, Parsippany, NJ, DAWN M. GALAYDA, St. Bonaventure University, NY, ROBERT C. GEBHARDT, Hopatcong, NJ, S. GENDLER, Clarion

University of Pennsylvania, RICHARDI. HESS, Rancho Palos Verdes, CA, NATHAN JASPEN, Stevens Institute of Technology, Hoboken, NJ, LOWELL F. LYNDE, JR., University of Arkansas at Monticello, WADE H. SHERARD, Furnan University, Greenville, SC, REX H. WU, New York, NY, and DAVID YAVENDITTI, Alma, MI.
746. [Spring 1991] Proposed by Gregory Wulczyn, Bucknell University, Lewisburg Pennsylvania,

Find the least positive integer $\boldsymbol{n}$ that will have remainder 1 when divided by $r$, the quotient will have remainder 2 when divided by $r$, the new quotient will have remainder 3 when divided by $\boldsymbol{r}$, and so forth through $\boldsymbol{r}-1$ divisions. That is, $\boldsymbol{n}=\boldsymbol{q}_{0}$, and $\boldsymbol{q}_{\mathrm{k}-1}=\boldsymbol{q}_{\boldsymbol{r}} \boldsymbol{r}+\mathrm{k}$ for $\mathrm{k}=1,2, \ldots, \boldsymbol{r}-1, r$ a positive integer greater than 1 .
I. Solution by John Putu, Alma College, Alma, Michigan.

We have $\boldsymbol{q}_{0}=1+\boldsymbol{q}_{1} r_{,} \boldsymbol{q}_{1}=2+\boldsymbol{q}_{\boldsymbol{r}}, \ldots, \boldsymbol{q}_{\boldsymbol{r}-2}=\boldsymbol{r}-\mathbf{1}+\boldsymbol{q}_{r-1} \boldsymbol{r}$. Multiply the second equation by $r$, the third by $r^{2}$, and so forth, and then substitute to get

$$
n=q_{0}=1+2 r+3 r^{2}+\ldots+(r-1) r^{r-2}+q_{r-1} r^{r-1}
$$

To minimize $\boldsymbol{n}$, we choose $\boldsymbol{q}_{\boldsymbol{r}-1}=0$ since $\boldsymbol{r}>1$. So

$$
n=1+2 r+3 r^{2}+\ldots+(r-1) r^{-2}
$$

and

$$
m=r+2 r^{2}+3 r^{3}+\ldots+(r-1) r^{r-1}
$$

Subtracting, we have

$$
(1-r) n=1+r+r^{2}+\ldots+r^{r-2}-(r-1) r^{r-1}=\frac{1-r^{r-1}}{1-r}+(1-r) r^{r-1}
$$

so

$$
n=r^{r-1}-\frac{r^{r-1}-1}{(r-1)^{2}}
$$

II. Solution by Stephen I. Gendler, Clarion University of Pennsylvania, Clarion, Pennsylvania The numbers described seem to be no more than

$$
(r-1)(r-2)(r-3) \ldots(3)(2)(1) \text { in baser, }
$$

where each pair of parentheses is a digit, since the repeated division is a method for changing bases. Such numbers appear in any positive base $\boldsymbol{r}$, the smallest being $\boldsymbol{r}=2, \boldsymbol{n}=1$. A few larger examples are listed below.

| $n$ (base $r$ ) | $n$ (base ten) |
| :---: | :---: |
| 21 | 7 |
| 321 | 57 |
| 4321 | 586 |
| 54321 | 7456 |
| 654321 | 114381 |
| 7654321 | 2054353 |
| 87654321 | 42374116 |
| 987654321 | 987654321 |

## HI. Comment by the Proposer.

This problem is a variant of the ordinary cocoanut-monkey problem in which the cocoanuts in a given pile are to be divided equally among $\boldsymbol{r}$ people the next morning. During the night one of the people sneaks out to the pile and divides it into $r$ equal piles with exactly $s$ cocoanuts left over, where is $\boldsymbol{s}<\boldsymbol{r}$. The $\boldsymbol{s}$ cocoanuts are thrown to a waiting monkey and the person hides one of the equal piles as his/her share. The remaining cocoanuts are repiled into one pile. As the night progresses, each of the $\boldsymbol{r}$ people in turn sneaks out to the pile and repeats the procedure of dividing the cocoanuts into $\boldsymbol{r}$ equal piles with exactlys cocoanuts left over, throwing the $s$ cocoanuts to the monkey, and hiding one pile as that person's share. In the morning, there remain just enough cocoanuts to be divided equally among the $\boldsymbol{r}$ people. Heres cocoanuts plus $1 / r$ of the remaining pile are removed exactly $\boldsymbol{r}$ times to leave a multiple of $\boldsymbol{r}$ cocoanuts in the pile. See problem 3242 in The American Mathematical Monthly (January 1928).

Also solved by CHARLES ASHBACHER, Hiawatha, LA, SEUNG-JIN BANG, Seoul, Korea, JAMES E. CAMPBELL, Indiana University at Bloomington, CAVELAND MATH GROUP, Western Kentucky University. Bowling Green, WILLIAM CHAU, New York, NY, CHARLES R. DIMINNIE, St. Bonaventure University, NY, MARK EVANS, Louisville, KY, VICTOR G. FESER, University of Mary, Bismarck, ND, HOWARD FORMAN, Parsippany, NJ, RICHARD I. HESS, Rancho Palos Verdes, CA, HENRY S. LIEBERMAN, Waban, MA, MOHAMMAD P. SHAIKH, Western Michigan University, Kalamazoo, KENNETH M. WILKE, Topeka, KS, REX H. WU, New York, NY, DAVID YAVENDITTI, Alma, MI, and the PROPOSER.
747. [Spring 1991] Proposed by the late Jack Garfunkel, Flushing, New York.

Let $\boldsymbol{A B C}$ be a triangle with inscribed circle $(\boldsymbol{I})$ and let the line segmentsAI, $\boldsymbol{B I}$, and $\boldsymbol{C I}$ cut the incircle at $\boldsymbol{A}^{\prime}, \boldsymbol{B}^{\prime}$, and $\boldsymbol{C}^{\prime}$ respectively. Prove that

$$
\sin A^{\prime}+\sin B^{\prime}+\sin C^{\prime} \geq \cos \frac{A}{2}+\cos \frac{B}{2}+\cos \frac{C}{2}
$$

where $\boldsymbol{A}^{\prime}, \boldsymbol{B}^{\prime}$, and $\boldsymbol{C}^{\prime}$ are the angles of triangle $\boldsymbol{A}^{\prime} \boldsymbol{B}^{\prime} \boldsymbol{C}^{\prime}$.
Solution by William Chau, New York, New York.
Since $\mathbf{2} \angle A^{\prime}=L B I C=\boldsymbol{\pi}^{-}(\angle B / 2+\angle C / 2)$, it follows that

$$
\sin 2 A^{\prime}=\sin \frac{B+C}{2}=\sin \frac{\pi-A}{2}=\cos \frac{A}{2},
$$

with similar equalities for $\boldsymbol{B}$ and for C . Now the stated equation is equivalent to

$$
\sin \boldsymbol{A}^{\prime}+\sin \boldsymbol{B}^{\prime}+\sin \boldsymbol{C}^{\prime} \geq \sin 2 \boldsymbol{A}^{\prime}+\sin 2 B^{\prime}+\sin 2 C^{\prime}
$$

which is item 2.4 on p. 18 of O. Bottema et al, Geometric Inequalities, Wolters-Noordhoff, Gronigen, 1968.


Also solved by SCOTT H. BROWN, Stuart Middle School, FL, RUSSELLEULER, Northwest Missouri State University, Maryville, MURRAY S. KLAMKIN, University of Alberta, Canada. YOSHINOBU MURAYOSHI, Eugene, OR, BOB PRIELIPP, University of Wisconsin-Oshkosh, REX H. WU, New Yark, NY, and the PROPOSER, Prielipp pointed out that this same problem appeared with his solution as problem 4274 in the March 1991 School Science and Mathematics.
748. [Spring 1991] Proposed by the late John Howell, Littlerock, California.
a) An urn contains $\boldsymbol{n}$ balls numbered $\boldsymbol{I}$ to $\boldsymbol{n}$. Algernon, Beauregard, and Chauncey draw a ball one after another with replacement. The game is terminated when two consecutive drawings produce the same ball. Find the probabilities of terminating on Algernon's draw, on Beauregard's draw, and on Chauncey's draw.
b) Repeat the problem for the case that the game terminates when three consecutive drawings produce the same ball.

Amalgam of solutions by David Yavenditti, Alma, Michigan, and Morris Katr, Macwahoc, Maine. Instead, we generalize the result fork consecutive drawings of the same ball, $\mathrm{k}>1$. We denote by $P(X)$ the probability the game terminates on $X^{\prime}$ s draw. Let the players be denoted by $X_{\mathrm{b}}$, where $X_{1}=$ Algernon, Beauregard, or Chauncey according as $i \equiv 1,2$, or $3(\bmod 3)$. The first person who could $\boldsymbol{\operatorname { w i n }}$ (terminate the game) is $\boldsymbol{X}_{\mathbf{k}}$ on the kth play with probabilityp $=\mathbf{1} / \boldsymbol{n}^{\boldsymbol{k}-1}$. Let $\mathrm{q}=1-\boldsymbol{p}$. If $\boldsymbol{X}_{\mathbf{k}}$ does not win at that turn, then $X_{\mathbf{k}+1}$ could win on the $(\boldsymbol{k}+1)$ st play with probability $\boldsymbol{p q}$, or $X_{\mathbf{k}+2}$ could win on the $(\boldsymbol{k}+2)$ nd play with probability $\boldsymbol{p} \boldsymbol{q}^{2}$, etc. Hence

$$
\begin{gathered}
P\left(X_{k}\right)=p+p q^{3}+p q^{6}+p q^{9}+\ldots \\
=\frac{p}{1-q^{3}}=\frac{1 / n^{k-1}}{1-\left(1-1 / n^{k-1}\right)^{3}}=\frac{n^{2 k-2}}{3 n^{2 k-2}-3 n^{k-1}+1}
\end{gathered}
$$

Now, let $\boldsymbol{X}_{\mathbf{k}}$ lose on the kth play (with probability q). Then $\boldsymbol{X}_{\mathbf{k}+\boldsymbol{1}}$ faces the same conditions that $\boldsymbol{X}_{\mathbf{k}}$ did on the $k$ th play, and we have

$$
P\left(X_{k+1}\right)=q \cdot P\left(X_{k}\right)=\frac{n^{2 k-2}\left(1-1 / n^{k-1}\right)}{3 n^{2 k-2}-3 n^{k-1}+1}=\frac{n^{2 k-2}-n^{k-1}}{3 n^{2 k-2}-3 n^{k-2}+1} .
$$

## Similarly,

$P\left(X_{k+2}\right)=q^{2} \cdot P\left(X_{k}\right)=\frac{n^{2 k-2}\left(1-1 / n^{k-1}\right)^{2}}{3 n^{2 k-2}-3 n^{k-1}+1}=\frac{n^{2 k-2}-2 n^{k-1}+1}{3 n^{2 k-2}-3 n^{k-1}+1}$.
It follows that, for part (a),

$$
P(B)=\frac{n^{2}}{3 n^{2}-3 n+1}, P(C)=\frac{n^{2}-n}{3 n^{2}-3 n+1}, P(A)=\frac{n^{2}-2 n+1}{3 n^{2}-3 n+1}
$$

Similarly, for part (b) we have that

$$
P(C)=\frac{n^{4}}{3 n^{4}-3 n^{2}+1}, P(A)=\frac{n^{4}-n^{2}}{3 n^{4}-3 n^{2}+1}, P(B)=\frac{n^{4}-2 n^{2}+1}{3 n^{4}-3 n^{2}+1}
$$

Also solved by CHARLES ASHBACHER, (Part (a) only), Hiawatha, LA, JAMES E. CAMPBELL,(Part (a)only), IndianaUniversityat Bloomington, WILLIAMCHAU,(Part (a) only). New York, NY, MARK EVANS, Louisville, KY, HOWARD FORMAN, (Part (a) only), Parsippany, NJ, RICHARD I. HESS, Rancho Palos Verdes, CA, MOHAMMAD P. SHAIKH, Western Michigan University, Kalamazoo, REX H. WU, (Part (a) only), New York, $\mathbf{N Y}$, and the PROPOSER. Not all solutions agreed with that of Morris Kate.
749. [Spring 1991] Proposer by R. S. Luthar, University of Wisconsin Center at Janesville, Janesville, Wisconsin.

$$
\text { If } \sin x+\sin y+\sin z=0 \text {, then prove that }
$$

$$
|\sin 3 x+\sin 3 y+\sin 3 z| \leq 12|x y z| .
$$

Solution by Bob Prielipp, University of Wisconsin-Oshkosh, Oshkosh,
Wisconsin.

$$
\begin{aligned}
& \text { If } a+b+\boldsymbol{c}=\mathbf{0} \text {, then } \\
& \qquad \mathbf{a}^{3}+\boldsymbol{b}^{3}+c^{3} \cdot \mathbf{3 a b c}=(a+b+c)\left(a^{2}+\boldsymbol{b}^{2}+c^{2} \cdot b c \cdot c a=a b\right)=0 .
\end{aligned}
$$

Thus, $\operatorname{since} \sin x+\sin y+\sin z=0$, we have

$$
\sin ^{3} x+\sin ^{3} y+\sin ^{3} z=3 \sin x \sin y \sin z
$$

Also, since $\sin 3 t=3 \sin t-4 \sin$ t, we get that

$$
\sin 3 x+\sin 3 y+\sin 3 z=-4\left(\sin ^{3} x+\sin ^{3} y+\sin ^{3} z\right)=-12 \sin x \sin y \sin z .
$$

Hence the required inequality holds if and only if

## $|\sin x \sin y \sin z| \leq|x y z|$,

which follows immediately from the vell known inequality $|\sin t| \leq|t|$.
Also solved by SEUNG-JIN BANG, Seoul, Korea, MARTIN BAZANT, Tucson, AZ, WILLIAM CHAU, New York, NY, RUSSELL EULER, Northwest Missouri State University. Maryville, MURRAY S. KLAMKIN, University of Alberta, Canada, YOSHINOBU MURAYOSHI, Eugene, OR, MOHAMMADP. SHAIKH, Western Michigan University. Kalamazoo, REX H. WU, New York, NY, and the PROPOSER.
*750. [Spring 1991] Proposed by Dmitry P. Mavlo, Moscow, U.S.S.R.
Solve the system of equations

$$
2^{x} y+\left(3^{x}\right) \sqrt{1-y^{2}}=\sqrt{3} \quad \text { and } 3^{x} y-\left(2^{x}\right) \sqrt{1-y^{2}}=\sqrt{2} .
$$

This problem appeared in the SYMP-86 Entrance Exam Mathematical Problems.
Solution by Parush Saxena, Massachusetts Maritime Academy, Buzzards Bay, Massachusetts. Lettingy $=\sin 8$, we must solve the equations

$$
2^{x} \sin \theta+3^{x} \cos \theta=\sqrt{3} \text { and } 3^{x} \sin 0 \cdot 2^{x} \cos 0=\sqrt{2} \text {. }
$$

Now square these equations and add to get $2^{2 x}+3^{2 x}=5$, which has the unique solution $x=1 / 2$ since the left side of the equation is an increasing function of x . Substituting $\mathrm{x}=1 / 2$ into the first of the original equations yields

$$
\sqrt{2} y+\sqrt{3} \sqrt{1-y^{2}}=\sqrt{3}
$$

which, upon squaring and simplifying, reduces to

$$
y(5 y-2 \sqrt{6})=0, \text { s o } y=0 \text { or } y=\frac{2 \sqrt{6}}{5}
$$

Only the latter value checks in the given equations, so the unique solution is $x=1 / 2$ and $y=2 \sqrt{2} / 5$.
Also solved by CHARLES ASHBACHER, Hiawatha, LA, SEUNG-JIN BANG, Seoul, Korea. FRANK P. BATTLES, Massachusetts Maritime Academy, Buzzards Bay, BARRY BRUNSON, Western Kentucky University, Bowling Green, JAMES E. CAMPBELL, Indiana University at Bloomington, WIITIAM CHAU, New York, NY, CHARLES R. DIMINNIE, St. Bonavenure University, NY, ROBERT O. DOWNES, Long Beach, CA, MARK EVANS, Louisville, KY, HOWARD FORMAN, Parsippany, NJ, ROBERTC. GEBHARDT, Hopatcong NI, RICHARD I. HESS, RanchoPalos Verdes, CA, HENRYS. LIEBERMAN, Waban, MA, PETER A. LINDSTROM, North Lake College, Iving, TX, G. MAVRIGIAN, Youngstown State University. OH, YOSHINOBU MURAYOSHI, Eugene, OR. WILIIAM $H$. PEIRCE, Stonington, CT, GEORGE $W$. RAINEY, Los Angeles, CA, MOHAMMADP. SHAIKH, Western Michigan University, Kalamazoo, DAVE SMITH, Messiah College, Grantham, PA, REX H. WU, New York, $N$ Y, and DAVID YAVENDITTI, Alma, MI. One faulty solution was also' submitted,
751. [Spring 1991] Proposed by Murray S. K]amkin, University of Alberta, Edmonton, Alberta,

## Canada.

Determine all pairs of positive numbers $x$ and $y$ such that

$$
9(x+y)+\frac{1}{x}+\frac{1}{y} \geq 10+\frac{x}{y}+\frac{y}{x} .
$$

## I. Solution by Seung-Jin Bane, Seoul, Republic of Korea

Multiply the stated inequality by the positive quantity $x y$ and rearrange the result to get the equivalent inequality
(1)

$$
(9 y-1) x^{2}+(9 y-1)(y-1) x+y(1-y) \geq 0
$$

Case 1. The left side is a quadratic polynomial whose discriminant is

$$
D=(9 y-1)(y-1)(3 y-1)^{2}
$$

soD $\leq 0$ for $1 / 9 \leq y \leq 1$, whence the original inequality holds for all $x>0$ when $1 / 9 \leq y \leq 1$. By the symmetry of the original inequality, it holds also for ally $>0$ when $1 / 9 \leq x \leq 1$.

Case 2. Since Inequality (1) can be rewritten in the form

$$
(1-9 y) x(1-x)+(1-9 x) y(1-y)+8 x y \geq 0
$$

## we see it is true for $0<x<1 / 9$ and $0<y<1 / 9$, and also for $x>1$ and $y>1$.

Case 3 . We need only consider the region $0<x<1 / 9$ and $y>1$ and by symmetry the region $0<y<1 / 9$ andx $>1$. Apply the quadratic formula to the quadratic polynomial of Inequality (1) to get that, when $\mathrm{y}>1$, we must have

$$
x>\frac{-(9 y-1)(y-1)+(3 y-1) \sqrt{(9 y-1)(y-1)}}{2(9 y-1)}
$$

Now interchange x and y to get the corresponding inequality for the region where $\mathrm{x}>1$.


II, Solution by Robert C, Gebhardt, Hopatcong, New Jersey.
Multiply the inequality by $x y$ and rearrange to get

$$
9 x^{2} y+9 x y^{2}+x+y-x^{2}-y^{2}-10 x y \geq 0 .
$$

We graph the equation, as shown in the accompanying figure. The curve has intercepts (0,0), (1,0), and $\mathbf{( 0 , 1 )}$. Then any address in the first quadrant between the curves (e.g. $\mathbf{( 1 , 1 )}$ or ( $\mathbf{0 . 3 2 1 , 6 5 . 4 3 2 \text { ) ) will satisfy }}$ the given inequality.
III. Comment by Elizabeth Andy, Limerick, Maine.

The graph in the accompanying figure clearly shows asymptotes $\mathrm{x}=\mathbf{1 / 9}, \mathrm{y}=\mathbf{1 / 9}$, andx +y $=8 / 9$. Replace the original inequality by equality and rewrite it in the form

$$
9(x+y)=10+\frac{x-1}{y}+\frac{y-1}{x}
$$

and firally

$$
9(x+y)=8+\frac{(x+y)(x+y-1)}{x y}
$$

We see that, when $x=-y+1$ and $|x|$ is large, the fraction is approximately zero, and we get the equation of the oblique asymptotex $+y=8 / 9$. Similarly, the equation can be rewritten in the form

$$
9 y=1+\frac{(9 y-1)(1-x)}{y}+\frac{y-1}{x}
$$

If $|x|$ is large and $\mathrm{y}=\mathbf{1 / 9}$, then both fractions on the right are approximately zero, and the equation reduces to the horizontal asymptote $9 y=1$. Thus we get the asymptotes algebraically, too.

Furthermore the original inequality shows there is symmetry in $x$ and $y$, that is, in the line $\boldsymbol{y}=$ $x$. Therefore, in Inequality (1) of Solution I above, replace the inequality by equality, interchangex and y , and apply the quadratic formula to get that

$$
y=\frac{-(9 x-1)(x-1) \pm \sqrt{(9 x-1)(x-1)(3 x-1)^{2}}}{2(9 x-1)}
$$

which reduces to

$$
y=-\frac{x-1}{2} \pm \frac{3 x-1}{2} \sqrt{\frac{x-1}{9 x-1}}
$$

When $|x|$ is large, then the quantity in the radical is approximately $1 / 9$, so the solution becomes

$$
y=-\frac{x-1}{2} \pm \frac{3 x-1}{6}
$$

that is,

$$
y=\frac{1}{3} \quad \text { or } \quad y \approx-x+\frac{2}{3}
$$

These are the equations of the asymptotes! What is wrong? Why does the application of the quadratic-formula show incorrect asymptotes?

When a problem is most paradoxical,
Then don't let it become cardiotaxical. You just say it aloud
To the P M E crowd.
And the answer you'll get from some foxy gal.
Also solved by WIILIAM CHAU, New York, $N Y$, MARK EVANS, Louisville, $K Y$, RICHARD I. HESS, Rancho Palos Verdes, CA, and the PROPOSER Two other solvers sent in a pair of faulty solutions and a partial solution.
752. [Spring 1991] Proposed by the late Charles W. Trigg San Diego, California.

Martin Gardner ("MathematicalGames," Scientific American, April 1964, page 135) has shown that the minimum sum of three 3-digit primes that contain the nine non-zero digits is 999 . Bind a set of three such primes that sums to another multiple of 37 .

Solution by William Chau, New York, New York.
Let $\boldsymbol{S}$ be the sum of the three primes. The units digit of each prime is one of $\mathbf{1 , 3}, \mathbf{7}$, and 9 , so the units digit of $\boldsymbol{S}$ will be $\mathbf{1 , 3 , 7}$, or 9 according as $\mathbf{9 , 7 , 3}$, or $\mathbf{1}$ is not used as the units digit of one of the primes. The maximum value of $S$ is $100(6+8+9)+10(2+4+5)+(1+3+7)=2421<$ 6637 and the minimum value is 999 .

Each prime is congruent to the sum of its digits modulo 9. Therefore $S \equiv 1+2+3+\ldots+$ $9=45 \equiv 0(\bmod 9)$, so $\boldsymbol{S}$ is a multiple of $9-37=333$. The only multiple of 333 between 999 and 2421 that terminates in $\mathbf{1}, \mathbf{3}, \mathbf{7}$, or 9 is 2331 , so $\boldsymbol{S}=2331$ and the units digits of the primes are $\mathbf{1}, 3$, and 7 . For $S$ to be that large, the hundreds digits of the primes must be $\mathbf{9}, \mathbf{8}$, and $\mathbf{4}$ or 5 or 6 . Since their tens digits are then 2 , and two of 4 and 5 and 6 , the tens column produces a carry of exactly $\mathbf{1}$ to the hundreds column, so the hundreds digits of the primes are $\mathbf{9 , 8}$, and 5 . The tens digits are thus 2, 4, and 6. From a table of primes we see that there are just 12 primes that meet the requirements for the addends: $521,523,541,547,563,821,823,827,863,941,947$, and 967 . It is easy now to find that there are exactly the four solutions $\{521,863,947\},\{541,823,967\},\{563,821,947\}$, and $\{563,827,941\}$.

Also solved by CHARLES ASHBACHER, Hiawatha, LA, MARK EVANS, Louisville, KY, VICTOR G. FESER, University of Mary, Bismarck, ND, DAWN M. GALAYDA, St. Bonaventure University, NY, STEPHEN I. GENDLER, Clarion University of Pennsylvania, RICHARD I. HESS, Rancho Palm Verdes, CA, HENRY S. LIEBERMAN, Waban, MA, BOB PRIELIPP, University of Wisconsin-Oshkosh, KENNETH M. WILKE, Topeka, KS, REX H. WU, New York, NY, and the PROPOSER. Ashbacher, Evans, Lieberman, and Wu each found the solution to the originalproblem 149 $+263+587=999$.
753. [Spring 1991] Proposed by R. S. Luthar, University of Wisconsin Center at Janesville, Janesville, Wisconsin.

Solve simultaneously
$e^{4 x}+e^{4 y}=82$ and $e^{x} \cdot e^{y}=2$

Solution by George P. Evanovich, Saint Peter's College, Jersey City, New Jersey
Let $\boldsymbol{e}^{\boldsymbol{x}}=u+v$ and $\boldsymbol{e}^{\boldsymbol{y}}=u=v$. From the second given equation, $v=1$, so we have

$$
(u+1)^{4}+(u-1)^{4}=82, \text { or } u^{4}+6 u^{2}-40=0
$$

and hence $\boldsymbol{u}= \pm \mathbf{2}$ or $\pm i \sqrt{ } 10$. Then $\left(e^{x}, e^{y}\right)=(3,1),(-1,-3),(1+i \sqrt{ } 10,-1+i \sqrt{ } 10)$, or $(1-i \sqrt{ } 10-1$. $i \sqrt{10})$. Now $(x, y)=(\ln 3,0)$ (the only real solution), $(\pi i+2 \pi k$, In $3+2 \pi k i)$, or $\left(\ln \sqrt{11}+i\left(\tan ^{-1}( \pm \sqrt{10})\right.\right.$ $+2 \pi k)$, In $\left.\sqrt{11}+i\left(\tan ^{-1}(\mp \sqrt{ } 10)+2 \pi k\right)\right)$, where $k$ is an integer.

Also solved by JOHN T. ANNULIS, University of Arkansas-Monticello, CHARLES ASHBACHER, Hiawatha,LA, SEUNG-JIN BANG, Seoul, Korea, FRANKP. BATTLES,Massachusetts Maritime Academy, Buzzards Bay, DIETER BENNEWITZ, Koblenz, Germany, SCOTT H. BROWN, Stuart Middle School, FL, BARRY BRUNSON, Western Kentucky University, Bowling Green, JAMES E. CAMPBELL, Indiana University at Bloomington, CAVELAND MATH GROUP, Western Kentucky University, Bowling Grew, JEAN CHAPMAN, Creston, LA, WILLIAM CHAU, New York. NY, PATRICK COSTELLO, Eastern Kentucky University, Richmond, CHARLES R. DIMINNIE, St. Bonaventure University, $\boldsymbol{N} Y$, ROBERT O. DOWNES, Long Beach, CA, RUSSELLEULER, Northwest Missouri State University, Maryville, MARK EVANS, Louisville. KY, VICTOR G. FESER, Universityof Mary, Bismarck, ND, HOWARD FORMAN, Parsippany, NJ, DAWN M. GALAYDA, St. Bonaventure University, $N Y$, ROBERTC. GEBHARDT, Hopatcong, $N$, STEPHEN I. GENDLER, Clarion University of Pennsylvania, RICHARD A. GIBBS, Fort Lewis College, Durango, CO, STAN HARTZLER, Messiah College, Grantham, PA, RICHARD I. HESS, Rancho Palos Verdes, CA, NATHAN JASPEN, Stevens Institute of Technology, Hoboken, NJ, HENRY S. UEBERMAN, Waban, AM, PETER A. LINDSTROM, North Lake College, Irving, TX, LOWELL F. LYNDE, JR., University of Arkansas at Monticello, G. MAVRIGIAN, YoungstownState University, OH, YOSHINOBU MURAYOSHI, Eugene, OR, WILLIAM H. PEIRCE, Stonington, CT, BOB PRIELIPP, University of Wisconsin-Oshkosh, GEORGE W. RAINEY, Los Angeles, CA, PARUSH SAXENA, Massachusetts Maritime Academy, Buzzards Bay, MOHAMMAD P. SHAIKH, Western Michigan University, Kalamazoo, WADE $H$. SHERARD, Furman University, Greenville, SC, KENNETHM. WILKE, Topeka, KS, REX H. WU, New York, NY, DAVID YAVENDITTI, Alma, MI. and the PROPOSER.
754. [Spring 1991] Proposed by Seung-Jin Bang, Seoul, Korea.

Let $a,=a,=1, a_{3}=2$, and $\boldsymbol{a}_{\mathrm{n}+1}=a, \cdot \boldsymbol{a}_{\mathrm{n}-1}+\boldsymbol{a}_{\mathrm{n}-2}$ for $n>3$. Show that

$$
a_{n+2} a_{n} a_{n-2}-a_{n+2} a_{n-1}^{2}-a_{n+1}^{2} a_{n-2}+2 a_{n+1} a_{n} a_{n-1}-a_{n}^{3}+3=0
$$

I. Solution by the Proposer.

It suffices to show that

$$
\left|\begin{array}{ccc}
a_{n+2} & a_{n+1} & a_{n} \\
a_{n+1} & a_{n} & a_{n-1} \\
a_{n} & a_{n-1} & a_{n-2}
\end{array}\right|=-3 .
$$

To that end, we have

$$
\left[\begin{array}{ccc}
a_{n+2} & a_{n+1} & a_{n} \\
a_{n+1} & a_{n} & a_{n-1} \\
a_{n} & a_{n-1} & a_{n-2}
\end{array}\right]=\left[\begin{array}{ccc}
1 & -1 & 1 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right]\left[\begin{array}{ccc}
a_{n+1} & a_{n} & a_{n-1} \\
a_{n} & a_{n-1} & a_{n-2} \\
a_{n-1} & a_{n-2} \cdot & a_{n-3}
\end{array}\right],
$$

and hence

$$
\left[\begin{array}{ccc}
a_{n+2} & a_{n+1} & a_{n} \\
a_{n+1} & a_{n} & a_{n-1} \\
a_{n} & a_{n-1} & a_{n-2}
\end{array}\right]=\left[\begin{array}{ccc}
1 & -1 & 1 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right]^{n-3}\left[\begin{array}{lll}
1 & 2 & 2 \\
2 & 2 & 1 \\
2 & 1 & 1
\end{array}\right] .
$$

Since the two matrices on the right side of this last equation have determinants 1 and -3 respectively, the determinant of the matrix on the left side is -3 .
II. Solution by Murray S. Klamkin, University of Alberta, Edmonton, Alberta, Canada.

As stated, $\boldsymbol{a}_{\mathbf{4}}$ is not determined, so that the given equation is not necessarily true for $n=3$. If the recursion formula is valid for $n \geq 3$, then $a,=\boldsymbol{a}_{\mathbf{3}} \cdot a,+a,=2$ and $\boldsymbol{a}_{\mathbf{5}}=1$. Since we now have $\boldsymbol{a}_{\mathrm{a}+2}$ $=\boldsymbol{a}_{\mathrm{s} 2}$ for $n \geq 3$, it suffices to check the desired equation for $n=3,4,5$, and 6 , for which values it is true.
III. Solution by Rex H. Wu, New York, New York

To express $\boldsymbol{a}_{\mathrm{a}}$ in terms of $\boldsymbol{n}$, let $\boldsymbol{a}_{\mathrm{n}}=\boldsymbol{\lambda}$. Then

$$
\begin{aligned}
& \lambda^{n+1}=\lambda^{n}-\lambda^{n-1}+\lambda^{n-2}, \text { so that } \lambda^{3}=A^{2}-A+1, \\
& (A \cdot 1)\left(\lambda^{2}+1\right)=0, \text { and finally } A=1, \pm i .
\end{aligned}
$$

Any linear combination of solutions is also a solution, so we have that

$$
a_{\mathrm{n}}=\alpha+\beta i^{\mathrm{n}}+\gamma(-i)^{\mathrm{n}}
$$

for some complex constants $\boldsymbol{a}, \beta$, and $\boldsymbol{\gamma}$. By hypothesis we must have

$$
\begin{gathered}
a_{1}=\alpha+\beta i-\gamma i=1, a_{2}=\alpha-\beta \cdot \gamma=1, \\
\text { and } a_{3}=\alpha-\beta i+\gamma i=2,
\end{gathered}
$$

which we solve simultaneously to get that

$$
\alpha=\frac{3}{2}, \quad \beta=\frac{i=1}{4 i}, \quad \text { and } \gamma=\frac{i+1}{4 i} .
$$

Therefore,

$$
a_{n}=\frac{3}{2}+\frac{i-1}{4 i} i^{n}+\frac{i+1}{4 i}(-i)^{n} .
$$

It follows immediately that $\boldsymbol{a}_{\boldsymbol{n}+4}=a$., By tedious but straightforward algebra one can show that $a_{\mathrm{a}} \cdot \boldsymbol{a}_{\mathrm{n}+2}=2$ and that $\boldsymbol{a}_{\mathrm{n}}^{2}+\boldsymbol{a}_{\mathrm{n}+2}^{2}=5$. Now we have

$$
\begin{aligned}
& a_{n+2} a_{n} a_{n-2}-a_{n+2} a_{n-1}^{2}-a_{n+1}^{2} a_{n-2}+2 a_{n+1} a_{n} a_{n-1}-a_{n}^{3}+3 \\
&=2 a_{n-2}-a_{n-2} a_{n-1}^{2}-a_{n+1}^{2} a_{n-2}+4 a_{n}-a_{n}^{3}+3 \\
&=2 a_{n-2}-5 a_{n-2}+4 a_{n}-a_{n}^{3}+3 \\
&=-3 a_{n-2}+4 a_{n}-a_{n}^{3}+3 .
\end{aligned}
$$

Since $\boldsymbol{a}_{\mathbf{a}}=\mathbf{1}$ ifn $\equiv \mathbf{1}$ or $2(\bmod 4)$ and $\boldsymbol{a}_{\mathrm{n}}=2$ ifn $\equiv 0$ or $3(\bmod 4)$, we need verify only that the last displayed line is zero for $\boldsymbol{a}_{\mathrm{r} 2}=1$ and $a,=2$, and for $\boldsymbol{a}_{\mathrm{m} 2}=2$ and $\mathrm{a},=1$, which is easily accomplished.

Editor's comment. Only Klamkin spotted the omission, which was my error. The proposer had stated anly the defining equations; I added the inequality. So I shall do my penance at least $n>3$ times.

Also solved by CHARLES ASHBACHER, Hiawatha, LA, SCOTT H. BROWN, Stuart Middle School, FL, JAMES E. CAMPBELL, Indiana Universityat Bloomington, CAVELAND MATH GROUP, Western Kentucky University, Bowling Green, WILLIAM CHAU, New York, NY, RUSSELL EULER, Northwest Missouri State Universit, Maryville, MARK EVANS, Louisville, KY, VICTOR G. FESER, Universityof Mary, Bismarck, ND, HOWARD FORMAN, Parsippany, NJ, ROBERT C. GEBHARDT, Hopatcong NJ, RICHARD I. HESS, Rancho Palos Verdes, CA, HENRYS. LIEBERMAN, Waban, MA, WILLIAM H. PEIRCE, Stonington, CT, MOHAMMAD P. SHAIKH, Western Michigan University, Kalamazoo, KENNETH M. WILKE, Topeka, KS, and DAVID YAVENDITTI, Alma, MI.
755. [Spring 1991] Proposed by Stanley Rabinowit, Alliant Computer Systems Corp, Littleton, Massachusetts.

In triangle $A B C$, a circle of radiusp is inscribed in the wedge bounded by sides $A B$ and $B C$ and the incircle $(I)$ of the triangle. A circle of radius $q$ is inscribed in the wedge bounded by sides $A C$ and $B C$ and the incircle. If $\mathrm{p}=\mathrm{q}$, prove that $\boldsymbol{A B}=A C$.


## I. Solution by Richard I. Hess, Rancho Palos Verdes, California.

Let the incircle (7) touch $B C$ at X . If the two side circles have the same radius, then a reflection about the line $X X$ leaves the picture unchanged, whence $\boldsymbol{A B}=A C$.
II. Solution by Murray S. Klamkin, University of Alberta, Edmonton, Alberta, Canada.

Draw the angle bisector $B I$ of angle $B$ and let $r$ denote the inradius. Let the parallel to $B C$ through the center $P$ of the circle of radiusp cut $I X$ at $Y$. Then $I Y=r=\mathrm{p}$ and $\mathbb{X}=r+\mathrm{p}$, and it now follows easily that

$$
\sin \frac{B}{2}=\frac{r-p}{r+p} \text { and similarly } \quad \sin \frac{C}{2}=\frac{r-q}{r+q}
$$

Finally, $\boldsymbol{p}=\mathrm{q}$ implies that $\sin (B / 2)=\sin (\boldsymbol{C} / 2)$ and hence that $\boldsymbol{A} \boldsymbol{B}=A C$.
Also solved by SEUNG-JIN BANG, Seoul, Korea, DIETER BENNEWITZ, Koblenz, Germany, SCOTT H. BROWN, Stuart Middle School, FL, WILLIAM CHAU, New York, NY, STEPHEN I. GENDLER, Clarion University of Pennsylvania, HENRY S. LIEBERMAN, Waban, MA, T. R. K PAPPU, OccidentalCollege, Los Angeles, CA, MOHAMMAD P. SHAIKH, Western Michigan University, Kalamazoo, WADE H. SHERARD, Funnan University, Greenville, SC, REX H. WU, New York, NY, DAVID YAVENDITTI, Alma, MI, and the PROPOSER.
756. [Spring 1991] Proposed by Basil Rennie, Bumside, South Australia.

Consider covering the unit interval $[0,1]$ with $\boldsymbol{n}$ measurable subsets, under the constraint that all $\boldsymbol{n}$ subsets must have the same centroid. The centroid $m$ of a set E may be defined by $\int_{\mathbf{E}}(x-m) d r$ $=0$. How can you choose then sets to minimize $m$ ?

For example, if $n=4$, it is possible to make $m=7 / 20$ by choosing the four sets $[0,2 / 5] \cup[9 / 10,1],[0,1 / 5] \cup[4 / 5,9 / 10],[1 / 20,1 / 4] \cup[7 / 10,4 / 5]$, and $[0,7 / 10]$.

Solution by the Proposer.
The smallest value of $m$ is $\mathbf{1}\left(\mathbf{1}+\boldsymbol{v} \boldsymbol{n}\right.$ ), which we denote by $c$. For, let $\boldsymbol{E}_{\mathrm{s}}$ (for $r=1,2, \ldots, \boldsymbol{n}$ ) consist of the union of the two intervals $[0, c]$ and $[c+c \sqrt{ }(r \cdot 1), c+c \sqrt{ } r]$. Each $E_{\mathbf{r}}$ has centroid $c$ and together they cover the interval.

To show this value $c$ is best possible, take $\boldsymbol{n}$ sets $\boldsymbol{E}_{\mathbf{r}}$ covering the interval and with centroids at $m$. Divide each set into $E_{s}^{\prime}$ to the left of $m$ and $\boldsymbol{E}_{r}^{w}$ to the right of $m$. The first moments of the two subsets about $m$ must add to zero, and therefore the moment of $\boldsymbol{E}_{\mathrm{r}}^{\prime \prime}$ about $m$ can be no more than $\boldsymbol{m}^{\mathbf{2}} / \mathbf{2}$, but the sum of all these moments over then sets is at least $(1-m)^{2} / 2$. Hence, $n m^{2} a(1-m)^{2}$, or $[m(v n+1)-1][m(v n-1)+1] a 0$, which is true whenever the quantity in the first brackets is nonnegative; that is, when $m \geq c$.
757. [Spring 1991] Proposed by Paul Anthony Courtney, graduate student, San Diego State University, San Diego, California.

Find the overall height of the pyramid formed from four spherical balls of radius $r$. Student solutions are especially solicited.

Solution by David Yavenditti, high school student, Alma, Michigan.
Consider instead the pyramid formed by $\boldsymbol{n}$ triangular stacks of spheres, each of radius $r$. Let $h$ be the overall height and let $A, B, \boldsymbol{C}, D$ be the centers of the four comer spheres, which determine a regular tetrahedron of edge $2 r(n-1)$, as shown in the accompanying figure. We must find the length of the altitude $B O$ of the tetrahedron, O being the center of the equilateral triangle $A C D$. Then $O M$ is perpendicular to $A C$ and triangle $A O M$ is a $30^{\circ}-60^{\circ}$ right triangle, so $A O=\mathbf{2 r}(n-1) / \sqrt{3}$. Since $A B O$ is a right triangle, then $B O=2 r(n-1) \sqrt{ } 6 / 3$, and the overall height is given by

$$
h=B O+2 r=2 r\left(1+\frac{(n-1) \sqrt{6}}{3}\right)
$$

The case we seek is $\boldsymbol{n}=2$, soh $=\mathbf{2 r}(1+\sqrt{6} / \mathbf{3})$.


Also solved by CHARLES ASHBACHER, Hiawatha, IA, MARTIN BAZANT, Tucson, AZ, WILLIAM CHAU, New York, NY, ROB DOWNES, Long Beach, CA, RUSSELL EULER, Northwest Missouri Stale University, Maryville, MARK EVANS, Louisville, KY, RICHARD I. HESS, Rancho Palos Verdes, CA, HENRY S. LIEBERMAN,Waban, MA, LOWELL F. LYNDE, JR., University of Arkansas at Monticello, MOHAMMADP. SHAIKH, Western Michigan University, Kalamazoo, REX H. WU, New York, $N Y$, and the PROPOSER.

## MESSAGE FROM THE SECRETARY-TREASURER

Copies of the new, revised Constitution and Bylaws are now available. The prices are: $\$ 1.50$ for each of the first four copies and $\$ 1$ for each copy thereafter. I.e., $\$(1.50 n)$ for $n<4$ and $\$(n+2)$ for $n \geq 4$.

The videotape of Professor Joseph A. Gallian's AMS-MAA-PME Invited Address, "The Mathematics of Identification Numbers," given as part of PME's 75th Anniversary Celebration at Boulder, CO, in August, 1989, is also still available. The tape may be borrowed free of charge by PME chapters, and by others upon an advance payment of $\$ 10$. Please contact $m y$ office if you desire to borrow the tape, telling me the date on which you would like to use it. I prefer to mail the tape directly to faculty advisors, and expect them to take responsibility for returning it to my office. Please submit your request in writing and include a phone number and a time that I might reach you if there are problems. Robert M. Woodside, Secretary-Treasurer, Department of Mathematics, East Carolina University, Greenville, NC 27858.

## UPCOMING PI MU EPSILON 1992 NATIONAL MEETING

The-national meeting for Pi Mu Epsilon this year will be very special. Usually, the national meetingeis held in conjunction with the national meetings of the American Mathematical Society and the Mathematical Association of America. In 1992, however, the International Congress of Mathematics Educators (ICME) will hold its annual meeting in Quebec City, in Canada. It has been the policy of the AMS and MAA that in order to avoid a conflict in scheduling, summer meetings are not held in years when an international mathematics meeting (e.g., ICME or the International Congress of Mathematicians) takes place in North America. For this reason, there will be no AMS-MAA national meeting this summer.

Because of these special circumstances, Pi Mu Epsilon will hold its summer meeting in conjunction with the meeting of the MAA Student Chapters. The meeting will take place August 5-8, at Miami University in Oxford, Ohio.

The meeting will begin on the evening of Wednesday, August 5, with a Student Pizza Party and Reception. (Registration and room check-in will begin in the afternoon and continue throughout the evening.)

Highlights of Thursday's program (August 6) will include the MAA Invited Lecture, by Peter Hilton; a reception for Professor Hilton; sessions for contributed papers by PME and MAA student chapter members; presentations by the MAA Modeling Contest winners; meetings of the PME Council and the MAA Student Chapter Committee; and an excursion to the nearby King's Island Theme Park.

The program on Friday, August 7, will feature more student presentations; a choice of two minicourses (open to students and faculty); a panel discussion and display on careers; the Pi Mu Epsilon Banquet; and, finally, the J. Sutherland Frame lecture. This year's Frame lecturer will be Underwood Dudley. The Pi Mu Epsilon portion of the meetings will conclude with informal gatherings after the Frame Lecture.

The meeting will conclude on Saturday, August 8, with the final session of MAA student papers and a choice of two minicourses.

## TRAVEL SUPPORT FOR THE SUMMER MEETING

Pi Mu Epsilon will provide travel support for one student speaker from each chapter. If a chapter is not represented by a student speaker, Pi Mu Epsilon will provide one-half support for a student delegate. Full support is defined to be full round-trip air fare (including ground transportation) from the student's school or home to Oxford, Ohio, up to a maximum of $\$ 600$. (Delegates will receive up to $\$ 300$.) A student who chooses to drive will receive 25 cents per mile for the round trip from school or home to Oxford, up to $\$ 600$. (Delegates will receive $12 \frac{1}{2}$ cents per mile, up to $\$ 300$. Travel support will be provided for only one student per chapter. However, if several students from the same chapter wish to attend, they may share the travel support, if they choose to do so. (Special discounted group airline tickets are available on Delta Airlines through Travel Unlimited, the official travel agency for the conference: 1-800-466-7555.)

For further information about the meeting and the travel support:
SEE YOUR PI MU EPSILON ADVISOR

## GLEANINGSFROM THE CHAPTER REPORTS

CONNECTICUT GAMMA (Fairfield University) During the fall semester the chapter sponsored the second annual Math Bowl Contest. Eight teams of four students competed in a "GE College Bowl ${ }^{\mathrm{n}}$ type of competition, in which all the questions were mathematical. In the spring, members of Pi Mu Epsilon assisted the Mathematics Department in coordinating the activities for Math Counts,-which is a mathematics contest for junior high school students. At the annual spring initiation ceremony thirty-two new members were inducted. "Biostatistics: Who, What, Why, and When?" by Kerrie Eileen Boyle of the Research Triangle Institute was the title of the Pi Mu Epsilon Lecture during the induction ceremony. Dr. Boyle, a 1974 graduate of Fairfield, was also inducted. During the Annual Arts and Sciences Awards Ceremony, two members, Thomas Lipka and Francis Maurais received recognition for their outstanding performance in mathematics. Each was given a Pi Mu Epsilon certificate of achievement, a book each selected in an area of mathematics, and a one-year membership in the Mathematical Association of America.

INDIANA GAMMA (Rose-Hulman Institute of Technology) In the fall of 1990, six students attended the Miami University Conference, with John O'Bryan, Jeff Dierckman, and Omar Zaidi presenting papers

The chapter helped administer the RHIT-St. Mary of the Woods Mathematics Competition (for area high school students) and the 2nd Annual Alfred R. Schmidt Freshmen Mathematics Competition at Rose-Hulman. Mark Roseberry took first place and Jonathan Atkins took second. Our chapter helped the Rose-Hulman Mathematics Department stage the Annual Rose-Hulman Undergraduate Mathematics Conference, which involved over 80 participants and 25 papers. Seven of our students gave papers: John O'Bryan, Jeff Dierckman, Omar Zaidi, Mark Roseberry, Jonathan Atkins, Ben Nicholson, and Tony Hinrichs. Five teams of three members each participated in the Indiana College Mathematics Competition, with the RHIT team of John O'Bryan, Mark Roseberry, and Jonathan Atkins taking first place.

On April 24, 32 new members were initiated into the Indiana Gamma Chapter. It was the 25th anniversary of the founding of the Chapter. The speaker at our initiation banquet was Dr. David Ballew, President of Pi Mu Epsilon and Chairman of the Computer Science Department at Western Illinois University.

KANSAS GAMMA (The Wichita State University) The chapter had a number of speakers during the year. The speakers were: Joseph Stafford, "Tilings," Abdelmalek Kemmou, 'Fuzzy Set Theory \& the Logic of the Continuum;" Ming Liu "Design of Experiment;" Rajiv Bagai, "Formal Logic as a Programming Language;" Dewi Saleh, "Some Mathematical Puzzles." Members of the chapter held free help sessions for undergraduate courses. One of the members, Abdelmalek Kemmou, gave a talk at the joint meetings of the Kansas Section of the MAA and the Kansas Association of Teachers in Mathematics, held at Southwestern College in April, 1991. The chapter also started a publication, called ALEPH TWO. The publication is intended to contain mathematical investigations, mostly by students.

MINNESOTA ZETA (St. Mary's College) The Chapter conducted a number of mathematics colloquia and several chapter-wide business meetings. The Chapter celebrated Math Awareness Week in April with two main activities: Professor Ken Kasin, of St. Mary's presented a talk entitled "Elementary Concepts of Mathematical Chaos"; and eleven new members were initiated into the chapter.

OHIO ZETA (University of Dayton) The chapter continued to be active this year. Among other activities, the members presented several talks at various meetings and conferences. Five students presented talks at the Pi Mu Epsilon Meeting in Columbus, Ohio, in August. Four of them presented the results of the research they conducted in the program "Research Experiences for Undergraduates in Algebraic Graph Theory at the University of Dayton." This program was sponsored by the NSF and Professors Higgins and Mushenheim conducted the program during the summer of 1990. All of these five students also gave talks at the Pi Mu Epsilon Regional Conference held at Miami University, Oxford, Ohio, in September, 1990. These students are Marjorie August, David Gebhard, Tom Bohman, Chicako Mese, and Colleen Hoover. David Gebhard and David Kaas presented talks at the Spring Meeting of the Ohio Section of the MAA held at Bowling Green, Ohio, in April, 1991.

Chicako Mese, Tom Bohman, and Colleen Hoover jointly received UD's Faculty Award for Excellence in Mathematics, while Kristen Toft and Kristine Fromm shared this year's Sophomore Class Award.

VIRGINIA ALPHA (University of Richmond) In the fall, in addition to an initiation ceremony, the Chapter co-sponsored a Math/Computer Science Department colloquium on October 22. The speaker, Professor Jim Kuzmanovich, from Wake Forest University, spoke on "The Lore of Infinity." In the spring, the Chapter held a research forum where four student members, who were engaged in independent research projects, gave 15-minute talks on their projects and how they got started. The speakers were Fran Centofante, Jeff Michel, John Murphy, and David Flader. David Flader presented his paper on Game Theory and Pseudo-Boolean Functions at the National Pi Mu Epsilon Meeting in Orono, Maine, in August. The final event of the year was the annual Pi Mu Epsilon picnic (co-sponsored with the Computer Science Club). At this picnic, Jeff Michel was presented with the award for Outstanding Computer Science Student and David Flader was presented with the award for Outstanding Mathematics Student. Freshman. Kelly Donnellon, was presented with the Pi Mu Epsilon Book Award for outstanding work in Calculus I and II.

WISCONSIN DELTA (St. Norbert College) In August, 1990, three students attended the Pi Mu Epsilon National Conference in Columbus, Ohio. Amy Krebsbach, Mike Lang, and Dave Olson were in attendance, with Mike Lang presenting a paper. In April, 1991, Amy Krebsbach presented a paper at the St. John's University Regional Math Conference. Also in attendance were Dawn Boyung, Chris Cypcar, Amy Gerrits, Mike Lang, and Mike Zittlow.

SNC was host to several speakers during the year. Dr. Bill Shay (UW-Green Bay) spoke on "Cyclic Redundancy Check - Error Detection Using Polynomial Division." Other speakers were: Dr. Alan Parks (Lawrence University) on "Genetic Assembly Line Balancing" and Richard Witalka and John Towne (Schneider National Corporation) on 'The C Programming Language."

Perhaps the biggest event of the year for the chapter was hosting the Fifth Annual Pi Mu Epsilon Regional Conference in November. The featured speaker was Dr. Jeanne LaDuke, of DePaul University, who spoke about the role of women in American mathematics. Laura Donzelli of SNC presented one of the 14 student papers at the conference.

Other significant events included the Ninth Annual SNC High School Math Meet held in conjunction with SNC's math club, Sigma Nu Delta. The annual Brenda Roebke VolleyballTournament was also held in cooperation with Sigma Nu Delta. Part of the proceeds from this tournament go toward a scholarship fund for SNC students majoring in mathematics. This year's winner was Linda Mueller.

## IN MEMORIAM

John T. O'Bryan, the president of the Indiana Gamma Chapter of Pi Mu Epsilon, at RoseHulman Institute of Technology, died December 16, 1991, as a result of injuries he received in a car accident.

John was one of the most active and productive members the Indiana Gamma Chapter has ever had. Between April,1990, and September, 1991, he presented five different papers at six different conferences and meetings. John's most outstanding work resulted from his participation in an NSFfunded REU project at Rose-Hulman, which he attended between his sophomore and junior years. His paper "Maximal Order Three-Rewriteable Subgroups of Symmetric Groups" became the initial Rose-Hulman Institute of Technology Technical Report. John also presented this paper at a special session during the 1991 Winter Meeting of the MAA held in San Francisco.

John's REU experience also led to a paper titled "Large 'Almost Abelian' Subgroups of the Symmetric Group," which he presented at the 17th Annual Regional Pi Mu Epsilon Meeting held at Miami University, Oxford, Ohio. This paper became part of a joint paper written with Dr. Gary Sherman, of Rose-Hulman, titled "Undergraduates, CALEY, and Mathematics," which has been submitted to the Journal of Technology in Mathematics.

But John's work wasn't limited to pure mathematics. During the summer between his junior and senior year, John worked as a summer researcher, in applied mathematics, as a member of the Outstanding Student Summer Program sponsored by Sandia National Laboratories in New Mexico. This experience led to his paper "Parallelization of a Parameter Identification Problem," which he gave at the 18th Annual Pi Mu Epsilon Conference at Miami University, in September, 1991. This was to be John's final Pi Mu Epsilon paper.

As a mathematics/physics double major at Rose-Hulman, John participated in many mathematics competitions as a leading member of the Rose-Hulman team. As a scholar, he received numerous awards, including the top freshman mathematics award, the top freshman student award, the top sophomore student award, and the top junior student award. This spring he will be awarded, posthumously, the Clarence P. Sousley Award as "a graduating mathematics major who has demonstrated exceptional performance in his field.'

John truly lived by the Pi Mu Epsilon pledge "... I will exert my best efforts to promote true scholarship, particularly in mathematics; and that I will support the objectives of the Pi Mu Epsilon Society."
(Elton Graves, RHIT Mathematics Department)

NINETEENTH ANNUAL
PI MU EPSILON
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We invite you to join us. There will be sessions of the student conference on Friday evening and Saturday afternoon. Free overnight lodging for all students will be arranged with Miami students. Each student should bring a sleeping bag. All student guests are invited to a free Friday evening pizza party
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We need your title, presentation time ( 15 or 30 min .), preferred date (Fri. or Sat.) and a 50 (approx.) word abstract by September 25, 1992. Please send to

> Professor Milton D. Cox

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Contact us for more details.

## П M E

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For information, contact:
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* See page 414 for farther details.
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